## ELEMENTARY DIFFERENTIAL EQUATIONS



William F. Trench
Andrew G. Cowles Distinguished Professor Emeritus Department of Mathematics

Trinity University
San Antonio, Texas, USA
wtrench@trinity.edu
This book has been judged to meet the evaluation criteria set by the Editorial Board of the American Institute of Mathematics in connection with the Institute's Open Textbook Initiative. It may be copied, modified, redistributed, translated, and built upon subject to the Creative Commons Attribution-NonCommercial-ShareAlike 3.0 Unported License.

Free Edition 1.01 (December 2013)
This book was published previously by Brooks/Cole Thomson Learning, 2001. This free edition is made available in the hope that it will be useful as a textbook or reference. Reproduction is permitted for any valid noncommercial educational, mathematical, or scientific purpose. However, charges for profit beyond reasonable printing costs are prohibited.

NOTE: This version of the textbook has been edited from its original for use by the University of Central Oklahoma. Several sections (and one whole chapter) have been removed from the text. References to removed material still appear in this text, as indicated by "??".

## TO BEVERLY

## Contents

Chapter 1 Introduction
1.1 Applications Leading to Differential Equations 1.2 First Order Equations ..... 5
1.3 Direction Fields for First Order Equations ..... 14
Chapter 2 First Order Equations
2.1 Linear First Order Equations ..... 27
2.2 Separable Equations ..... 39
2.3 Existence and Uniqueness of Solutions of Nonlinear Equations ..... 48
2.5 Exact Equations ..... 55
2.6 Integrating Factors ..... 63
Chapter 3 Numerical Methods
3.1 Euler's Method ..... 74
3.2 The Improved Euler Method and Related Methods ..... 85
Chapter 4 Applications of First Order Equations
4.2 Cooling and Mixing ..... 96
4.3 Elementary Mechanics ..... 105
4.4 Autonomous Second Order Equations ..... 115
Chapter 5 Linear Second Order Equations
5.1 Homogeneous Linear Equations ..... 132
5.2 Constant Coefficient Homogeneous Equations ..... 146
5.3 Nonhomgeneous Linear Equations ..... 155
5.4 The Method of Undetermined Coefficients I ..... 162
5.5 The Method of Undetermined Coefficients II ..... 169
5.7 Variation of Parameters ..... 178
Chapter 6 Applcations of Linear Second Order Equations
6.1 Spring Problems I ..... 188
6.2 Spring Problems II ..... 197
Chapter 7 Series Solutions of Linear Second Order Equations
7.1 Review of Power Series ..... 208
7.2 Series Solutions Near an Ordinary Point I ..... 219
7.3 Series Solutions Near an Ordinary Point II ..... 231
Chapter 8 Laplace Transforms
8.1 Introduction to the Laplace Transform ..... 240
8.2 The Inverse Laplace Transform ..... 250
8.3 Solution of Initial Value Problems ..... 257
8.4 The Unit Step Function ..... 263
8.5 Constant Coefficient Equations with Piecewise Continuous Forcing Functions ..... 272
8.6 Convolution ..... 280
8.7 Constant Cofficient Equations with Impulses ..... 290
8.8 A Brief Table of Laplace Transforms
Chapter 10 Linear Systems of Differential Equations
10.1 Introduction to Systems of Differential Equations ..... 301
10.2 Linear Systems of Differential Equations ..... 308
10.3 Basic Theory of Homogeneous Linear Systems ..... 313
10.4 Constant Coefficient Homogeneous Systems I ..... 320
10.5 Constant Coefficient Homogeneous Systems II ..... 331
10.6 Constant Coefficient Homogeneous Systems II ..... 344
10.7 Variation of Parameters for Nonhomogeneous Linear Systems ..... 354

## Preface

Elementary Differential Equations with Boundary Value Problems is written for students in science, engineering, and mathematics who have completed calculus through partial differentiation. If your syllabus includes Chapter 10 (Linear Systems of Differential Equations), your students should have some preparation in linear algebra.

In writing this book I have been guided by the these principles:

- An elementary text should be written so the student can read it with comprehension without too much pain. I have tried to put myself in the student's place, and have chosen to err on the side of too much detail rather than not enough.
- An elementary text can't be better than its exercises. This text includes 1695 numbered exercises, many with several parts. They range in difficulty from routine to very challenging.
- An elementary text should be written in an informal but mathematically accurate way, illustrated by appropriate graphics. I have tried to formulate mathematical concepts succinctly in language that students can understand. I have minimized the number of explicitly stated theorems and definitions, preferring to deal with concepts in a more conversational way, copiously illustrated by 250 completely worked out examples. Where appropriate, concepts and results are depicted in 144 figures.

Although I believe that the computer is an immensely valuable tool for learning, doing, and writing mathematics, the selection and treatment of topics in this text reflects my pedagogical orientation along traditional lines. However, I have incorporated what I believe to be the best use of modern technology, so you can select the level of technology that you want to include in your course. The text includes 336 exercises - identified by the symbols C and $\mathrm{C} / \mathrm{G}$ - that call for graphics or computation and graphics. There are also 73 laboratory exercises - identified by L - that require extensive use of technology. In addition, several sections include informal advice on the use of technology. If you prefer not to emphasize technology, simply ignore these exercises and the advice.

There are two schools of thought on whether techniques and applications should be treated together or separately. I have chosen to separate them; thus, Chapter 2 deals with techniques for solving first order equations, and Chapter 4 deals with applications. Similarly, Chapter 5 deals with techniques for solving second order equations, and Chapter 6 deals with applications. However, the exercise sets of the sections dealing with techniques include some applied problems.

Traditionally oriented elementary differential equations texts are occasionally criticized as being collections of unrelated methods for solving miscellaneous problems. To some extent this is true; after all, no single method applies to all situations. Nevertheless, I believe that one idea can go a long way toward unifying some of the techniques for solving diverse problems: variation of parameters. I use variation of parameters at the earliest opportunity in Section 2.1, to solve the nonhomogeneous linear equation, given a nontrivial solution of the complementary equation. You may find this annoying, since most of us learned that one should use integrating factors for this task, while perhaps mentioning the variation of parameters option in an exercise. However, there's little difference between the two approaches, since an integrating factor is nothing more than the reciprocal of a nontrivial solution of the complementary equation. The advantage of using variation of parameters here is that it introduces the concept in its simplest form and focuses the student's attention on the idea of seeking a solution $y$ of a differential equation by writing it as $y=u y_{1}$, where $y_{1}$ is a known solution of related equation and $u$ is a function to be determined. I use this idea in nonstandard ways, as follows:

- In Section 2.4 to solve nonlinear first order equations, such as Bernoulli equations and nonlinear homogeneous equations.
- In Chapter 3 for numerical solution of semilinear first order equations.
- In Section 5.2 to avoid the necessity of introducing complex exponentials in solving a second order constant coefficient homogeneous equation with characteristic polynomials that have complex zeros.
- In Sections 5.4, 5.5, and 9.3 for the method of undetermined coefficients. (If the method of annihilators is your preferred approach to this problem, compare the labor involved in solving, for example, $y^{\prime \prime}+y^{\prime}+y=x^{4} e^{x}$ by the method of annihilators and the method used in Section 5.4.)

Introducing variation of parameters as early as possible (Section 2.1) prepares the student for the concept when it appears again in more complex forms in Section 5.6, where reduction of order is used not merely to find a second solution of the complementary equation, but also to find the general solution of the nonhomogeneous equation, and in Sections 5.7, 9.4, and 10.7, that treat the usual variation of parameters problem for second and higher order linear equations and for linear systems.

You may also find the following to be of interest:

- Section 2.6 deals with integrating factors of the form $\mu=p(x) q(y)$, in addition to those of the form $\mu=p(x)$ and $\mu=q(y)$ discussed in most texts.
- Section 4.4 makes phase plane analysis of nonlinear second order autonomous equations accessible to students who have not taken linear algebra, since eigenvalues and eigenvectors do not enter into the treatment. Phase plane analysis of constant coefficient linear systems is included in Sections 10.4-6.
- Section 4.5 presents an extensive discussion of applications of differential equations to curves.
- Section 6.4 studies motion under a central force, which may be useful to students interested in the mathematics of satellite orbits.
- Sections 7.5-7 present the method of Frobenius in more detail than in most texts. The approach is to systematize the computations in a way that avoids the necessity of substituting the unknown Frobenius series into each equation. This leads to efficiency in the computation of the coefficients of the Frobenius solution. It also clarifies the case where the roots of the indicial equation differ by an integer (Section 7.7).
- The free Student Solutions Manual contains solutions of most of the even-numbered exercises.
- The free Instructor's Solutions Manual is available by email to wtrench @trinity.edu, subject to verification of the requestor's faculty status.

The following observations may be helpful as you choose your syllabus:

- Section 2.3 is the only specific prerequisite for Chapter 3. To accomodate institutions that offer a separate course in numerical analysis, Chapter 3 is not a prerequisite for any other section in the text.
- The sections in Chapter 4 are independent of each other, and are not prerequisites for any of the later chapters. This is also true of the sections in Chapter 6, except that Section 6.1 is a prerequisite for Section 6.2.
- Chapters 7, 8, and 9 can be covered in any order after the topics selected from Chapter 5. For example, you can proceed directly from Chapter 5 to Chapter 9.
- The second order Euler equation is discussed in Section 7.4, where it sets the stage for the method of Frobenius. As noted at the beginning of Section 7.4, if you want to include Euler equations in your syllabus while omitting the method of Frobenius, you can skip the introductory paragraphs in Section 7.4 and begin with Definition 7.4.2. You can then cover Section 7.4 immediately after Section 5.2.

William F. Trench

## CHAPTER 1 <br> Introduction

IN THIS CHAPTER we begin our study of differential equations.
SECTION 1.1 presents examples of applications that lead to differential equations.
SECTION 1.2 introduces basic concepts and definitions concerning differential equations.
SECTION 1.3 presents a geometric method for dealing with differential equations that has been known for a very long time, but has become particularly useful and important with the proliferation of readily available differential equations software.

### 1.1 APPLICATIONS LEADING TO DIFFERENTIAL EQUATIONS

In order to apply mathematical methods to a physical or "real life" problem, we must formulate the problem in mathematical terms; that is, we must construct a mathematical model for the problem. Many physical problems concern relationships between changing quantities. Since rates of change are represented mathematically by derivatives, mathematical models often involve equations relating an unknown function and one or more of its derivatives. Such equations are differential equations. They are the subject of this book.

Much of calculus is devoted to learning mathematical techniques that are applied in later courses in mathematics and the sciences; you wouldn't have time to learn much calculus if you insisted on seeing a specific application of every topic covered in the course. Similarly, much of this book is devoted to methods that can be applied in later courses. Only a relatively small part of the book is devoted to the derivation of specific differential equations from mathematical models, or relating the differential equations that we study to specific applications. In this section we mention a few such applications.

The mathematical model for an applied problem is almost always simpler than the actual situation being studied, since simplifying assumptions are usually required to obtain a mathematical problem that can be solved. For example, in modeling the motion of a falling object, we might neglect air resistance and the gravitational pull of celestial bodies other than Earth, or in modeling population growth we might assume that the population grows continuously rather than in discrete steps.

A good mathematical model has two important properties:

- It's sufficiently simple so that the mathematical problem can be solved.
- It represents the actual situation sufficiently well so that the solution to the mathematical problem predicts the outcome of the real problem to within a useful degree of accuracy. If results predicted by the model don't agree with physical observations, the underlying assumptions of the model must be revised until satisfactory agreement is obtained.

We'll now give examples of mathematical models involving differential equations. We'll return to these problems at the appropriate times, as we learn how to solve the various types of differential equations that occur in the models.

All the examples in this section deal with functions of time, which we denote by $t$. If $y$ is a function of $t, y^{\prime}$ denotes the derivative of $y$ with respect to $t$; thus,

$$
y^{\prime}=\frac{d y}{d t}
$$

## Population Growth and Decay

Although the number of members of a population (people in a given country, bacteria in a laboratory culture, wildflowers in a forest, etc.) at any given time $t$ is necessarily an integer, models that use differential equations to describe the growth and decay of populations usually rest on the simplifying assumption that the number of members of the population can be regarded as a differentiable function $P=P(t)$. In most models it is assumed that the differential equation takes the form

$$
\begin{equation*}
P^{\prime}=a(P) P \tag{1.1.1}
\end{equation*}
$$

where $a$ is a continuous function of $P$ that represents the rate of change of population per unit time per individual. In the Malthusian model, it is assumed that $a(P)$ is a constant, so (1.1.1) becomes

$$
\begin{equation*}
P^{\prime}=a P . \tag{1.1.2}
\end{equation*}
$$

(When you see a name in blue italics, just click on it for information about the person.) This model assumes that the numbers of births and deaths per unit time are both proportional to the population. The constants of proportionality are the birth rate (births per unit time per individual) and the death rate (deaths per unit time per individual); $a$ is the birth rate minus the death rate. You learned in calculus that if $c$ is any constant then

$$
\begin{equation*}
P=c e^{a t} \tag{1.1.3}
\end{equation*}
$$

satisfies (1.1.2), so (1.1.2) has infinitely many solutions. To select the solution of the specific problem that we're considering, we must know the population $P_{0}$ at an initial time, say $t=0$. Setting $t=0$ in (1.1.3) yields $c=P(0)=P_{0}$, so the applicable solution is

$$
P(t)=P_{0} e^{a t} .
$$

This implies that

$$
\lim _{t \rightarrow \infty} P(t)=\left\{\begin{array}{cl}
\infty & \text { if } a>0 \\
0 & \text { if } a<0
\end{array}\right.
$$

that is, the population approaches infinity if the birth rate exceeds the death rate, or zero if the death rate exceeds the birth rate.

To see the limitations of the Malthusian model, suppose we're modeling the population of a country, starting from a time $t=0$ when the birth rate exceeds the death rate (so $a>0$ ), and the country's resources in terms of space, food supply, and other necessities of life can support the existing population. Then the prediction $P=P_{0} e^{a t}$ may be reasonably accurate as long as it remains within limits that the country's resources can support. However, the model must inevitably lose validity when the prediction exceeds these limits. (If nothing else, eventually there won't be enough space for the predicted population!)

This flaw in the Malthusian model suggests the need for a model that accounts for limitations of space and resources that tend to oppose the rate of population growth as the population increases. Perhaps the most famous model of this kind is the Verhulst model, where (1.1.2) is replaced by

$$
\begin{equation*}
P^{\prime}=a P(1-\alpha P), \tag{1.1.4}
\end{equation*}
$$

where $\alpha$ is a positive constant. As long as $P$ is small compared to $1 / \alpha$, the ratio $P^{\prime} / P$ is approximately equal to $a$. Therefore the growth is approximately exponential; however, as $P$ increases, the ratio $P^{\prime} / P$ decreases as opposing factors become significant.

Equation (1.1.4) is the logistic equation. You will learn how to solve it in Section 1.2. (See Exercise 2.2.28.) The solution is

$$
P=\frac{P_{0}}{\alpha P_{0}+\left(1-\alpha P_{0}\right) e^{-a t}},
$$

where $P_{0}=P(0)>0$. Therefore $\lim _{t \rightarrow \infty} P(t)=1 / \alpha$, independent of $P_{0}$.
Figure 1.1.1 shows typical graphs of $P$ versus $t$ for various values of $P_{0}$.


Figure 1.1.1 Solutions of the logistic equation

## Newton's Law of Cooling

According to Newton's law of cooling, the temperature of a body changes at a rate proportional to the difference between the temperature of the body and the temperature of the surrounding medium. Thus, if $T_{m}$ is the temperature of the medium and $T=T(t)$ is the temperature of the body at time $t$, then

$$
\begin{equation*}
T^{\prime}=-k\left(T-T_{m}\right), \tag{1.1.5}
\end{equation*}
$$

where $k$ is a positive constant and the minus sign indicates; that the temperature of the body increases with time if it's less than the temperature of the medium, or decreases if it's greater. We'll see in Section 4.2 that if $T_{m}$ is constant then the solution of (1.1.5) is

$$
\begin{equation*}
T=T_{m}+\left(T_{0}-T_{m}\right) e^{-k t} \tag{1.1.6}
\end{equation*}
$$

where $T_{0}$ is the temperature of the body when $t=0$. Therefore $\lim _{t \rightarrow \infty} T(t)=T_{m}$, independent of $T_{0}$. (Common sense suggests this. Why?)

Figure 1.1.2 shows typical graphs of $T$ versus $t$ for various values of $T_{0}$.


Figure 1.1.2 Temperature according to Newton's Law of Cooling
Assuming that the medium remains at constant temperature seems reasonable if we're considering a cup of coffee cooling in a room, but not if we're cooling a huge cauldron of molten metal in the same room. The difference between the two situations is that the heat lost by the coffee isn't likely to raise the temperature of the room appreciably, but the heat lost by the cooling metal is. In this second situation we must use a model that accounts for the heat exchanged between the object and the medium. Let $T=T(t)$ and $T_{m}=T_{m}(t)$ be the temperatures of the object and the medium respectively, and let $T_{0}$ and $T_{m 0}$ be their initial values. Again, we assume that $T$ and $T_{m}$ are related by (1.1.5). We also assume that the change in heat of the object as its temperature changes from $T_{0}$ to $T$ is $a\left(T-T_{0}\right)$ and the change in heat of the medium as its temperature changes from $T_{m 0}$ to $T_{m}$ is $a_{m}\left(T_{m}-T_{m 0}\right)$, where $a$ and $a_{m}$ are positive constants depending upon the masses and thermal properties of the object and medium respectively. If we assume that the total heat of the in the object and the medium remains constant (that is, energy is conserved), then

$$
a\left(T-T_{0}\right)+a_{m}\left(T_{m}-T_{m 0}\right)=0 .
$$

Solving this for $T_{m}$ and substituting the result into (1.1.6) yields the differential equation

$$
T^{\prime}=-k\left(1+\frac{a}{a_{m}}\right) T+k\left(T_{m 0}+\frac{a}{a_{m}} T_{0}\right)
$$

for the temperature of the object. After learning to solve linear first order equations, you'll be able to show (Exercise 4.2.17) that

$$
T=\frac{a T_{0}+a_{m} T_{m 0}}{a+a_{m}}+\frac{a_{m}\left(T_{0}-T_{m 0}\right)}{a+a_{m}} e^{-k\left(1+a / a_{m}\right) t} .
$$

Glucose Absorption by the Body

Glucose is absorbed by the body at a rate proportional to the amount of glucose present in the bloodstream. Let $\lambda$ denote the (positive) constant of proportionality. Suppose there are $G_{0}$ units of glucose in the bloodstream when $t=0$, and let $G=G(t)$ be the number of units in the bloodstream at time $t>0$. Then, since the glucose being absorbed by the body is leaving the bloodstream, $G$ satisfies the equation

$$
\begin{equation*}
G^{\prime}=-\lambda G \tag{1.1.7}
\end{equation*}
$$

From calculus you know that if $c$ is any constant then

$$
\begin{equation*}
G=c e^{-\lambda t} \tag{1.1.8}
\end{equation*}
$$

satisfies (1.1.7), so (1.1.7) has infinitely many solutions. Setting $t=0$ in (1.1.8) and requiring that $G(0)=G_{0}$ yields $c=G_{0}$, so

$$
G(t)=G_{0} e^{-\lambda t}
$$

Now let's complicate matters by injecting glucose intravenously at a constant rate of $r$ units of glucose per unit of time. Then the rate of change of the amount of glucose in the bloodstream per unit time is

$$
\begin{equation*}
G^{\prime}=-\lambda G+r, \tag{1.1.9}
\end{equation*}
$$

where the first term on the right is due to the absorption of the glucose by the body and the second term is due to the injection. After you've studied Section 2.1, you'll be able to show (Exercise 2.1.43) that the solution of (1.1.9) that satisfies $G(0)=G_{0}$ is

$$
G=\frac{r}{\lambda}+\left(G_{0}-\frac{r}{\lambda}\right) e^{-\lambda t} .
$$

Graphs of this function are similar to those in Figure 1.1.2. (Why?)

## Spread of Epidemics

One model for the spread of epidemics assumes that the number of people infected changes at a rate proportional to the product of the number of people already infected and the number of people who are susceptible, but not yet infected. Therefore, if $S$ denotes the total population of susceptible people and $I=I(t)$ denotes the number of infected people at time $t$, then $S-I$ is the number of people who are susceptible, but not yet infected. Thus,

$$
I^{\prime}=r I(S-I),
$$

where $r$ is a positive constant. Assuming that $I(0)=I_{0}$, the solution of this equation is

$$
I=\frac{S I_{0}}{I_{0}+\left(S-I_{0}\right) e^{-r S t}}
$$

(Exercise 2.2.29). Graphs of this function are similar to those in Figure 1.1.1. (Why?) Since $\lim _{t \rightarrow \infty} I(t)=$ $S$, this model predicts that all the susceptible people eventually become infected.

## Newton's Second Law of Motion

According to Newton's second law of motion, the instantaneous acceleration $a$ of an object with constant mass $m$ is related to the force $F$ acting on the object by the equation $F=m a$. For simplicity, let's assume that $m=1$ and the motion of the object is along a vertical line. Let $y$ be the displacement of the object from some reference point on Earth's surface, measured positive upward. In many applications, there are three kinds of forces that may act on the object:
(a) A force such as gravity that depends only on the position $y$, which we write as $-p(y)$, where $p(y)>0$ if $y \geq 0$.
(b) A force such as atmospheric resistance that depends on the position and velocity of the object, which we write as $-q\left(y, y^{\prime}\right) y^{\prime}$, where $q$ is a nonnegative function and we've put $y^{\prime}$ "outside" to indicate that the resistive force is always in the direction opposite to the velocity.
(c) A force $f=f(t)$, exerted from an external source (such as a towline from a helicopter) that depends only on $t$.
In this case, Newton's second law implies that

$$
y^{\prime \prime}=-q\left(y, y^{\prime}\right) y^{\prime}-p(y)+f(t)
$$

which is usually rewritten as

$$
y^{\prime \prime}+q\left(y, y^{\prime}\right) y^{\prime}+p(y)=f(t)
$$

Since the second (and no higher) order derivative of $y$ occurs in this equation, we say that it is a second order differential equation.

Interacting Species: Competition
Let $P=P(t)$ and $Q=Q(t)$ be the populations of two species at time $t$, and assume that each population would grow exponentially if the other didn't exist; that is, in the absence of competition we would have

$$
\begin{equation*}
P^{\prime}=a P \quad \text { and } \quad Q^{\prime}=b Q \tag{1.1.10}
\end{equation*}
$$

where $a$ and $b$ are positive constants. One way to model the effect of competition is to assume that the growth rate per individual of each population is reduced by an amount proportional to the other population, so (1.1.10) is replaced by

$$
\begin{aligned}
P^{\prime} & =a P-\alpha Q \\
Q^{\prime} & =-\beta P+b Q
\end{aligned}
$$

where $\alpha$ and $\beta$ are positive constants. (Since negative population doesn't make sense, this system works only while $P$ and $Q$ are both positive.) Now suppose $P(0)=P_{0}>0$ and $Q(0)=Q_{0}>0$. It can be shown (Exercise 10.4.42) that there's a positive constant $\rho$ such that if $\left(P_{0}, Q_{0}\right)$ is above the line $L$ through the origin with slope $\rho$, then the species with population $P$ becomes extinct in finite time, but if $\left(P_{0}, Q_{0}\right)$ is below $L$, the species with population $Q$ becomes extinct in finite time. Figure 1.1.3 illustrates this. The curves shown there are given parametrically by $P=P(t), Q=Q(t), t>0$. The arrows indicate direction along the curves with increasing $t$.


Figure 1.1.3 Populations of competing species

### 1.2 BASIC CONCEPTS

A differential equation is an equation that contains one or more derivatives of an unknown function. The order of a differential equation is the order of the highest derivative that it contains. A differential equation is an ordinary differential equation if it involves an unknown function of only one variable, or a partial differential equation if it involves partial derivatives of a function of more than one variable. For now we'll consider only ordinary differential equations, and we'll just call them differential equations.

Throughout this text, all variables and constants are real unless it's stated otherwise. We'll usually use $x$ for the independent variable unless the independent variable is time; then we'll use $t$.
The simplest differential equations are first order equations of the form

$$
\frac{d y}{d x}=f(x) \quad \text { or, equivalently, } \quad y^{\prime}=f(x)
$$

where $f$ is a known function of $x$. We already know from calculus how to find functions that satisfy this kind of equation. For example, if

$$
y^{\prime}=x^{3}
$$

then

$$
y=\int x^{3} d x=\frac{x^{4}}{4}+c
$$

where $c$ is an arbitrary constant. If $n>1$ we can find functions $y$ that satisfy equations of the form

$$
\begin{equation*}
y^{(n)}=f(x) \tag{1.2.1}
\end{equation*}
$$

by repeated integration. Again, this is a calculus problem.
Except for illustrative purposes in this section, there's no need to consider differential equations like (1.2.1).We'll usually consider differential equations that can be written as

$$
\begin{equation*}
y^{(n)}=f\left(x, y, y^{\prime}, \ldots, y^{(n-1)}\right) \tag{1.2.2}
\end{equation*}
$$

where at least one of the functions $y, y^{\prime}, \ldots, y^{(n-1)}$ actually appears on the right. Here are some examples:

$$
\begin{array}{rlll}
\frac{d y}{d x}-x^{2} & =0 & & (\text { first order }) \\
\frac{d y}{d x}+2 x y^{2} & =-2 & & (\text { first order }) \\
\frac{d^{2} y}{d x^{2}}+2 \frac{d y}{d x}+y & =2 x & & \text { (second order) } \\
x y^{\prime \prime \prime}+y^{2} & =\sin x & & \text { (third order) } \\
y^{(n)}+x y^{\prime}+3 y & =x & & (n \text {-th order })
\end{array}
$$

Although none of these equations is written as in (1.2.2), all of them can be written in this form:

$$
\begin{aligned}
y^{\prime} & =x^{2} \\
y^{\prime} & =-2-2 x y^{2} \\
y^{\prime \prime} & =2 x-2 y^{\prime}-y \\
y^{\prime \prime \prime} & =\frac{\sin x-y^{2}}{x} \\
y^{(n)} & =x-x y^{\prime}-3 y
\end{aligned}
$$

## Solutions of Differential Equations

A solution of a differential equation is a function that satisfies the differential equation on some open interval; thus, $y$ is a solution of (1.2.2) if $y$ is $n$ times differentiable and

$$
y^{(n)}(x)=f\left(x, y(x), y^{\prime}(x), \ldots, y^{(n-1)}(x)\right)
$$

for all $x$ in some open interval $(a, b)$. In this case, we also say that $y$ is a solution of (1.2.2) on $(a, b)$. Functions that satisfy a differential equation at isolated points are not interesting. For example, $y=x^{2}$ satisfies

$$
x y^{\prime}+x^{2}=3 x
$$

if and only if $x=0$ or $x=1$, but it's not a solution of this differential equation because it does not satisfy the equation on an open interval.

The graph of a solution of a differential equation is a solution curve. More generally, a curve $C$ is said to be an integral curve of a differential equation if every function $y=y(x)$ whose graph is a segment of $C$ is a solution of the differential equation. Thus, any solution curve of a differential equation is an integral curve, but an integral curve need not be a solution curve.

Example 1.2.1 If $a$ is any positive constant, the circle

$$
\begin{equation*}
x^{2}+y^{2}=a^{2} \tag{1.2.3}
\end{equation*}
$$

is an integral curve of

$$
\begin{equation*}
y^{\prime}=-\frac{x}{y} \tag{1.2.4}
\end{equation*}
$$

To see this, note that the only functions whose graphs are segments of (1.2.3) are

$$
y_{1}=\sqrt{a^{2}-x^{2}} \quad \text { and } \quad y_{2}=-\sqrt{a^{2}-x^{2}} .
$$

We leave it to you to verify that these functions both satisfy (1.2.4) on the open interval $(-a, a)$. However, (1.2.3) is not a solution curve of (1.2.4), since it's not the graph of a function.

Example 1.2.2 Verify that

$$
\begin{equation*}
y=\frac{x^{2}}{3}+\frac{1}{x} \tag{1.2.5}
\end{equation*}
$$

is a solution of

$$
\begin{equation*}
x y^{\prime}+y=x^{2} \tag{1.2.6}
\end{equation*}
$$

on $(0, \infty)$ and on $(-\infty, 0)$.

Solution Substituting (1.2.5) and

$$
y^{\prime}=\frac{2 x}{3}-\frac{1}{x^{2}}
$$

into (1.2.6) yields

$$
x y^{\prime}(x)+y(x)=x\left(\frac{2 x}{3}-\frac{1}{x^{2}}\right)+\left(\frac{x^{2}}{3}+\frac{1}{x}\right)=x^{2}
$$

for all $x \neq 0$. Therefore $y$ is a solution of (1.2.6) on $(-\infty, 0)$ and $(0, \infty)$. However, $y$ isn't a solution of the differential equation on any open interval that contains $x=0$, since $y$ is not defined at $x=0$.

Figure 1.2.1 shows the graph of (1.2.5). The part of the graph of (1.2.5) on $(0, \infty)$ is a solution curve of $(1.2 .6)$, as is the part of the graph on $(-\infty, 0)$.

Example 1.2.3 Show that if $c_{1}$ and $c_{2}$ are constants then

$$
\begin{equation*}
y=\left(c_{1}+c_{2} x\right) e^{-x}+2 x-4 \tag{1.2.7}
\end{equation*}
$$

is a solution of

$$
\begin{equation*}
y^{\prime \prime}+2 y^{\prime}+y=2 x \tag{1.2.8}
\end{equation*}
$$

on $(-\infty, \infty)$.

Solution Differentiating (1.2.7) twice yields

$$
y^{\prime}=-\left(c_{1}+c_{2} x\right) e^{-x}+c_{2} e^{-x}+2
$$

and

$$
y^{\prime \prime}=\left(c_{1}+c_{2} x\right) e^{-x}-2 c_{2} e^{-x}
$$

so

$$
\begin{aligned}
y^{\prime \prime}+2 y^{\prime}+y= & \left(c_{1}+c_{2} x\right) e^{-x}-2 c_{2} e^{-x} \\
& +2\left[-\left(c_{1}+c_{2} x\right) e^{-x}+c_{2} e^{-x}+2\right] \\
& +\left(c_{1}+c_{2} x\right) e^{-x}+2 x-4 \\
= & (1-2+1)\left(c_{1}+c_{2} x\right) e^{-x}+(-2+2) c_{2} e^{-x} \\
& +4+2 x-4=2 x
\end{aligned}
$$

for all values of $x$. Therefore $y$ is a solution of (1.2.8) on $(-\infty, \infty)$.


Figure 1.2.1 $y=\frac{x^{2}}{3}+\frac{1}{x}$

Example 1.2.4 Find all solutions of

$$
\begin{equation*}
y^{(n)}=e^{2 x} . \tag{1.2.9}
\end{equation*}
$$

Solution Integrating (1.2.9) yields

$$
y^{(n-1)}=\frac{e^{2 x}}{2}+k_{1}
$$

where $k_{1}$ is a constant. If $n \geq 2$, integrating again yields

$$
y^{(n-2)}=\frac{e^{2 x}}{4}+k_{1} x+k_{2}
$$

If $n \geq 3$, repeatedly integrating yields

$$
\begin{equation*}
y=\frac{e^{2 x}}{2^{n}}+k_{1} \frac{x^{n-1}}{(n-1)!}+k_{2} \frac{x^{n-2}}{(n-2)!}+\cdots+k_{n} \tag{1.2.10}
\end{equation*}
$$

where $k_{1}, k_{2}, \ldots, k_{n}$ are constants. This shows that every solution of (1.2.9) has the form (1.2.10) for some choice of the constants $k_{1}, k_{2}, \ldots, k_{n}$. On the other hand, differentiating (1.2.10) $n$ times shows that if $k_{1}, k_{2}, \ldots, k_{n}$ are arbitrary constants, then the function $y$ in (1.2.10) satisfies (1.2.9).

Since the constants $k_{1}, k_{2}, \ldots, k_{n}$ in (1.2.10) are arbitrary, so are the constants

$$
\frac{k_{1}}{(n-1)!}, \frac{k_{2}}{(n-2)!}, \cdots, k_{n}
$$

Therefore Example 1.2.4 actually shows that all solutions of (1.2.9) can be written as

$$
y=\frac{e^{2 x}}{2^{n}}+c_{1}+c_{2} x+\cdots+c_{n} x^{n-1}
$$

where we renamed the arbitrary constants in (1.2.10) to obtain a simpler formula. As a general rule, arbitrary constants appearing in solutions of differential equations should be simplified if possible. You'll see examples of this throughout the text.
Initial Value Problems
In Example 1.2.4 we saw that the differential equation $y^{(n)}=e^{2 x}$ has an infinite family of solutions that depend upon the $n$ arbitrary constants $c_{1}, c_{2}, \ldots, c_{n}$. In the absence of additional conditions, there's no
reason to prefer one solution of a differential equation over another. However, we'll often be interested in finding a solution of a differential equation that satisfies one or more specific conditions. The next example illustrates this.

Example 1.2.5 Find a solution of

$$
y^{\prime}=x^{3}
$$

such that $y(1)=2$.

Solution At the beginning of this section we saw that the solutions of $y^{\prime}=x^{3}$ are

$$
y=\frac{x^{4}}{4}+c
$$

To determine a value of $c$ such that $y(1)=2$, we set $x=1$ and $y=2$ here to obtain

$$
2=y(1)=\frac{1}{4}+c, \quad \text { so } \quad c=\frac{7}{4}
$$

Therefore the required solution is

$$
y=\frac{x^{4}+7}{4}
$$

Figure 1.2.2 shows the graph of this solution. Note that imposing the condition $y(1)=2$ is equivalent to requiring the graph of $y$ to pass through the point $(1,2)$.

We can rewrite the problem considered in Example 1.2.5 more briefly as

$$
y^{\prime}=x^{3}, \quad y(1)=2 .
$$

We call this an initial value problem. The requirement $y(1)=2$ is an initial condition. Initial value problems can also be posed for higher order differential equations. For example,

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+3 y=e^{x}, \quad y(0)=1, \quad y^{\prime}(0)=2 \tag{1.2.11}
\end{equation*}
$$

is an initial value problem for a second order differential equation where $y$ and $y^{\prime}$ are required to have specified values at $x=0$. In general, an initial value problem for an $n$-th order differential equation requires $y$ and its first $n-1$ derivatives to have specified values at some point $x_{0}$. These requirements are the initial conditions.


Figure 1.2.2 $y=\frac{x^{2}+7}{4}$

We'll denote an initial value problem for a differential equation by writing the initial conditions after the equation, as in (1.2.11). For example, we would write an initial value problem for (1.2.2) as

$$
\begin{equation*}
y^{(n)}=f\left(x, y, y^{\prime}, \ldots, y^{(n-1)}\right), y\left(x_{0}\right)=k_{0}, y^{\prime}\left(x_{0}\right)=k_{1}, \ldots, y^{(n-1)}=k_{n-1} \tag{1.2.12}
\end{equation*}
$$

Consistent with our earlier definition of a solution of the differential equation in (1.2.12), we say that $y$ is a solution of the initial value problem (1.2.12) if $y$ is $n$ times differentiable and

$$
y^{(n)}(x)=f\left(x, y(x), y^{\prime}(x), \ldots, y^{(n-1)}(x)\right)
$$

for all $x$ in some open interval $(a, b)$ that contains $x_{0}$, and $y$ satisfies the initial conditions in (1.2.12). The largest open interval that contains $x_{0}$ on which $y$ is defined and satisfies the differential equation is the interval of validity of $y$.

Example 1.2.6 In Example 1.2 .5 we saw that

$$
\begin{equation*}
y=\frac{x^{4}+7}{4} \tag{1.2.13}
\end{equation*}
$$

is a solution of the initial value problem

$$
y^{\prime}=x^{3}, \quad y(1)=2 .
$$

Since the function in (1.2.13) is defined for all $x$, the interval of validity of this solution is $(-\infty, \infty)$.
Example 1.2.7 In Example 1.2.2 we verified that

$$
\begin{equation*}
y=\frac{x^{2}}{3}+\frac{1}{x} \tag{1.2.14}
\end{equation*}
$$

is a solution of

$$
x y^{\prime}+y=x^{2}
$$

on $(0, \infty)$ and on $(-\infty, 0)$. By evaluating (1.2.14) at $x= \pm 1$, you can see that (1.2.14) is a solution of the initial value problems

$$
\begin{equation*}
x y^{\prime}+y=x^{2}, \quad y(1)=\frac{4}{3} \tag{1.2.15}
\end{equation*}
$$

and

$$
\begin{equation*}
x y^{\prime}+y=x^{2}, \quad y(-1)=-\frac{2}{3} . \tag{1.2.16}
\end{equation*}
$$

The interval of validity of (1.2.14) as a solution of $(1.2 .15)$ is $(0, \infty)$, since this is the largest interval that contains $x_{0}=1$ on which (1.2.14) is defined. Similarly, the interval of validity of (1.2.14) as a solution of (1.2.16) is $(-\infty, 0)$, since this is the largest interval that contains $x_{0}=-1$ on which (1.2.14) is defined.

## Free Fall Under Constant Gravity

The term initial value problem originated in problems of motion where the independent variable is $t$ (representing elapsed time), and the initial conditions are the position and velocity of an object at the initial (starting) time of an experiment.

Example 1.2.8 An object falls under the influence of gravity near Earth's surface, where it can be assumed that the magnitude of the acceleration due to gravity is a constant $g$.
(a) Construct a mathematical model for the motion of the object in the form of an initial value problem for a second order differential equation, assuming that the altitude and velocity of the object at time $t=0$ are known. Assume that gravity is the only force acting on the object.
(b) Solve the initial value problem derived in (a) to obtain the altitude as a function of time.

SOLUTION(a) Let $y(t)$ be the altitude of the object at time $t$. Since the acceleration of the object has constant magnitude $g$ and is in the downward (negative) direction, $y$ satisfies the second order equation

$$
y^{\prime \prime}=-g,
$$

## 12 Chapter 1 Introduction

where the prime now indicates differentiation with respect to $t$. If $y_{0}$ and $v_{0}$ denote the altitude and velocity when $t=0$, then $y$ is a solution of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}=-g, \quad y(0)=y_{0}, \quad y^{\prime}(0)=v_{0} . \tag{1.2.17}
\end{equation*}
$$

$\underline{\text { SOLUTION(b) Integrating (1.2.17) twice yields }}$

$$
\begin{aligned}
y^{\prime} & =-g t+c_{1} \\
y & =-\frac{g t^{2}}{2}+c_{1} t+c_{2} .
\end{aligned}
$$

Imposing the initial conditions $y(0)=y_{0}$ and $y^{\prime}(0)=v_{0}$ in these two equations shows that $c_{1}=v_{0}$ and $c_{2}=y_{0}$. Therefore the solution of the initial value problem (1.2.17) is

$$
y=-\frac{g t^{2}}{2}+v_{0} t+y_{0}
$$

### 1.2 Exercises

1. Find the order of the equation.
(a) $\frac{d^{2} y}{d x^{2}}+2 \frac{d y}{d x} \frac{d^{3} y}{d x^{3}}+x=0$
(b) $y^{\prime \prime}-3 y^{\prime}+2 y=x^{7}$
(c) $y^{\prime}-y^{7}=0$
(d) $y^{\prime \prime} y-\left(y^{\prime}\right)^{2}=2$
2. Verify that the function is a solution of the differential equation on some interval, for any choice of the arbitrary constants appearing in the function.
(a) $y=c e^{2 x} ; \quad y^{\prime}=2 y$
(b) $y=\frac{x^{2}}{3}+\frac{c}{x} ; \quad x y^{\prime}+y=x^{2}$
(c) $y=\frac{1}{2}+c e^{-x^{2}} ; \quad y^{\prime}+2 x y=x$
(d) $y=\left(1+c e^{-x^{2} / 2}\right) ;\left(1-c e^{-x^{2} / 2}\right)^{-1} \quad 2 y^{\prime}+x\left(y^{2}-1\right)=0$
(e) $y=\tan \left(\frac{x^{3}}{3}+c\right) ; \quad y^{\prime}=x^{2}\left(1+y^{2}\right)$
(f) $y=\left(c_{1}+c_{2} x\right) e^{x}+\sin x+x^{2} ; \quad y^{\prime \prime}-2 y^{\prime}+y=-2 \cos x+x^{2}-4 x+2$
(g) $y=c_{1} e^{x}+c_{2} x+\frac{2}{x} ; \quad(1-x) y^{\prime \prime}+x y^{\prime}-y=4\left(1-x-x^{2}\right) x^{-3}$
(h) $y=x^{-1 / 2}\left(c_{1} \sin x+c_{2} \cos x\right)+4 x+8$;

$$
x^{2} y^{\prime \prime}+x y^{\prime}+\left(x^{2}-\frac{1}{4}\right) y=4 x^{3}+8 x^{2}+3 x-2
$$

3. Find all solutions of the equation.
(a) $y^{\prime}=-x$
(b) $y^{\prime}=-x \sin x$
(c) $y^{\prime}=x \ln x$
(d) $y^{\prime \prime}=x \cos x$
(e) $y^{\prime \prime}=2 x e^{x}$
(f) $y^{\prime \prime}=2 x+\sin x+e^{x}$
(g) $y^{\prime \prime \prime}=-\cos x$
(h) $y^{\prime \prime \prime}=-x^{2}+e^{x}$
(i) $y^{\prime \prime \prime}=7 e^{4 x}$
4. Solve the initial value problem.
(a) $y^{\prime}=-x e^{x}, \quad y(0)=1$
(b) $y^{\prime}=x \sin x^{2}, \quad y\left(\sqrt{\frac{\pi}{2}}\right)=1$
(c) $y^{\prime}=\tan x, \quad y(\pi / 4)=3$
(d) $y^{\prime \prime}=x^{4}, \quad y(2)=-1, \quad y^{\prime}(2)=-1$
(e) $y^{\prime \prime}=x e^{2 x}, \quad y(0)=7, \quad y^{\prime}(0)=1$
(f) $y^{\prime \prime}=-x \sin x, \quad y(0)=1, \quad y^{\prime}(0)=-3$
(g) $y^{\prime \prime \prime}=x^{2} e^{x}, \quad y(0)=1, \quad y^{\prime}(0)=-2, \quad y^{\prime \prime}(0)=3$
(h) $y^{\prime \prime \prime}=2+\sin 2 x, \quad y(0)=1, \quad y^{\prime}(0)=-6, \quad y^{\prime \prime}(0)=3$
(i) $y^{\prime \prime \prime}=2 x+1, \quad y(2)=1, \quad y^{\prime}(2)=-4, \quad y^{\prime \prime}(2)=7$
5. Verify that the function is a solution of the initial value problem.
(a) $y=x \cos x ; \quad y^{\prime}=\cos x-y \tan x, \quad y(\pi / 4)=\frac{\pi}{4 \sqrt{2}}$
(b) $y=\frac{1+2 \ln x}{x^{2}}+\frac{1}{2} ; \quad y^{\prime}=\frac{x^{2}-2 x^{2} y+2}{x^{3}}, \quad y(1)=\frac{3}{2}$
(c) $y=\tan \left(\frac{x^{2}}{2}\right) ; \quad y^{\prime}=x\left(1+y^{2}\right), \quad y(0)=0$
(d) $y=\frac{2}{x-2} ; \quad y^{\prime}=\frac{-y(y+1)}{x}, \quad y(1)=-2$
6. Verify that the function is a solution of the initial value problem.
(a) $y=x^{2}(1+\ln x) ; \quad y^{\prime \prime}=\frac{3 x y^{\prime}-4 y}{x^{2}}, \quad y(e)=2 e^{2}, \quad y^{\prime}(e)=5 e$
(b) $y=\frac{x^{2}}{3}+x-1 ; \quad y^{\prime \prime}=\frac{x^{2}-x y^{\prime}+y+1}{x^{2}}, \quad y(1)=\frac{1}{3}, \quad y^{\prime}(1)=\frac{5}{3}$
(c) $y=\left(1+x^{2}\right)^{-1 / 2} ; \quad y^{\prime \prime}=\frac{\left(x^{2}-1\right) y-x\left(x^{2}+1\right) y^{\prime}}{\left(x^{2}+1\right)^{2}}, \quad y(0)=1$, $y^{\prime}(0)=0$
(d) $y=\frac{x^{2}}{1-x} ; \quad y^{\prime \prime}=\frac{2(x+y)\left(x y^{\prime}-y\right)}{x^{3}}, \quad y(1 / 2)=1 / 2, \quad y^{\prime}(1 / 2)=3$
7. Suppose an object is launched from a point 320 feet above the earth with an initial velocity of 128 $\mathrm{ft} / \mathrm{sec}$ upward, and the only force acting on it thereafter is gravity. Take $g=32 \mathrm{ft} / \mathrm{sec}^{2}$.
(a) Find the highest altitude attained by the object.
(b) Determine how long it takes for the object to fall to the ground.
8. Let $a$ be a nonzero real number.
(a) Verify that if $c$ is an arbitrary constant then

$$
\begin{equation*}
y=(x-c)^{a} \tag{A}
\end{equation*}
$$

is a solution of

$$
\begin{equation*}
y^{\prime}=a y^{(a-1) / a} \tag{B}
\end{equation*}
$$

on $(c, \infty)$.
(b) Suppose $a<0$ or $a>1$. Can you think of a solution of (B) that isn't of the form (A)?
9. Verify that

$$
y=\left\{\begin{array}{cl}
e^{x}-1, & x \geq 0 \\
1-e^{-x}, & x<0
\end{array}\right.
$$

is a solution of

$$
y^{\prime}=|y|+1
$$

on $(-\infty, \infty)$. Hint: Use the definition of derivative at $x=0$.
10. (a) Verify that if $c$ is any real number then

$$
\begin{equation*}
y=c^{2}+c x+2 c+1 \tag{A}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
y^{\prime}=\frac{-(x+2)+\sqrt{x^{2}+4 x+4 y}}{2} \tag{B}
\end{equation*}
$$

on some open interval. Identify the open interval.
(b) Verify that

$$
y_{1}=\frac{-x(x+4)}{4}
$$

also satisfies (B) on some open interval, and identify the open interval. (Note that $y_{1}$ can't be obtained by selecting a value of $c$ in (A).)

### 1.3 DIRECTION FIELDS FOR FIRST ORDER EQUATIONS

It's impossible to find explicit formulas for solutions of some differential equations. Even if there are such formulas, they may be so complicated that they're useless. In this case we may resort to graphical or numerical methods to get some idea of how the solutions of the given equation behave.

In Section 2.3 we'll take up the question of existence of solutions of a first order equation

$$
\begin{equation*}
y^{\prime}=f(x, y) \tag{1.3.1}
\end{equation*}
$$

In this section we'll simply assume that (1.3.1) has solutions and discuss a graphical method for approximating them. In Chapter 3 we discuss numerical methods for obtaining approximate solutions of (1.3.1).

Recall that a solution of (1.3.1) is a function $y=y(x)$ such that

$$
y^{\prime}(x)=f(x, y(x))
$$

for all values of $x$ in some interval, and an integral curve is either the graph of a solution or is made up of segments that are graphs of solutions. Therefore, not being able to solve (1.3.1) is equivalent to not knowing the equations of integral curves of (1.3.1). However, it's easy to calculate the slopes of these curves. To be specific, the slope of an integral curve of (1.3.1) through a given point $\left(x_{0}, y_{0}\right)$ is given by the number $f\left(x_{0}, y_{0}\right)$. This is the basis of the method of direction fields.

If $f$ is defined on a set $R$, we can construct a direction field for (1.3.1) in $R$ by drawing a short line segment through each point $(x, y)$ in $R$ with slope $f(x, y)$. Of course, as a practical matter, we can't actually draw line segments through every point in $R$; rather, we must select a finite set of points in $R$. For example, suppose $f$ is defined on the closed rectangular region

$$
R:\{a \leq x \leq b, c \leq y \leq d\}
$$

Let

$$
a=x_{0}<x_{1}<\cdots<x_{m}=b
$$

be equally spaced points in $[a, b]$ and

$$
c=y_{0}<y_{1}<\cdots<y_{n}=d
$$

be equally spaced points in $[c, d]$. We say that the points

$$
\left(x_{i}, y_{j}\right), \quad 0 \leq i \leq m, \quad 0 \leq j \leq n,
$$

form a rectangular grid (Figure 1.3.1). Through each point in the grid we draw a short line segment with slope $f\left(x_{i}, y_{j}\right)$. The result is an approximation to a direction field for (1.3.1) in $R$. If the grid points are sufficiently numerous and close together, we can draw approximate integral curves of (1.3.1) by drawing curves through points in the grid tangent to the line segments associated with the points in the grid.


Figure 1.3.1 A rectangular grid

Unfortunately, approximating a direction field and graphing integral curves in this way is too tedious to be done effectively by hand. However, there is software for doing this. As you'll see, the combination of direction fields and integral curves gives useful insights into the behavior of the solutions of the differential equation even if we can't obtain exact solutions.

We'll study numerical methods for solving a single first order equation (1.3.1) in Chapter 3. These methods can be used to plot solution curves of (1.3.1) in a rectangular region $R$ if $f$ is continuous on $R$. Figures 1.3.2, 1.3.3, and 1.3.4 show direction fields and solution curves for the differential equations

$$
y^{\prime}=\frac{x^{2}-y^{2}}{1+x^{2}+y^{2}}, \quad y^{\prime}=1+x y^{2}, \quad \text { and } \quad y^{\prime}=\frac{x-y}{1+x^{2}}
$$

which are all of the form (1.3.1) with $f$ continuous for all $(x, y)$.


Figure 1.3.2 A direction field and integral curves

$$
\text { for } y=\frac{x^{2}-y^{2}}{1+x^{2}+y^{2}}
$$



Figure 1.3.3 A direction field and integral curves for

$$
y^{\prime}=1+x y^{2}
$$

The methods of Chapter 3 won't work for the equation

$$
\begin{equation*}
y^{\prime}=-x / y \tag{1.3.2}
\end{equation*}
$$

if $R$ contains part of the $x$-axis, since $f(x, y)=-x / y$ is undefined when $y=0$. Similarly, they won't work for the equation

$$
\begin{equation*}
y^{\prime}=\frac{x^{2}}{1-x^{2}-y^{2}} \tag{1.3.3}
\end{equation*}
$$

if $R$ contains any part of the unit circle $x^{2}+y^{2}=1$, because the right side of (1.3.3) is undefined if $x^{2}+y^{2}=1$. However, (1.3.2) and (1.3.3) can written as

$$
\begin{equation*}
y^{\prime}=\frac{A(x, y)}{B(x, y)} \tag{1.3.4}
\end{equation*}
$$

where $A$ and $B$ are continuous on any rectangle $R$. Because of this, some differential equation software is based on numerically solving pairs of equations of the form

$$
\begin{equation*}
\frac{d x}{d t}=B(x, y), \quad \frac{d y}{d t}=A(x, y) \tag{1.3.5}
\end{equation*}
$$

where $x$ and $y$ are regarded as functions of a parameter $t$. If $x=x(t)$ and $y=y(t)$ satisfy these equations, then

$$
y^{\prime}=\frac{d y}{d x}=\frac{d y}{d t} / \frac{d x}{d t}=\frac{A(x, y)}{B(x, y)},
$$

so $y=y(x)$ satisfies (1.3.4).
Eqns. (1.3.2) and (1.3.3) can be reformulated as in (1.3.4) with

$$
\frac{d x}{d t}=-y, \quad \frac{d y}{d t}=x
$$



Figure 1.3.4 A direction and integral curves for $y^{\prime}=\frac{x-y}{1+x^{2}}$
and

$$
\frac{d x}{d t}=1-x^{2}-y^{2}, \quad \frac{d y}{d t}=x^{2}
$$

respectively. Even if $f$ is continuous and otherwise "nice" throughout $R$, your software may require you to reformulate the equation $y^{\prime}=f(x, y)$ as

$$
\frac{d x}{d t}=1, \quad \frac{d y}{d t}=f(x, y)
$$

which is of the form (1.3.5) with $A(x, y)=f(x, y)$ and $B(x, y)=1$.
Figure 1.3.5 shows a direction field and some integral curves for (1.3.2). As we saw in Example 1.2.1 and will verify again in Section 2.2, the integral curves of (1.3.2) are circles centered at the origin.


Figure 1.3.5 A direction field and integral curves for $y^{\prime}=-\frac{x}{y}$
Figure 1.3.6 shows a direction field and some integral curves for (1.3.3). The integral curves near the top and bottom are solution curves. However, the integral curves near the middle are more complicated. For example, Figure 1.3.7 shows the integral curve through the origin. The vertices of the dashed rectangle are on the circle $x^{2}+y^{2}=1(a \approx .846, b \approx .533)$, where all integral curves of (1.3.3) have infinite slope. There are three solution curves of (1.3.3) on the integral curve in the figure: the segment above the level $y=b$ is the graph of a solution on $(-\infty, a)$, the segment below the level $y=-b$ is the graph of a solution on $(-a, \infty)$, and the segment between these two levels is the graph of a solution on $(-a, a)$.

## USING TECHNOLOGY

As you study from this book, you'll often be asked to use computer software and graphics. Exercises with this intent are marked as C (computer or calculator required), $\mathrm{C} / \mathrm{G}$ (computer and/or graphics required), or L (laboratory work requiring software and/or graphics). Often you may not completely understand how the software does what it does. This is similar to the situation most people are in when they drive automobiles or watch television, and it doesn't decrease the value of using modern technology as an aid to learning. Just be careful that you use the technology as a supplement to thought rather than a substitute for it.


Figure 1.3.6 A direction field and integral curves for

$$
y^{\prime}=\frac{x^{2}}{1-x^{2}-y^{2}}
$$



Figure 1.3.7
1.3 Exercises

In Exercises 1-11 a direction field is drawn for the given equation. Sketch some integral curves.










10 A direction field for $y^{\prime}=x^{3} y^{2}+x y^{3}$


In Exercises 12-22 construct a direction field and plot some integral curves in the indicated rectangular region.
12. $\mathrm{C} / \mathrm{G} y^{\prime}=y(y-1) ; \quad\{-1 \leq x \leq 2,-2 \leq y \leq 2\}$
13. C/G $y^{\prime}=2-3 x y ; \quad\{-1 \leq x \leq 4,-4 \leq y \leq 4\}$
14. $\mathrm{C} / \mathrm{G} y^{\prime}=x y(y-1) ; \quad\{-2 \leq x \leq 2,-4 \leq y \leq 4\}$
15. $\mathrm{C} / \mathrm{G} \quad y^{\prime}=3 x+y ; \quad\{-2 \leq x \leq 2,0 \leq y \leq 4\}$
16. $\mathrm{C} / \mathrm{G} y^{\prime}=y-x^{3} ; \quad\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$
17. $\mathrm{C} / \mathrm{G} y^{\prime}=1-x^{2}-y^{2} ; \quad\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$
18. C/G $y^{\prime}=x\left(y^{2}-1\right) ; \quad\{-3 \leq x \leq 3,-3 \leq y \leq 2\}$
19. $\mathrm{C} / \mathrm{G} y^{\prime}=\frac{x}{y\left(y^{2}-1\right)} ; \quad\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$
20. $\mathrm{C} / \mathrm{G} y^{\prime}=\frac{x y^{2}}{y-1} ; \quad\{-2 \leq x \leq 2,-1 \leq y \leq 4\}$
21. $\mathrm{C} / \mathrm{G} y^{\prime}=\frac{x\left(y^{2}-1\right)}{y} ; \quad\{-1 \leq x \leq 1,-2 \leq y \leq 2\}$
22. $\mathrm{C} / \mathrm{G} y^{\prime}=-\frac{x^{2}+y^{2}}{1-x^{2}-y^{2}} ; \quad\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$
23. L By suitably renaming the constants and dependent variables in the equations

$$
\begin{equation*}
T^{\prime}=-k\left(T-T_{m}\right) \tag{A}
\end{equation*}
$$

and

$$
\begin{equation*}
G^{\prime}=-\lambda G+r \tag{B}
\end{equation*}
$$

discussed in Section 1.2 in connection with Newton's law of cooling and absorption of glucose in the body, we can write both as

$$
\begin{equation*}
y^{\prime}=-a y+b \tag{C}
\end{equation*}
$$

where $a$ is a positive constant and $b$ is an arbitrary constant. Thus, (A) is of the form (C) with $y=T, a=k$, and $b=k T_{m}$, and (B) is of the form (C) with $y=G, a=\lambda$, and $b=r$. We'll encounter equations of the form (C) in many other applications in Chapter 2.
Choose a positive $a$ and an arbitrary $b$. Construct a direction field and plot some integral curves for $(\mathrm{C})$ in a rectangular region of the form

$$
\{0 \leq t \leq T, c \leq y \leq d\}
$$

of the $t y$-plane. Vary $T, c$, and $d$ until you discover a common property of all the solutions of (C). Repeat this experiment with various choices of $a$ and $b$ until you can state this property precisely in terms of $a$ and $b$.
24. L By suitably renaming the constants and dependent variables in the equations

$$
\begin{equation*}
P^{\prime}=a P(1-\alpha P) \tag{A}
\end{equation*}
$$

and

$$
\begin{equation*}
I^{\prime}=r I(S-I) \tag{B}
\end{equation*}
$$

discussed in Section 1.1 in connection with Verhulst's population model and the spread of an epidemic, we can write both in the form

$$
\begin{equation*}
y^{\prime}=a y-b y^{2} \tag{C}
\end{equation*}
$$

where $a$ and $b$ are positive constants. Thus, (A) is of the form (C) with $y=P, a=a$, and $b=a \alpha$, and (B) is of the form (C) with $y=I, a=r S$, and $b=r$. In Chapter 2 we'll encounter equations of the form (C) in other applications..
(a) Choose positive numbers $a$ and $b$. Construct a direction field and plot some integral curves for $(\mathrm{C})$ in a rectangular region of the form

$$
\{0 \leq t \leq T, 0 \leq y \leq d\}
$$

of the $t y$-plane. Vary $T$ and $d$ until you discover a common property of all solutions of (C) with $y(0)>0$. Repeat this experiment with various choices of $a$ and $b$ until you can state this property precisely in terms of $a$ and $b$.
(b) Choose positive numbers $a$ and $b$. Construct a direction field and plot some integral curves for $(\mathrm{C})$ in a rectangular region of the form

$$
\{0 \leq t \leq T, c \leq y \leq 0\}
$$

of the $t y$-plane. Vary $a, b, T$ and $c$ until you discover a common property of all solutions of (C) with $y(0)<0$.

You can verify your results later by doing Exercise 2.2.27.

## CHAPTER 2 First Order Equations

IN THIS CHAPTER we study first order equations for which there are general methods of solution.
SECTION 2.1 deals with linear equations, the simplest kind of first order equations. In this section we introduce the method of variation of parameters. The idea underlying this method will be a unifying theme for our approach to solving many different kinds of differential equations throughout the book.

SECTION 2.2 deals with separable equations, the simplest nonlinear equations. In this section we introduce the idea of implicit and constant solutions of differential equations, and we point out some differences between the properties of linear and nonlinear equations.

SECTION 2.3 discusses existence and uniqueness of solutions of nonlinear equations. Although it may seem logical to place this section before Section 2.2, we presented Section 2.2 first so we could have illustrative examples in Section 2.3.
SECTION 2.4 deals with nonlinear equations that are not separable, but can be transformed into separable equations by a procedure similar to variation of parameters.

SECTION 2.5 covers exact differential equations, which are given this name because the method for solving them uses the idea of an exact differential from calculus.

SECTION 2.6 deals with equations that are not exact, but can made exact by multiplying them by a function known called integrating factor.

### 2.1 LINEAR FIRST ORDER EQUATIONS

A first order differential equation is said to be linear if it can be written as

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{2.1.1}
\end{equation*}
$$

A first order differential equation that can't be written like this is nonlinear. We say that (2.1.1) is homogeneous if $f \equiv 0$; otherwise it's nonhomogeneous. Since $y \equiv 0$ is obviously a solution of the homgeneous equation

$$
y^{\prime}+p(x) y=0
$$

we call it the trivial solution. Any other solution is nontrivial.
Example 2.1.1 The first order equations

$$
\begin{aligned}
x^{2} y^{\prime}+3 y & =x^{2}, \\
x y^{\prime}-8 x^{2} y & =\sin x, \\
x y^{\prime}+(\ln x) y & =0, \\
y^{\prime} & =x^{2} y-2,
\end{aligned}
$$

are not in the form (2.1.1), but they are linear, since they can be rewritten as

$$
\begin{aligned}
y^{\prime}+\frac{3}{x^{2}} y & =1 \\
y^{\prime}-8 x y & =\frac{\sin x}{x} \\
y^{\prime}+\frac{\ln x}{x} y & =0 \\
y^{\prime}-x^{2} y & =-2 .
\end{aligned}
$$

Example 2.1.2 Here are some nonlinear first order equations:

$$
\begin{array}{rlc}
x y^{\prime}+3 y^{2} & =2 x & \text { (because } y \text { is squared) }, \\
y y^{\prime} & =3 & \text { (because of the product } \left.y y^{\prime}\right), \\
y^{\prime}+x e^{y} & =12 & \text { (because of } \left.e^{y}\right) .
\end{array}
$$

General Solution of a Linear First Order Equation
To motivate a definition that we'll need, consider the simple linear first order equation

$$
\begin{equation*}
y^{\prime}=\frac{1}{x^{2}} \tag{2.1.2}
\end{equation*}
$$

From calculus we know that $y$ satisfies this equation if and only if

$$
\begin{equation*}
y=-\frac{1}{x}+c \tag{2.1.3}
\end{equation*}
$$

where $c$ is an arbitrary constant. We call $c$ a parameter and say that (2.1.3) defines a one-parameter family of functions. For each real number $c$, the function defined by (2.1.3) is a solution of (2.1.2) on $(-\infty, 0)$ and $(0, \infty)$; moreover, every solution of (2.1.2) on either of these intervals is of the form (2.1.3) for some choice of $c$. We say that (2.1.3) is the general solution of (2.1.2).

We'll see that a similar situation occurs in connection with any first order linear equation

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{2.1.4}
\end{equation*}
$$

that is, if $p$ and $f$ are continuous on some open interval $(a, b)$ then there's a unique formula $y=y(x, c)$ analogous to (2.1.3) that involves $x$ and a parameter $c$ and has the these properties:

- For each fixed value of $c$, the resulting function of $x$ is a solution of (2.1.4) on $(a, b)$.
- If $y$ is a solution of (2.1.4) on $(a, b)$, then $y$ can be obtained from the formula by choosing $c$ appropriately.
We'll call $y=y(x, c)$ the general solution of (2.1.4).
When this has been established, it will follow that an equation of the form

$$
\begin{equation*}
P_{0}(x) y^{\prime}+P_{1}(x) y=F(x) \tag{2.1.5}
\end{equation*}
$$

has a general solution on any open interval $(a, b)$ on which $P_{0}, P_{1}$, and $F$ are all continuous and $P_{0}$ has no zeros, since in this case we can rewrite (2.1.5) in the form (2.1.4) with $p=P_{1} / P_{0}$ and $f=F / P_{0}$, which are both continuous on $(a, b)$.

To avoid awkward wording in examples and exercises, we won't specify the interval $(a, b)$ when we ask for the general solution of a specific linear first order equation. Let's agree that this always means that we want the general solution on every open interval on which $p$ and $f$ are continuous if the equation is of the form (2.1.4), or on which $P_{0}, P_{1}$, and $F$ are continuous and $P_{0}$ has no zeros, if the equation is of the form (2.1.5). We leave it to you to identify these intervals in specific examples and exercises.

For completeness, we point out that if $P_{0}, P_{1}$, and $F$ are all continuous on an open interval $(a, b)$, but $P_{0}$ does have a zero in $(a, b)$, then (2.1.5) may fail to have a general solution on $(a, b)$ in the sense just defined. Since this isn't a major point that needs to be developed in depth, we won't discuss it further; however, see Exercise 44 for an example.

Homogeneous Linear First Order Equations
We begin with the problem of finding the general solution of a homogeneous linear first order equation. The next example recalls a familiar result from calculus.

Example 2.1.3 Let $a$ be a constant.
(a) Find the general solution of

$$
\begin{equation*}
y^{\prime}-a y=0 \tag{2.1.6}
\end{equation*}
$$

(b) Solve the initial value problem

$$
y^{\prime}-a y=0, \quad y\left(x_{0}\right)=y_{0}
$$

SOLUTION(a) You already know from calculus that if $c$ is any constant, then $y=c e^{a x}$ satisfies (2.1.6). However, let's pretend you've forgotten this, and use this problem to illustrate a general method for solving a homogeneous linear first order equation.

We know that (2.1.6) has the trivial solution $y \equiv 0$. Now suppose $y$ is a nontrivial solution of (2.1.6). Then, since a differentiable function must be continuous, there must be some open interval $I$ on which $y$ has no zeros. We rewrite (2.1.6) as

$$
\frac{y^{\prime}}{y}=a
$$

for $x$ in $I$. Integrating this shows that

$$
\ln |y|=a x+k, \quad \text { so } \quad|y|=e^{k} e^{a x}
$$

where $k$ is an arbitrary constant. Since $e^{a x}$ can never equal zero, $y$ has no zeros, so $y$ is either always positive or always negative. Therefore we can rewrite $y$ as

$$
\begin{equation*}
y=c e^{a x} \tag{2.1.7}
\end{equation*}
$$

where

$$
c=\left\{\begin{aligned}
e^{k} & \text { if } y>0 \\
-e^{k} & \text { if } y<0
\end{aligned}\right.
$$

This shows that every nontrivial solution of (2.1.6) is of the form $y=c e^{a x}$ for some nonzero constant $c$. Since setting $c=0$ yields the trivial solution, all solutions of (2.1.6) are of the form (2.1.7). Conversely, (2.1.7) is a solution of (2.1.6) for every choice of $c$, since differentiating (2.1.7) yields $y^{\prime}=a c e^{a x}=a y$.
$\underline{\text { SOLUTION(b) }}$ Imposing the initial condition $y\left(x_{0}\right)=y_{0}$ yields $y_{0}=c e^{a x_{0}}$, so $c=y_{0} e^{-a x_{0}}$ and

$$
y=y_{0} e^{-a x_{0}} e^{a x}=y_{0} e^{a\left(x-x_{0}\right)}
$$

Figure 2.1.1 show the graphs of this function with $x_{0}=0, y_{0}=1$, and various values of $a$.


Figure 2.1.1 Solutions of $y^{\prime}-a y=0, y(0)=1$

Example 2.1.4 (a) Find the general solution of

$$
\begin{equation*}
x y^{\prime}+y=0 . \tag{2.1.8}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
x y^{\prime}+y=0, \quad y(1)=3 . \tag{2.1.9}
\end{equation*}
$$

SOLUTION(a) We rewrite (2.1.8) as

$$
\begin{equation*}
y^{\prime}+\frac{1}{x} y=0 \tag{2.1.10}
\end{equation*}
$$

where $x$ is restricted to either $(-\infty, 0)$ or $(0, \infty)$. If $y$ is a nontrivial solution of (2.1.10), there must be some open interval I on which $y$ has no zeros. We can rewrite (2.1.10) as

$$
\frac{y^{\prime}}{y}=-\frac{1}{x}
$$

for $x$ in $I$. Integrating shows that

$$
\ln |y|=-\ln |x|+k, \quad \text { so } \quad|y|=\frac{e^{k}}{|x|}
$$

Since a function that satisfies the last equation can't change sign on either $(-\infty, 0)$ or $(0, \infty)$, we can rewrite this result more simply as

$$
\begin{equation*}
y=\frac{c}{x} \tag{2.1.11}
\end{equation*}
$$

where

$$
c=\left\{\begin{aligned}
e^{k} & \text { if } y>0 \\
-e^{k} & \text { if } y<0
\end{aligned}\right.
$$

We've now shown that every solution of (2.1.10) is given by (2.1.11) for some choice of $c$. (Even though we assumed that $y$ was nontrivial to derive (2.1.11), we can get the trivial solution by setting $c=0$ in (2.1.11).) Conversely, any function of the form (2.1.11) is a solution of (2.1.10), since differentiating (2.1.11) yields

$$
y^{\prime}=-\frac{c}{x^{2}},
$$



Figure 2.1.2 Solutions of $x y^{\prime}+y=0$ on $(0, \infty)$ and $(-\infty, 0)$
and substituting this and (2.1.11) into (2.1.10) yields

$$
\begin{aligned}
y^{\prime}+\frac{1}{x} y & =-\frac{c}{x^{2}}+\frac{1}{x} \frac{c}{x} \\
& =-\frac{c}{x^{2}}+\frac{c}{x^{2}}=0 .
\end{aligned}
$$

Figure 2.1.2 shows the graphs of some solutions corresponding to various values of $c$
$\underline{\text { SOLUTION }(\mathbf{b})}$ Imposing the initial condition $y(1)=3$ in (2.1.11) yields $c=3$. Therefore the solution of (2.1.9) is

$$
y=\frac{3}{x} .
$$

The interval of validity of this solution is $(0, \infty)$.
The results in Examples 2.1.3(a) and 2.1.4(b) are special cases of the next theorem.
Theorem 2.1.1 If p is continuous on $(a, b)$, then the general solution of the homogeneous equation

$$
\begin{equation*}
y^{\prime}+p(x) y=0 \tag{2.1.12}
\end{equation*}
$$

on $(a, b)$ is

$$
y=c e^{-P(x)},
$$

where

$$
\begin{equation*}
P(x)=\int p(x) d x \tag{2.1.13}
\end{equation*}
$$

is any antiderivative of $p$ on $(a, b)$; that is,

$$
\begin{equation*}
P^{\prime}(x)=p(x), \quad a<x<b . \tag{2.1.14}
\end{equation*}
$$

Proof If $y=c e^{-P(x)}$, differentiating $y$ and using (2.1.14) shows that

$$
y^{\prime}=-P^{\prime}(x) c e^{-P(x)}=-p(x) c e^{-P(x)}=-p(x) y,
$$

so $y^{\prime}+p(x) y=0$; that is, $y$ is a solution of (2.1.12), for any choice of $c$.

Now we'll show that any solution of (2.1.12) can be written as $y=c e^{-P(x)}$ for some constant $c$. The trivial solution can be written this way, with $c=0$. Now suppose $y$ is a nontrivial solution. Then there's an open subinterval $I$ of $(a, b)$ on which $y$ has no zeros. We can rewrite (2.1.12) as

$$
\begin{equation*}
\frac{y^{\prime}}{y}=-p(x) \tag{2.1.15}
\end{equation*}
$$

for $x$ in $I$. Integrating (2.1.15) and recalling (2.1.13) yields

$$
\ln |y|=-P(x)+k,
$$

where $k$ is a constant. This implies that

$$
|y|=e^{k} e^{-P(x)}
$$

Since $P$ is defined for all $x$ in $(a, b)$ and an exponential can never equal zero, we can take $I=(a, b)$, so $y$ has zeros on $(a, b)(a, b)$, so we can rewrite the last equation as $y=c e^{-P(x)}$, where

$$
c=\left\{\begin{aligned}
e^{k} & \text { if } y>0 \text { on }(a, b), \\
-e^{k} & \text { if } y<0 \text { on }(a, b) .
\end{aligned}\right.
$$

REMARK: Rewriting a first order differential equation so that one side depends only on $y$ and $y^{\prime}$ and the other depends only on $x$ is called separation of variables. We did this in Examples 2.1.3 and 2.1.4, and in rewriting (2.1.12) as (2.1.15).We'llapply this method to nonlinear equations in Section 2.2.

Linear Nonhomogeneous First Order Equations
We'll now solve the nonhomogeneous equation

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{2.1.16}
\end{equation*}
$$

When considering this equation we call

$$
y^{\prime}+p(x) y=0
$$

the complementary equation.
We'll find solutions of (2.1.16) in the form $y=u y_{1}$, where $y_{1}$ is a nontrivial solution of the complementary equation and $u$ is to be determined. This method of using a solution of the complementary equation to obtain solutions of a nonhomogeneous equation is a special case of a method called variation of parameters, which you'll encounter several times in this book. (Obviously, $u$ can't be constant, since if it were, the left side of (2.1.16) would be zero. Recognizing this, the early users of this method viewed $u$ as a "parameter" that varies; hence, the name "variation of parameters.")
If

$$
y=u y_{1}, \quad \text { then } \quad y^{\prime}=u^{\prime} y_{1}+u y_{1}^{\prime} .
$$

Substituting these expressions for $y$ and $y^{\prime}$ into (2.1.16) yields

$$
u^{\prime} y_{1}+u\left(y_{1}^{\prime}+p(x) y_{1}\right)=f(x)
$$

which reduces to

$$
\begin{equation*}
u^{\prime} y_{1}=f(x), \tag{2.1.17}
\end{equation*}
$$

since $y_{1}$ is a solution of the complementary equation; that is,

$$
y_{1}^{\prime}+p(x) y_{1}=0 .
$$

In the proof of Theorem 2.2.1 we saw that $y_{1}$ has no zeros on an interval where $p$ is continuous. Therefore we can divide (2.1.17) through by $y_{1}$ to obtain

$$
u^{\prime}=f(x) / y_{1}(x) .
$$

We can integrate this (introducing a constant of integration), and multiply the result by $y_{1}$ to get the general solution of (2.1.16). Before turning to the formal proof of this claim, let's consider some examples.


Figure 2.1.3 A direction field and integral curves for $y^{\prime}+2 y=x^{2} e^{-2 x}$

Example 2.1.5 Find the general solution of

$$
\begin{equation*}
y^{\prime}+2 y=x^{3} e^{-2 x} \tag{2.1.18}
\end{equation*}
$$

By applying (a) of Example 2.1.3 with $a=-2$, we see that $y_{1}=e^{-2 x}$ is a solution of the complementary equation $y^{\prime}+2 y=0$. Therefore we seek solutions of (2.1.18) in the form $y=u e^{-2 x}$, so that

$$
\begin{equation*}
y^{\prime}=u^{\prime} e^{-2 x}-2 u e^{-2 x} \quad \text { and } \quad y^{\prime}+2 y=u^{\prime} e^{-2 x}-2 u e^{-2 x}+2 u e^{-2 x}=u^{\prime} e^{-2 x} \tag{2.1.19}
\end{equation*}
$$

Therefore $y$ is a solution of (2.1.18) if and only if

$$
u^{\prime} e^{-2 x}=x^{3} e^{-2 x} \quad \text { or, equivalently, } \quad u^{\prime}=x^{3}
$$

Therefore

$$
u=\frac{x^{4}}{4}+c
$$

and

$$
y=u e^{-2 x}=e^{-2 x}\left(\frac{x^{4}}{4}+c\right)
$$

is the general solution of (2.1.18).
Figure 2.1.3 shows a direction field and some integral curves for (2.1.18).

## Example 2.1.6

(a) Find the general solution

$$
\begin{equation*}
y^{\prime}+(\cot x) y=x \csc x \tag{2.1.20}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime}+(\cot x) y=x \csc x, \quad y(\pi / 2)=1 \tag{2.1.21}
\end{equation*}
$$

SOLUTION(a) Here $p(x)=\cot x$ and $f(x)=x \csc x$ are both continuous except at the points $x=r \pi$, where $r$ is an integer. Therefore we seek solutions of (2.1.20) on the intervals $(r \pi,(r+1) \pi)$. We need a nontrival solution $y_{1}$ of the complementary equation; thus, $y_{1}$ must satisfy $y_{1}^{\prime}+(\cot x) y_{1}=0$, which we rewrite as

$$
\begin{equation*}
\frac{y_{1}^{\prime}}{y_{1}}=-\cot x=-\frac{\cos x}{\sin x} \tag{2.1.22}
\end{equation*}
$$

Integrating this yields

$$
\ln \left|y_{1}\right|=-\ln |\sin x|,
$$

where we take the constant of integration to be zero since we need only one function that satisfies (2.1.22). Clearly $y_{1}=1 / \sin x$ is a suitable choice. Therefore we seek solutions of (2.1.20) in the form

$$
y=\frac{u}{\sin x},
$$

so that

$$
\begin{equation*}
y^{\prime}=\frac{u^{\prime}}{\sin x}-\frac{u \cos x}{\sin ^{2} x} \tag{2.1.23}
\end{equation*}
$$

and

$$
\begin{align*}
y^{\prime}+(\cot x) y & =\frac{u^{\prime}}{\sin x}-\frac{u \cos x}{\sin ^{2} x}+\frac{u \cot x}{\sin x} \\
& =\frac{u^{\prime}}{\sin x}-\frac{u \cos x}{\sin ^{2} x}+\frac{u \cos x}{\sin ^{2} x}  \tag{2.1.24}\\
& =\frac{u^{\prime}}{\sin x}
\end{align*}
$$

Therefore $y$ is a solution of (2.1.20) if and only if

$$
u^{\prime} / \sin x=x \csc x=x / \sin x \quad \text { or, equivalently, } \quad u^{\prime}=x
$$

Integrating this yields

$$
\begin{equation*}
u=\frac{x^{2}}{2}+c, \quad \text { and } \quad y=\frac{u}{\sin x}=\frac{x^{2}}{2 \sin x}+\frac{c}{\sin x} \tag{2.1.25}
\end{equation*}
$$

is the general solution of (2.1.20) on every interval $(r \pi,(r+1) \pi)(r=$ integer $)$.
$\underline{\operatorname{SOLUTION}(b)}$ Imposing the initial condition $y(\pi / 2)=1$ in (2.1.25) yields

$$
1=\frac{\pi^{2}}{8}+c \quad \text { or } \quad c=1-\frac{\pi^{2}}{8}
$$

Thus,

$$
y=\frac{x^{2}}{2 \sin x}+\frac{\left(1-\pi^{2} / 8\right)}{\sin x}
$$

is a solution of (2.1.21). The interval of validity of this solution is $(0, \pi)$; Figure 2.1.4 shows its graph.


Figure 2.1.4 Solution of $y^{\prime}+(\cot x) y=x \csc x, y(\pi / 2)=1$

REMARK: It wasn't necessary to do the computations (2.1.23) and (2.1.24) in Example 2.1.6, since we showed in the discussion preceding Example 2.1.5 that if $y=u y_{1}$ where $y_{1}^{\prime}+p(x) y_{1}=0$, then $y^{\prime}+$ $p(x) y=u^{\prime} y_{1}$. We did these computations so you would see this happen in this specific example. We

## 34 Chapter 2 First Order Equations

recommend that you include these "unnecesary" computations in doing exercises, until you're confident that you really understand the method. After that, omit them.

We summarize the method of variation of parameters for solving

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{2.1.26}
\end{equation*}
$$

as follows:
(a) Find a function $y_{1}$ such that

$$
\frac{y_{1}^{\prime}}{y_{1}}=-p(x)
$$

For convenience, take the constant of integration to be zero.
(b) Write

$$
\begin{equation*}
y=u y_{1} \tag{2.1.27}
\end{equation*}
$$

to remind yourself of what you're doing.
(c) Write $u^{\prime} y_{1}=f$ and solve for $u^{\prime}$; thus, $u^{\prime}=f / y_{1}$.
(d) Integrate $u^{\prime}$ to obtain $u$, with an arbitrary constant of integration.
(e) Substitute $u$ into (2.1.27) to obtain $y$.

To solve an equation written as

$$
P_{0}(x) y^{\prime}+P_{1}(x) y=F(x),
$$

we recommend that you divide through by $P_{0}(x)$ to obtain an equation of the form (2.1.26) and then follow this procedure.
Solutions in Integral Form
Sometimes the integrals that arise in solving a linear first order equation can't be evaluated in terms of elementary functions. In this case the solution must be left in terms of an integral.

## Example 2.1.7

(a) Find the general solution of

$$
y^{\prime}-2 x y=1
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime}-2 x y=1, \quad y(0)=y_{0} . \tag{2.1.28}
\end{equation*}
$$

SOLUTION(a) To apply variation of parameters, we need a nontrivial solution $y_{1}$ of the complementary equation; thus, $y_{1}^{\prime}-2 x y_{1}=0$, which we rewrite as

$$
\frac{y_{1}^{\prime}}{y_{1}}=2 x .
$$

Integrating this and taking the constant of integration to be zero yields

$$
\ln \left|y_{1}\right|=x^{2}, \quad \text { so } \quad\left|y_{1}\right|=e^{x^{2}}
$$

We choose $y_{1}=e^{x^{2}}$ and seek solutions of (2.1.28) in the form $y=u e^{x^{2}}$, where

$$
u^{\prime} e^{x^{2}}=1, \quad \text { so } \quad u^{\prime}=e^{-x^{2}}
$$

Therefore

$$
u=c+\int e^{-x^{2}} d x
$$

but we can't simplify the integral on the right because there's no elementary function with derivative equal to $e^{-x^{2}}$. Therefore the best available form for the general solution of (2.1.28) is

$$
\begin{equation*}
y=u e^{x^{2}}=e^{x^{2}}\left(c+\int e^{-x^{2}} d x\right) . \tag{2.1.29}
\end{equation*}
$$

$\underline{\text { SOLUTION(b) }}$ Since the initial condition in (2.1.28) is imposed at $x_{0}=0$, it is convenient to rewrite (2.1.29) as

$$
y=e^{x^{2}}\left(c+\int_{0}^{x} e^{-t^{2}} d t\right), \quad \text { since } \quad \int_{0}^{0} e^{-t^{2}} d t=0
$$

Setting $x=0$ and $y=y_{0}$ here shows that $c=y_{0}$. Therefore the solution of the initial value problem is

$$
\begin{equation*}
y=e^{x^{2}}\left(y_{0}+\int_{0}^{x} e^{-t^{2}} d t\right) \tag{2.1.30}
\end{equation*}
$$

For a given value of $y_{0}$ and each fixed $x$, the integral on the right can be evaluated by numerical methods. An alternate procedure is to apply the numerical integration procedures discussed in Chapter 3 directly to the initial value problem (2.1.28). Figure 2.1 .5 shows graphs of of (2.1.30) for several values of $y_{0}$.


Figure 2.1.5 Solutions of $y^{\prime}-2 x y=1, y(0)=y_{0}$

An Existence and Uniqueness Theorem
The method of variation of parameters leads to this theorem.
Theorem 2.1.2 Suppose $p$ and $f$ are continuous on an open interval $(a, b)$, and let $y_{1}$ be any nontrivial solution of the complementary equation

$$
y^{\prime}+p(x) y=0
$$

on $(a, b)$. Then:
(a) The general solution of the nonhomogeneous equation

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{2.1.31}
\end{equation*}
$$

on $(a, b)$ is

$$
\begin{equation*}
y=y_{1}(x)\left(c+\int f(x) / y_{1}(x) d x\right) . \tag{2.1.32}
\end{equation*}
$$

(b) If $x_{0}$ is an arbitrary point in $(a, b)$ and $y_{0}$ is an arbitrary real number, then the initial value problem

$$
y^{\prime}+p(x) y=f(x), \quad y\left(x_{0}\right)=y_{0}
$$

has the unique solution

$$
y=y_{1}(x)\left(\frac{y_{0}}{y_{1}\left(x_{0}\right)}+\int_{x_{0}}^{x} \frac{f(t)}{y_{1}(t)} d t\right)
$$

on $(a, b)$.

Proof (a) To show that (2.1.32) is the general solution of (2.1.31) on $(a, b)$, we must prove that:
(i) If $c$ is any constant, the function $y$ in (2.1.32) is a solution of (2.1.31) on $(a, b)$.
(ii) If $y$ is a solution of (2.1.31) on $(a, b)$ then $y$ is of the form (2.1.32) for some constant $c$.

To prove (i), we first observe that any function of the form (2.1.32) is defined on $(a, b)$, since $p$ and $f$ are continuous on $(a, b)$. Differentiating (2.1.32) yields

$$
y^{\prime}=y_{1}^{\prime}(x)\left(c+\int f(x) / y_{1}(x) d x\right)+f(x)
$$

Since $y_{1}^{\prime}=-p(x) y_{1}$, this and (2.1.32) imply that

$$
\begin{aligned}
y^{\prime} & =-p(x) y_{1}(x)\left(c+\int f(x) / y_{1}(x) d x\right)+f(x) \\
& =-p(x) y(x)+f(x),
\end{aligned}
$$

which implies that $y$ is a solution of (2.1.31).
To prove (ii), suppose $y$ is a solution of (2.1.31) on $(a, b)$. From the proof of Theorem 2.1.1, we know that $y_{1}$ has no zeros on $(a, b)$, so the function $u=y / y_{1}$ is defined on $(a, b)$. Moreover, since

$$
\begin{aligned}
y^{\prime} & =-p y+f \quad \text { and } \quad y_{1}^{\prime}=-p y_{1}, \\
u^{\prime} & =\frac{y_{1} y^{\prime}-y_{1}^{\prime} y}{y_{1}^{2}} \\
& =\frac{y_{1}(-p y+f)-\left(-p y_{1}\right) y}{y_{1}^{2}}=\frac{f}{y_{1}} .
\end{aligned}
$$

Integrating $u^{\prime}=f / y_{1}$ yields

$$
u=\left(c+\int f(x) / y_{1}(x) d x\right)
$$

which implies (2.1.32), since $y=u y_{1}$.
(b) We've proved (a), where $\int f(x) / y_{1}(x) d x$ in (2.1.32) is an arbitrary antiderivative of $f / y_{1}$. Now it's convenient to choose the antiderivative that equals zero when $x=x_{0}$, and write the general solution of (2.1.31) as

$$
y=y_{1}(x)\left(c+\int_{x_{0}}^{x} \frac{f(t)}{y_{1}(t)} d t\right)
$$

Since

$$
y\left(x_{0}\right)=y_{1}\left(x_{0}\right)\left(c+\int_{x_{0}}^{x_{0}} \frac{f(t)}{y_{1}(t)} d t\right)=c y_{1}\left(x_{0}\right)
$$

we see that $y\left(x_{0}\right)=y_{0}$ if and only if $c=y_{0} / y_{1}\left(x_{0}\right)$.

### 2.1 Exercises

In Exercises 1-5 find the general solution.

1. $y^{\prime}+a y=0(a=$ constant $)$
2. $y^{\prime}+3 x^{2} y=0$
3. $x y^{\prime}+(\ln x) y=0$
4. $x y^{\prime}+3 y=0$
5. $x^{2} y^{\prime}+y=0$

In Exercises 6-11 solve the initial value problem.
6. $\quad y^{\prime}+\left(\frac{1+x}{x}\right) y=0, \quad y(1)=1$
7. $x y^{\prime}+\left(1+\frac{1}{\ln x}\right) y=0, \quad y(e)=1$
8. $x y^{\prime}+(1+x \cot x) y=0, \quad y\left(\frac{\pi}{2}\right)=2$
9. $y^{\prime}-\left(\frac{2 x}{1+x^{2}}\right) y=0, \quad y(0)=2$
10. $\quad y^{\prime}+\frac{k}{x} y=0, \quad y(1)=3 \quad(k=$ constant $)$
11. $y^{\prime}+(\tan k x) y=0, \quad y(0)=2 \quad(k=$ constant $)$

In Exercises 12-15 find the general solution. Also, plot a direction field and some integral curves on the rectangular region $\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$.
12. $\mathrm{C} / \mathrm{G} y^{\prime}+3 y=1$
13. $\mathrm{C} / \mathrm{G} y^{\prime}+\left(\frac{1}{x}-1\right) y=-\frac{2}{x}$
14. $\mathrm{C} / \mathrm{G} y^{\prime}+2 x y=x e^{-x^{2}}$
15. $\mathrm{C} / \mathrm{G} y^{\prime}+\frac{2 x}{1+x^{2}} y=\frac{e^{-x}}{1+x^{2}}$

In Exercises 16-24 find the general solution.
16. $y^{\prime}+\frac{1}{x} y=\frac{7}{x^{2}}+3$
17. $y^{\prime}+\frac{4}{x-1} y=\frac{1}{(x-1)^{5}}+\frac{\sin x}{(x-1)^{4}}$
18. $x y^{\prime}+\left(1+2 x^{2}\right) y=x^{3} e^{-x^{2}}$
19. $x y^{\prime}+2 y=\frac{2}{x^{2}}+1$
20. $y^{\prime}+(\tan x) y=\cos x$
21. $(1+x) y^{\prime}+2 y=\frac{\sin x}{1+x}$
22. $(x-2)(x-1) y^{\prime}-(4 x-3) y=(x-2)^{3}$
23. $y^{\prime}+(2 \sin x \cos x) y=e^{-\sin ^{2} x}$
24. $x^{2} y^{\prime}+3 x y=e^{x}$

In Exercises 25-29 solve the initial value problem and sketch the graph of the solution.
25. C/G $y^{\prime}+7 y=e^{3 x}, \quad y(0)=0$
26. $\mathrm{C} / \mathrm{G}\left(1+x^{2}\right) y^{\prime}+4 x y=\frac{2}{1+x^{2}}, \quad y(0)=1$
27. $\mathrm{C} / \mathrm{G} x y^{\prime}+3 y=\frac{2}{x\left(1+x^{2}\right)}, \quad y(-1)=0$
28. $\mathrm{C} / \mathrm{G} y^{\prime}+(\cot x) y=\cos x, \quad y\left(\frac{\pi}{2}\right)=1$
29. $\mathrm{C} / \mathrm{G} y^{\prime}+\frac{1}{x} y=\frac{2}{x^{2}}+1, \quad y(-1)=0$

In Exercises 30-37 solve the initial value problem.
30. $(x-1) y^{\prime}+3 y=\frac{1}{(x-1)^{3}}+\frac{\sin x}{(x-1)^{2}}, \quad y(0)=1$
31. $x y^{\prime}+2 y=8 x^{2}, \quad y(1)=3$
32. $x y^{\prime}-2 y=-x^{2}, \quad y(1)=1$
33. $y^{\prime}+2 x y=x, \quad y(0)=3$
34. $(x-1) y^{\prime}+3 y=\frac{1+(x-1) \sec ^{2} x}{(x-1)^{3}}, \quad y(0)=-1$
35. $(x+2) y^{\prime}+4 y=\frac{1+2 x^{2}}{x(x+2)^{3}}, \quad y(-1)=2$
36. $\left(x^{2}-1\right) y^{\prime}-2 x y=x\left(x^{2}-1\right), \quad y(0)=4$
37. $\left(x^{2}-5\right) y^{\prime}-2 x y=-2 x\left(x^{2}-5\right), \quad y(2)=7$

In Exercises 38-42 solve the initial value problem and leave the answer in a form involving a definite integral. (You can solve these problems numerically by methods discussed in Chapter 3.)
38. $y^{\prime}+2 x y=x^{2}, \quad y(0)=3$
39. $y^{\prime}+\frac{1}{x} y=\frac{\sin x}{x^{2}}, \quad y(1)=2$
40. $y^{\prime}+y=\frac{e^{-x} \tan x}{x}, \quad y(1)=0$
41. $y^{\prime}+\frac{2 x}{1+x^{2}} y=\frac{e^{x}}{\left(1+x^{2}\right)^{2}}, \quad y(0)=1$
42. $x y^{\prime}+(x+1) y=e^{x^{2}}, \quad y(1)=2$
43. Experiments indicate that glucose is absorbed by the body at a rate proportional to the amount of glucose present in the bloodstream. Let $\lambda$ denote the (positive) constant of proportionality. Now suppose glucose is injected into a patient's bloodstream at a constant rate of $r$ units per unit of time. Let $G=G(t)$ be the number of units in the patient's bloodstream at time $t>0$. Then

$$
G^{\prime}=-\lambda G+r,
$$

where the first term on the right is due to the absorption of the glucose by the patient's body and the second term is due to the injection. Determine $G$ for $t>0$, given that $G(0)=G_{0}$. Also, find $\lim _{t \rightarrow \infty} G(t)$.
44. (a) L Plot a direction field and some integral curves for

$$
\begin{equation*}
x y^{\prime}-2 y=-1 \tag{A}
\end{equation*}
$$

on the rectangular region $\{-1 \leq x \leq 1,-.5 \leq y \leq 1.5\}$. What do all the integral curves have in common?
(b) Show that the general solution of $(\mathrm{A})$ on $(-\infty, 0)$ and $(0, \infty)$ is

$$
y=\frac{1}{2}+c x^{2} .
$$

(c) Show that $y$ is a solution of $(\mathrm{A})$ on $(-\infty, \infty)$ if and only if

$$
y= \begin{cases}\frac{1}{2}+c_{1} x^{2}, & x \geq 0 \\ \frac{1}{2}+c_{2} x^{2}, & x<0\end{cases}
$$

where $c_{1}$ and $c_{2}$ are arbitrary constants.
(d) Conclude from (c) that all solutions of (A) on $(-\infty, \infty)$ are solutions of the initial value problem

$$
x y^{\prime}-2 y=-1, \quad y(0)=\frac{1}{2}
$$

(e) Use (b) to show that if $x_{0} \neq 0$ and $y_{0}$ is arbitrary, then the initial value problem

$$
x y^{\prime}-2 y=-1, \quad y\left(x_{0}\right)=y_{0}
$$

has infinitely many solutions on $(-\infty, \infty)$. Explain why this does'nt contradict Theorem 2.1.1(b).
45. Suppose $f$ is continuous on an open interval $(a, b)$ and $\alpha$ is a constant.
(a) Derive a formula for the solution of the initial value problem

$$
\begin{equation*}
y^{\prime}+\alpha y=f(x), \quad y\left(x_{0}\right)=y_{0} \tag{A}
\end{equation*}
$$

where $x_{0}$ is in $(a, b)$ and $y_{0}$ is an arbitrary real number.
(b) Suppose $(a, b)=(a, \infty), \alpha>0$ and $\lim _{x \rightarrow \infty} f(x)=L$. Show that if $y$ is the solution of (A), then $\lim _{x \rightarrow \infty} y(x)=L / \alpha$.
46. Assume that all functions in this exercise are defined on a common interval $(a, b)$.
(a) Prove: If $y_{1}$ and $y_{2}$ are solutions of

$$
y^{\prime}+p(x) y=f_{1}(x)
$$

and

$$
y^{\prime}+p(x) y=f_{2}(x)
$$

respectively, and $c_{1}$ and $c_{2}$ are constants, then $y=c_{1} y_{1}+c_{2} y_{2}$ is a solution of

$$
y^{\prime}+p(x) y=c_{1} f_{1}(x)+c_{2} f_{2}(x) .
$$

(This is theprinciple of superposition.)
(b) Use (a) to show that if $y_{1}$ and $y_{2}$ are solutions of the nonhomogeneous equation

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{A}
\end{equation*}
$$

then $y_{1}-y_{2}$ is a solution of the homogeneous equation

$$
\begin{equation*}
y^{\prime}+p(x) y=0 \tag{B}
\end{equation*}
$$

(c) Use (a) to show that if $y_{1}$ is a solution of (A) and $y_{2}$ is a solution of (B), then $y_{1}+y_{2}$ is a solution of (A).
47. Some nonlinear equations can be transformed into linear equations by changing the dependent variable. Show that if

$$
g^{\prime}(y) y^{\prime}+p(x) g(y)=f(x)
$$

where $y$ is a function of $x$ and $g$ is a function of $y$, then the new dependent variable $z=g(y)$ satisfies the linear equation

$$
z^{\prime}+p(x) z=f(x)
$$

48. Solve by the method discussed in Exercise 47.
(a) $\left(\sec ^{2} y\right) y^{\prime}-3 \tan y=-1$
(b) $e^{y^{2}}\left(2 y y^{\prime}+\frac{2}{x}\right)=\frac{1}{x^{2}}$
(c) $\frac{x y^{\prime}}{y}+2 \ln y=4 x^{2}$
(d) $\frac{y^{\prime}}{(1+y)^{2}}-\frac{1}{x(1+y)}=-\frac{3}{x^{2}}$
49. We've shown that if $p$ and $f$ are continuous on $(a, b)$ then every solution of

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{A}
\end{equation*}
$$

on $(a, b)$ can be written as $y=u y_{1}$, where $y_{1}$ is a nontrivial solution of the complementary equation for (A) and $u^{\prime}=f / y_{1}$. Now suppose $f, f^{\prime}, \ldots, f^{(m)}$ and $p, p^{\prime}, \ldots, p^{(m-1)}$ are continuous on ( $a, b$ ), where $m$ is a positive integer, and define

$$
\begin{aligned}
& f_{0}=f \\
& f_{j}=f_{j-1}^{\prime}+p f_{j-1}, \quad 1 \leq j \leq m
\end{aligned}
$$

Show that

$$
u^{(j+1)}=\frac{f_{j}}{y_{1}}, \quad 0 \leq j \leq m
$$

### 2.2 SEPARABLE EQUATIONS

A first order differential equation is separable if it can be written as

$$
\begin{equation*}
h(y) y^{\prime}=g(x), \tag{2.2.1}
\end{equation*}
$$

where the left side is a product of $y^{\prime}$ and a function of $y$ and the right side is a function of $x$. Rewriting a separable differential equation in this form is called separation of variables. In Section 2.1 we used
separation of variables to solve homogeneous linear equations. In this section we'll apply this method to nonlinear equations.
To see how to solve (2.2.1), let's first assume that $y$ is a solution. Let $G(x)$ and $H(y)$ be antiderivatives of $g(x)$ and $h(y)$; that is,

$$
\begin{equation*}
H^{\prime}(y)=h(y) \quad \text { and } \quad G^{\prime}(x)=g(x) . \tag{2.2.2}
\end{equation*}
$$

Then, from the chain rule,

$$
\frac{d}{d x} H(y(x))=H^{\prime}(y(x)) y^{\prime}(x)=h(y) y^{\prime}(x) .
$$

Therefore (2.2.1) is equivalent to

$$
\frac{d}{d x} H(y(x))=\frac{d}{d x} G(x) .
$$

Integrating both sides of this equation and combining the constants of integration yields

$$
\begin{equation*}
H(y(x))=G(x)+c . \tag{2.2.3}
\end{equation*}
$$

Although we derived this equation on the assumption that $y$ is a solution of (2.2.1), we can now view it differently: Any differentiable function $y$ that satisfies (2.2.3) for some constant $c$ is a solution of (2.2.1). To see this, we differentiate both sides of (2.2.3), using the chain rule on the left, to obtain

$$
H^{\prime}(y(x)) y^{\prime}(x)=G^{\prime}(x),
$$

which is equivalent to

$$
h(y(x)) y^{\prime}(x)=g(x)
$$

because of (2.2.2).
In conclusion, to solve (2.2.1) it suffices to find functions $G=G(x)$ and $H=H(y)$ that satisfy (2.2.2). Then any differentiable function $y=y(x)$ that satisfies (2.2.3) is a solution of (2.2.1).

Example 2.2.1 Solve the equation

$$
y^{\prime}=x\left(1+y^{2}\right) .
$$

Solution Separating variables yields

$$
\frac{y^{\prime}}{1+y^{2}}=x
$$

Integrating yields

$$
\tan ^{-1} y=\frac{x^{2}}{2}+c
$$

Therefore

$$
y=\tan \left(\frac{x^{2}}{2}+c\right)
$$

## Example 2.2.2

(a) Solve the equation

$$
\begin{equation*}
y^{\prime}=-\frac{x}{y} . \tag{2.2.4}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime}=-\frac{x}{y}, \quad y(1)=1 \tag{2.2.5}
\end{equation*}
$$

(c) Solve the initial value problem

$$
\begin{equation*}
y^{\prime}=-\frac{x}{y}, \quad y(1)=-2 . \tag{2.2.6}
\end{equation*}
$$

$\underline{\text { SOLUTION(a) }}$ Separating variables in (2.2.4) yields

$$
y y^{\prime}=-x .
$$

Integrating yields

$$
\frac{y^{2}}{2}=-\frac{x^{2}}{2}+c, \quad \text { or, equivalently, } \quad x^{2}+y^{2}=2 c
$$

The last equation shows that $c$ must be positive if $y$ is to be a solution of (2.2.4) on an open interval. Therefore we let $2 c=a^{2}$ (with $a>0$ ) and rewrite the last equation as

$$
\begin{equation*}
x^{2}+y^{2}=a^{2} . \tag{2.2.7}
\end{equation*}
$$

This equation has two differentiable solutions for $y$ in terms of $x$ :

$$
\begin{equation*}
y=\sqrt{a^{2}-x^{2}}, \quad-a<x<a, \tag{2.2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
y=-\sqrt{a^{2}-x^{2}}, \quad-a<x<a . \tag{2.2.9}
\end{equation*}
$$

The solution curves defined by (2.2.8) are semicircles above the $x$-axis and those defined by (2.2.9) are semicircles below the $x$-axis (Figure 2.2.1).

SOLUTION(b) The solution of (2.2.5) is positive when $x=1$; hence, it is of the form (2.2.8). Substituting $x=1$ and $y=1$ into (2.2.7) to satisfy the initial condition yields $a^{2}=2$; hence, the solution of (2.2.5) is

$$
y=\sqrt{2-x^{2}}, \quad-\sqrt{2}<x<\sqrt{2} .
$$

SOLUTION(c) The solution of (2.2.6) is negative when $x=1$ and is therefore of the form (2.2.9). Substituting $x=1$ and $y=-2$ into (2.2.7) to satisfy the initial condition yields $a^{2}=5$. Hence, the solution of (2.2.6) is

$$
y=-\sqrt{5-x^{2}}, \quad-\sqrt{5}<x<\sqrt{5} .
$$



Figure 2.2.1 (a) $y=\sqrt{2-x^{2}},-\sqrt{2}<x<\sqrt{2} ;$ (b) $y=-\sqrt{5-x^{2}},-\sqrt{5}<x<\sqrt{5}$

Implicit Solutions of Separable Equations
In Examples 2.2.1 and 2.2.2 we were able to solve the equation $H(y)=G(x)+c$ to obtain explicit formulas for solutions of the given separable differential equations. As we'll see in the next example,

## 42 Chapter 2 First Order Equations

this isn't always possible. In this situation we must broaden our definition of a solution of a separable equation. The next theorem provides the basis for this modification. We omit the proof, which requires a result from advanced calculus called as the implicit function theorem.

Theorem 2.2.1 Suppose $g=g(x)$ is continous on $(a, b)$ and $h=h(y)$ are continuous on $(c, d)$. Let $G$ be an antiderivative of $g$ on $(a, b)$ and let $H$ be an antiderivative of $h$ on $(c, d)$. Let $x_{0}$ be an arbitrary point in $(a, b)$, let $y_{0}$ be a point in $(c, d)$ such that $h\left(y_{0}\right) \neq 0$, and define

$$
\begin{equation*}
c=H\left(y_{0}\right)-G\left(x_{0}\right) . \tag{2.2.10}
\end{equation*}
$$

Then there's a function $y=y(x)$ defined on some open interval ( $a_{1}, b_{1}$ ), where $a \leq a_{1}<x_{0}<b_{1} \leq b$, such that $y\left(x_{0}\right)=y_{0}$ and

$$
\begin{equation*}
H(y)=G(x)+c \tag{2.2.11}
\end{equation*}
$$

for $a_{1}<x<b_{1}$. Therefore $y$ is a solution of the initial value problem

$$
\begin{equation*}
h(y) y^{\prime}=g(x), \quad y\left(x_{0}\right)=x_{0} . \tag{2.2.12}
\end{equation*}
$$

It's convenient to say that (2.2.11) with $c$ arbitrary is an implicit solution of $h(y) y^{\prime}=g(x)$. Curves defined by (2.2.11) are integral curves of $h(y) y^{\prime}=g(x)$. If $c$ satisfies (2.2.10), we'll say that (2.2.11) is an implicit solution of the initial value problem (2.2.12). However, keep these points in mind:

- For some choices of $c$ there may not be any differentiable functions $y$ that satisfy (2.2.11).
- The function $y$ in (2.2.11) (not (2.2.11) itself) is a solution of $h(y) y^{\prime}=g(x)$.


## Example 2.2.3

(a) Find implicit solutions of

$$
\begin{equation*}
y^{\prime}=\frac{2 x+1}{5 y^{4}+1} . \tag{2.2.13}
\end{equation*}
$$

(b) Find an implicit solution of

$$
\begin{equation*}
y^{\prime}=\frac{2 x+1}{5 y^{4}+1}, \quad y(2)=1 \tag{2.2.14}
\end{equation*}
$$

SOLUTION(a) Separating variables yields

$$
\left(5 y^{4}+1\right) y^{\prime}=2 x+1 .
$$

Integrating yields the implicit solution

$$
\begin{equation*}
y^{5}+y=x^{2}+x+c . \tag{2.2.15}
\end{equation*}
$$

of (2.2.13).
 Therefore

$$
y^{5}+y=x^{2}+x-4
$$

is an implicit solution of the initial value problem (2.2.14). Although more than one differentiable function $y=y(x)$ satisfies 2.2.13) near $x=1$, it can be shown that there's only one such function that satisfies the initial condition $y(1)=2$.

Figure 2.2.2 shows a direction field and some integral curves for (2.2.13).
Constant Solutions of Separable Equations
An equation of the form

$$
y^{\prime}=g(x) p(y)
$$

is separable, since it can be rewritten as

$$
\frac{1}{p(y)} y^{\prime}=g(x)
$$

However, the division by $p(y)$ is not legitimate if $p(y)=0$ for some values of $y$. The next two examples show how to deal with this problem.


Figure 2.2.2 A direction field and integral curves for $y^{\prime}=\frac{2 x+1}{5 y^{4}+1}$

Example 2.2.4 Find all solutions of

$$
\begin{equation*}
y^{\prime}=2 x y^{2} \tag{2.2.16}
\end{equation*}
$$

Solution Here we must divide by $p(y)=y^{2}$ to separate variables. This isn't legitimate if $y$ is a solution of (2.2.16) that equals zero for some value of $x$. One such solution can be found by inspection: $y \equiv 0$. Now suppose $y$ is a solution of (2.2.16) that isn't identically zero. Since $y$ is continuous there must be an interval on which $y$ is never zero. Since division by $y^{2}$ is legitimate for $x$ in this interval, we can separate variables in (2.2.16) to obtain

$$
\frac{y^{\prime}}{y^{2}}=2 x
$$

Integrating this yields

$$
-\frac{1}{y}=x^{2}+c,
$$

which is equivalent to

$$
\begin{equation*}
y=-\frac{1}{x^{2}+c} . \tag{2.2.17}
\end{equation*}
$$

We've now shown that if $y$ is a solution of (2.2.16) that is not identically zero, then $y$ must be of the form (2.2.17). By substituting (2.2.17) into (2.2.16), you can verify that (2.2.17) is a solution of (2.2.16). Thus, solutions of (2.2.16) are $y \equiv 0$ and the functions of the form (2.2.17). Note that the solution $y \equiv 0$ isn't of the form (2.2.17) for any value of $c$.

Figure 2.2.3 shows a direction field and some integral curves for (2.2.16)
Example 2.2.5 Find all solutions of

$$
\begin{equation*}
y^{\prime}=\frac{1}{2} x\left(1-y^{2}\right) . \tag{2.2.18}
\end{equation*}
$$

Solution Here we must divide by $p(y)=1-y^{2}$ to separate variables. This isn't legitimate if $y$ is a solution of (2.2.18) that equals $\pm 1$ for some value of $x$. Two such solutions can be found by inspection: $y \equiv 1$ and $y \equiv-1$. Now suppose $y$ is a solution of (2.2.18) such that $1-y^{2}$ isn't identically zero. Since $1-y^{2}$ is continuous there must be an interval on which $1-y^{2}$ is never zero. Since division by $1-y^{2}$ is legitimate for $x$ in this interval, we can separate variables in (2.2.18) to obtain

$$
\frac{2 y^{\prime}}{y^{2}-1}=-x
$$



Figure 2.2.3 A direction field and integral curves for $y^{\prime}=2 x y^{2}$

A partial fraction expansion on the left yields

$$
\left[\frac{1}{y-1}-\frac{1}{y+1}\right] y^{\prime}=-x
$$

and integrating yields

$$
\ln \left|\frac{y-1}{y+1}\right|=-\frac{x^{2}}{2}+k
$$

hence,

$$
\left|\frac{y-1}{y+1}\right|=e^{k} e^{-x^{2} / 2}
$$

Since $y(x) \neq \pm 1$ for $x$ on the interval under discussion, the quantity $(y-1) /(y+1)$ can't change sign in this interval. Therefore we can rewrite the last equation as

$$
\frac{y-1}{y+1}=c e^{-x^{2} / 2}
$$

where $c= \pm e^{k}$, depending upon the sign of $(y-1) /(y+1)$ on the interval. Solving for $y$ yields

$$
\begin{equation*}
y=\frac{1+c e^{-x^{2} / 2}}{1-c e^{-x^{2} / 2}} \tag{2.2.19}
\end{equation*}
$$

We've now shown that if $y$ is a solution of (2.2.18) that is not identically equal to $\pm 1$, then $y$ must be as in (2.2.19). By substituting (2.2.19) into (2.2.18) you can verify that (2.2.19) is a solution of (2.2.18). Thus, the solutions of (2.2.18) are $y \equiv 1, y \equiv-1$ and the functions of the form (2.2.19). Note that the constant solution $y \equiv 1$ can be obtained from this formula by taking $c=0$; however, the other constant solution, $y \equiv-1$, can't be obtained in this way.

Figure 2.2.4 shows a direction field and some integrals for (2.2.18).


Figure 2.2.4 A direction field and integral curves for $y^{\prime}=\frac{x\left(1-y^{2}\right)}{2}$

Differences Between Linear and Nonlinear Equations
Theorem 2.1.2 states that if $p$ and $f$ are continuous on $(a, b)$ then every solution of

$$
y^{\prime}+p(x) y=f(x)
$$

on $(a, b)$ can be obtained by choosing a value for the constant $c$ in the general solution, and if $x_{0}$ is any point in ( $a, b$ ) and $y_{0}$ is arbitrary, then the initial value problem

$$
y^{\prime}+p(x) y=f(x), \quad y\left(x_{0}\right)=y_{0}
$$

has a solution on $(a, b)$.
The not true for nonlinear equations. First, we saw in Examples 2.2.4 and 2.2.5 that a nonlinear equation may have solutions that can't be obtained by choosing a specific value of a constant appearing in a one-parameter family of solutions. Second, it is in general impossible to determine the interval of validity of a solution to an initial value problem for a nonlinear equation by simply examining the equation, since the interval of validity may depend on the initial condition. For instance, in Example 2.2.2 we saw that the solution of

$$
\frac{d y}{d x}=-\frac{x}{y}, \quad y\left(x_{0}\right)=y_{0}
$$

is valid on $(-a, a)$, where $a=\sqrt{x_{0}^{2}+y_{0}^{2}}$.
Example 2.2.6 Solve the initial value problem

$$
y^{\prime}=2 x y^{2}, \quad y(0)=y_{0}
$$

and determine the interval of validity of the solution.

Solution First suppose $y_{0} \neq 0$. From Example 2.2.4, we know that $y$ must be of the form

$$
\begin{equation*}
y=-\frac{1}{x^{2}+c} . \tag{2.2.20}
\end{equation*}
$$

Imposing the initial condition shows that $c=-1 / y_{0}$. Substituting this into (2.2.20) and rearranging terms yields the solution

$$
y=\frac{y_{0}}{1-y_{0} x^{2}} .
$$

This is also the solution if $y_{0}=0$. If $y_{0}<0$, the denominator isn't zero for any value of $x$, so the the solution is valid on $(-\infty, \infty)$. If $y_{0}>0$, the solution is valid only on $\left(-1 / \sqrt{y_{0}}, 1 / \sqrt{y_{0}}\right)$.

### 2.2 Exercises

In Exercises 1-6 find all solutions.

1. $y^{\prime}=\frac{3 x^{2}+2 x+1}{y-2}$
2. $x y^{\prime}+y^{2}+y=0$
3. $\left(3 y^{3}+3 y \cos y+1\right) y^{\prime}+\frac{(2 x+1) y}{1+x^{2}}=0$
4. $x^{2} y y^{\prime}=\left(y^{2}-1\right)^{3 / 2}$

In Exercises 7-10 find all solutions. Also, plot a direction field and some integral curves on the indicated rectangular region.
7. C/G $y^{\prime}=x^{2}\left(1+y^{2}\right) ;\{-1 \leq x \leq 1,-1 \leq y \leq 1\}$
8. $\mathrm{C} / \mathrm{G} y^{\prime}\left(1+x^{2}\right)+x y=0 ;\{-2 \leq x \leq 2,-1 \leq y \leq 1\}$
9. $\mathrm{C} / \mathrm{G} y^{\prime}=(x-1)(y-1)(y-2) ;\{-2 \leq x \leq 2,-3 \leq y \leq 3\}$
10. C/G $(y-1)^{2} y^{\prime}=2 x+3 ;\{-2 \leq x \leq 2,-2 \leq y \leq 5\}$

In Exercises 11 and 12 solve the initial value problem.
11. $y^{\prime}=\frac{x^{2}+3 x+2}{y-2}, \quad y(1)=4$
12. $y^{\prime}+x\left(y^{2}+y\right)=0, \quad y(2)=1$

In Exercises 13-16 solve the initial value problem and graph the solution.
13. $\mathrm{C} / \mathrm{G}\left(3 y^{2}+4 y\right) y^{\prime}+2 x+\cos x=0, \quad y(0)=1$
14. $\mathrm{C} / \mathrm{G} y^{\prime}+\frac{(y+1)(y-1)(y-2)}{x+1}=0, \quad y(1)=0$
15. $\mathrm{C} / \mathrm{G} y^{\prime}+2 x(y+1)=0, \quad y(0)=2$
16. $\mathrm{C} / \mathrm{G} y^{\prime}=2 x y\left(1+y^{2}\right), \quad y(0)=1$

In Exercises 17-23 solve the initial value problem and find the interval of validity of the solution.
17. $y^{\prime}\left(x^{2}+2\right)+4 x\left(y^{2}+2 y+1\right)=0, \quad y(1)=-1$
18. $y^{\prime}=-2 x\left(y^{2}-3 y+2\right), \quad y(0)=3$
19. $y^{\prime}=\frac{2 x}{1+2 y}, \quad y(2)=0$
20. $y^{\prime}=2 y-y^{2}, \quad y(0)=1$
21. $x+y y^{\prime}=0, \quad y(3)=-4$
22. $y^{\prime}+x^{2}(y+1)(y-2)^{2}=0, \quad y(4)=2$
23. $(x+1)(x-2) y^{\prime}+y=0, \quad y(1)=-3$
24. Solve $y^{\prime}=\frac{\left(1+y^{2}\right)}{\left(1+x^{2}\right)}$ explicitly. HinT: Use the identity $\tan (A+B)=\frac{\tan A+\tan B}{1-\tan A \tan B}$.
25. Solve $y^{\prime} \sqrt{1-x^{2}}+\sqrt{1-y^{2}}=0$ explicitly. HinT: Use the identity $\sin (A-B)=\sin A \cos B-$ $\cos A \sin B$.
26. Solve $y^{\prime}=\frac{\cos x}{\sin y}, \quad y(\pi)=\frac{\pi}{2}$ explicitly. Hint: Use the identity $\cos (x+\pi / 2)=-\sin x$ and the periodicity of the cosine.
27. Solve the initial value problem

$$
y^{\prime}=a y-b y^{2}, \quad y(0)=y_{0} .
$$

Discuss the behavior of the solution if (a) $y_{0} \geq 0$; (b) $y_{0}<0$.
28. The population $P=P(t)$ of a species satisfies the logistic equation

$$
P^{\prime}=a P(1-\alpha P)
$$

and $P(0)=P_{0}>0$. Find $P$ for $t>0$, and find $\lim _{t \rightarrow \infty} P(t)$.
29. An epidemic spreads through a population at a rate proportional to the product of the number of people already infected and the number of people susceptible, but not yet infected. Therefore, if $S$ denotes the total population of susceptible people and $I=I(t)$ denotes the number of infected people at time $t$, then

$$
I^{\prime}=r I(S-I)
$$

where $r$ is a positive constant. Assuming that $I(0)=I_{0}$, find $I(t)$ for $t>0$, and show that $\lim _{t \rightarrow \infty} I(t)=S$.
30. L The result of Exercise 29 is discouraging: if any susceptible member of the group is initially infected, then in the long run all susceptible members are infected! On a more hopeful note, suppose the disease spreads according to the model of Exercise 29, but there's a medication that cures the infected population at a rate proportional to the number of infected individuals. Now the equation for the number of infected individuals becomes

$$
\begin{equation*}
I^{\prime}=r I(S-I)-q I \tag{A}
\end{equation*}
$$

where $q$ is a positive constant.
(a) Choose $r$ and $S$ positive. By plotting direction fields and solutions of (A) on suitable rectangular grids

$$
R=\{0 \leq t \leq T, 0 \leq I \leq d\}
$$

in the $(t, I)$-plane, verify that if $I$ is any solution of (A) such that $I(0)>0$, then $\lim _{t \rightarrow \infty} I(t)=$ $S-q / r$ if $q<r S$ and $\lim _{t \rightarrow \infty} I(t)=0$ if $q \geq r S$.
(b) To verify the experimental results of (a), use separation of variables to solve (A) with initial condition $I(0)=I_{0}>0$, and find $\lim _{t \rightarrow \infty} I(t)$. Hint: There are three cases to consider: (i) $q<r S$; (ii) $q>r S$; (iii) $q=r S$.
31. L Consider the differential equation

$$
\begin{equation*}
y^{\prime}=a y-b y^{2}-q, \tag{A}
\end{equation*}
$$

where $a, b$ are positive constants, and $q$ is an arbitrary constant. Suppose $y$ denotes a solution of this equation that satisfies the initial condition $y(0)=y_{0}$.
(a) Choose $a$ and $b$ positive and $q<a^{2} / 4 b$. By plotting direction fields and solutions of (A) on suitable rectangular grids

$$
\begin{equation*}
R=\{0 \leq t \leq T, c \leq y \leq d\} \tag{B}
\end{equation*}
$$

in the $(t, y)$-plane, discover that there are numbers $y_{1}$ and $y_{2}$ with $y_{1}<y_{2}$ such that if $y_{0}>y_{1}$ then $\lim _{t \rightarrow \infty} y(t)=y_{2}$, and if $y_{0}<y_{1}$ then $y(t)=-\infty$ for some finite value of $t$. (What happens if $y_{0}=y_{1}$ ?)
(b) Choose $a$ and $b$ positive and $q=a^{2} / 4 b$. By plotting direction fields and solutions of (A) on suitable rectangular grids of the form (B), discover that there's a number $y_{1}$ such that if $y_{0} \geq y_{1}$ then $\lim _{t \rightarrow \infty} y(t)=y_{1}$, while if $y_{0}<y_{1}$ then $y(t)=-\infty$ for some finite value of $t$.
(c) Choose positive $a, b$ and $q>a^{2} / 4 b$. By plotting direction fields and solutions of (A) on suitable rectangular grids of the form (B), discover that no matter what $y_{0}$ is, $y(t)=-\infty$ for some finite value of $t$.
(d) Verify your results experiments analytically. Start by separating variables in (A) to obtain

$$
\frac{y^{\prime}}{a y-b y^{2}-q}=1
$$

To decide what to do next you'll have to use the quadratic formula. This should lead you to see why there are three cases. Take it from there!
Because of its role in the transition between these three cases, $q_{0}=a^{2} / 4 b$ is called a bifurcation value of $q$. In general, if $q$ is a parameter in any differential equation, $q_{0}$ is said to be a bifurcation value of $q$ if the nature of the solutions of the equation with $q<q_{0}$ is qualitatively different from the nature of the solutions with $q>q_{0}$.
32. $L$ By plotting direction fields and solutions of

$$
y^{\prime}=q y-y^{3}
$$

convince yourself that $q_{0}=0$ is a bifurcation value of $q$ for this equation. Explain what makes you draw this conclusion.
33. Suppose a disease spreads according to the model of Exercise 29, but there's a medication that cures the infected population at a constant rate of $q$ individuals per unit time, where $q>0$. Then the equation for the number of infected individuals becomes

$$
I^{\prime}=r I(S-I)-q
$$

Assuming that $I(0)=I_{0}>0$, use the results of Exercise 31 to describe what happens as $t \rightarrow \infty$.
34. Assuming that $p \not \equiv 0$, state conditions under which the linear equation

$$
y^{\prime}+p(x) y=f(x)
$$

is separable. If the equation satisfies these conditions, solve it by separation of variables and by the method developed in Section 2.1.

Solve the equations in Exercises 35-38 using variation of parameters followed by separation of variables.
35. $y^{\prime}+y=\frac{2 x e^{-x}}{1+y e^{x}}$
36. $x y^{\prime}-2 y=\frac{x^{6}}{y+x^{2}}$
37. $y^{\prime}-y=\frac{(x+1) e^{4 x}}{\left(y+e^{x}\right)^{2}}$
38. $y^{\prime}-2 y=\frac{x e^{2 x}}{1-y e^{-2 x}}$
39. Use variation of parameters to show that the solutions of the following equations are of the form $y=u y_{1}$, where $u$ satisfies a separable equation $u^{\prime}=g(x) p(u)$. Find $y_{1}$ and $g$ for each equation.
(a) $x y^{\prime}+y=h(x) p(x y)$
(b) $x y^{\prime}-y=h(x) p\left(\frac{y}{x}\right)$
(c) $y^{\prime}+y=h(x) p\left(e^{x} y\right)$
(d) $x y^{\prime}+r y=h(x) p\left(x^{r} y\right)$
(e) $y^{\prime}+\frac{v^{\prime}(x)}{v(x)} y=h(x) p(v(x) y)$

### 2.3 EXISTENCE AND UNIQUENESS OF SOLUTIONS OF NONLINEAR EQUATIONS

Although there are methods for solving some nonlinear equations, it's impossible to find useful formulas for the solutions of most. Whether we're looking for exact solutions or numerical approximations, it's useful to know conditions that imply the existence and uniqueness of solutions of initial value problems
for nonlinear equations. In this section we state such a condition and illustrate it with examples.


Figure 2.3.1 An open rectangle

Some terminology: an open rectangle $R$ is a set of points $(x, y)$ such that

$$
a<x<b \quad \text { and } \quad c<y<d
$$

(Figure 2.3.1). We'll denote this set by $R:\{a<x<b, c<y<d\}$. "Open" means that the boundary rectangle (indicated by the dashed lines in Figure 2.3.1) isn't included in $R$.

The next theorem gives sufficient conditions for existence and uniqueness of solutions of initial value problems for first order nonlinear differential equations. We omit the proof, which is beyond the scope of this book.

## Theorem 2.3.1

(a) If $f$ is continuous on an open rectangle

$$
R:\{a<x<b, c<y<d\}
$$

that contains $\left(x_{0}, y_{0}\right)$ then the initial value problem

$$
\begin{equation*}
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0} \tag{2.3.1}
\end{equation*}
$$

has at least one solution on some open subinterval of $(a, b)$ that contains $x_{0}$.
(b) If both $f$ and $f_{y}$ are continuous on $R$ then (2.3.1) has a unique solution on some open subinterval of $(a, b)$ that contains $x_{0}$.

It's important to understand exactly what Theorem 2.3.1 says.

- (a) is an existence theorem. It guarantees that a solution exists on some open interval that contains $x_{0}$, but provides no information on how to find the solution, or to determine the open interval on which it exists. Moreover, (a) provides no information on the number of solutions that (2.3.1) may have. It leaves open the possibility that (2.3.1) may have two or more solutions that differ for values of $x$ arbitrarily close to $x_{0}$. We will see in Example 2.3.6 that this can happen.
- (b) is a uniqueness theorem. It guarantees that (2.3.1) has a unique solution on some open interval $(\mathrm{a}, \mathrm{b})$ that contains $x_{0}$. However, if $(a, b) \neq(-\infty, \infty)$, (2.3.1) may have more than one solution on a larger interval that contains $(a, b)$. For example, it may happen that $b<\infty$ and all solutions have the same values on $(a, b)$, but two solutions $y_{1}$ and $y_{2}$ are defined on some interval $\left(a, b_{1}\right)$ with $b_{1}>b$, and have different values for $b<x<b_{1}$; thus, the graphs of the $y_{1}$ and $y_{2}$ "branch off" in different directions at $x=b$. (See Example 2.3.7 and Figure 2.3.3). In this case, continuity implies that $y_{1}(b)=y_{2}(b)$ (call their common value $\bar{y}$ ), and $y_{1}$ and $y_{2}$ are both solutions of the initial value problem

$$
\begin{equation*}
y^{\prime}=f(x, y), \quad y(b)=\bar{y} \tag{2.3.2}
\end{equation*}
$$

that differ on every open interval that contains $b$. Therefore $f$ or $f_{y}$ must have a discontinuity at some point in each open rectangle that contains $(b, \bar{y})$, since if this were not so, (2.3.2) would have a unique solution on some open interval that contains $b$. We leave it to you to give a similar analysis of the case where $a>-\infty$.

Example 2.3.1 Consider the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{x^{2}-y^{2}}{1+x^{2}+y^{2}}, \quad y\left(x_{0}\right)=y_{0} \tag{2.3.3}
\end{equation*}
$$

Since

$$
f(x, y)=\frac{x^{2}-y^{2}}{1+x^{2}+y^{2}} \quad \text { and } \quad f_{y}(x, y)=-\frac{2 y\left(1+2 x^{2}\right)}{\left(1+x^{2}+y^{2}\right)^{2}}
$$

are continuous for all $(x, y)$, Theorem 2.3.1 implies that if $\left(x_{0}, y_{0}\right)$ is arbitrary, then (2.3.3) has a unique solution on some open interval that contains $x_{0}$.

Example 2.3.2 Consider the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{x^{2}-y^{2}}{x^{2}+y^{2}}, \quad y\left(x_{0}\right)=y_{0} \tag{2.3.4}
\end{equation*}
$$

Here

$$
f(x, y)=\frac{x^{2}-y^{2}}{x^{2}+y^{2}} \quad \text { and } \quad f_{y}(x, y)=-\frac{4 x^{2} y}{\left(x^{2}+y^{2}\right)^{2}}
$$

are continuous everywhere except at $(0,0)$. If $\left(x_{0}, y_{0}\right) \neq(0,0)$, there's an open rectangle $R$ that contains $\left(x_{0}, y_{0}\right)$ that does not contain $(0,0)$. Since $f$ and $f_{y}$ are continuous on $R$, Theorem 2.3.1 implies that if $\left(x_{0}, y_{0}\right) \neq(0,0)$ then (2.3.4) has a unique solution on some open interval that contains $x_{0}$.

Example 2.3.3 Consider the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{x+y}{x-y}, \quad y\left(x_{0}\right)=y_{0} . \tag{2.3.5}
\end{equation*}
$$

Here

$$
f(x, y)=\frac{x+y}{x-y} \quad \text { and } \quad f_{y}(x, y)=\frac{2 x}{(x-y)^{2}}
$$

are continuous everywhere except on the line $y=x$. If $y_{0} \neq x_{0}$, there's an open rectangle $R$ that contains $\left(x_{0}, y_{0}\right)$ that does not intersect the line $y=x$. Since $f$ and $f_{y}$ are continuous on $R$, Theorem 2.3.1 implies that if $y_{0} \neq x_{0}$, (2.3.5) has a unique solution on some open interval that contains $x_{0}$.

Example 2.3.4 In Example 2.2.4 we saw that the solutions of

$$
\begin{equation*}
y^{\prime}=2 x y^{2} \tag{2.3.6}
\end{equation*}
$$

are

$$
y \equiv 0 \quad \text { and } \quad y=-\frac{1}{x^{2}+c}
$$

where $c$ is an arbitrary constant. In particular, this implies that no solution of (2.3.6) other than $y \equiv 0$ can equal zero for any value of $x$. Show that Theorem 2.3.1(b) implies this.

Solution We'll obtain a contradiction by assuming that (2.3.6) has a solution $y_{1}$ that equals zero for some value of $x$, but isn't identically zero. If $y_{1}$ has this property, there's a point $x_{0}$ such that $y_{1}\left(x_{0}\right)=0$, but $y_{1}(x) \neq 0$ for some value of $x$ in every open interval that contains $x_{0}$. This means that the initial value problem

$$
\begin{equation*}
y^{\prime}=2 x y^{2}, \quad y\left(x_{0}\right)=0 \tag{2.3.7}
\end{equation*}
$$

has two solutions $y \equiv 0$ and $y=y_{1}$ that differ for some value of $x$ on every open interval that contains $x_{0}$. This contradicts Theorem 2.3.1(b), since in (2.3.6) the functions

$$
f(x, y)=2 x y^{2} \quad \text { and } \quad f_{y}(x, y)=4 x y
$$

are both continuous for all $(x, y)$, which implies that (2.3.7) has a unique solution on some open interval that contains $x_{0}$.

Example 2.3.5 Consider the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{10}{3} x y^{2 / 5}, \quad y\left(x_{0}\right)=y_{0} . \tag{2.3.8}
\end{equation*}
$$

(a) For what points $\left(x_{0}, y_{0}\right)$ does Theorem 2.3.1(a) imply that (2.3.8) has a solution?
(b) For what points $\left(x_{0}, y_{0}\right)$ does Theorem 2.3.1(b) imply that (2.3.8) has a unique solution on some open interval that contains $x_{0}$ ?
$\underline{\text { SOLUTION(a) Since }}$

$$
f(x, y)=\frac{10}{3} x y^{2 / 5}
$$

is continuous for all $(x, y)$, Theorem 2.3.1 implies that (2.3.8) has a solution for every $\left(x_{0}, y_{0}\right)$.
$\underline{\text { SOLUTION(b) }}$ Here

$$
f_{y}(x, y)=\frac{4}{3} x y^{-3 / 5}
$$

is continuous for all $(x, y)$ with $y \neq 0$. Therefore, if $y_{0} \neq 0$ there's an open rectangle on which both $f$ and $f_{y}$ are continuous, and Theorem 2.3.1 implies that (2.3.8) has a unique solution on some open interval that contains $x_{0}$.

If $y=0$ then $f_{y}(x, y)$ is undefined, and therefore discontinuous; hence, Theorem 2.3.1 does not apply to (2.3.8) if $y_{0}=0$.

Example 2.3.6 Example 2.3.5 leaves open the possibility that the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{10}{3} x y^{2 / 5}, \quad y(0)=0 \tag{2.3.9}
\end{equation*}
$$

has more than one solution on every open interval that contains $x_{0}=0$. Show that this is true.

Solution By inspection, $y \equiv 0$ is a solution of the differential equation

$$
\begin{equation*}
y^{\prime}=\frac{10}{3} x y^{2 / 5} \tag{2.3.10}
\end{equation*}
$$

Since $y \equiv 0$ satisfies the initial condition $y(0)=0$, it's a solution of (2.3.9).
Now suppose $y$ is a solution of (2.3.10) that isn't identically zero. Separating variables in (2.3.10) yields

$$
y^{-2 / 5} y^{\prime}=\frac{10}{3} x
$$

on any open interval where $y$ has no zeros. Integrating this and rewriting the arbitrary constant as $5 c / 3$ yields

$$
\frac{5}{3} y^{3 / 5}=\frac{5}{3}\left(x^{2}+c\right)
$$

Therefore

$$
\begin{equation*}
y=\left(x^{2}+c\right)^{5 / 3} \tag{2.3.11}
\end{equation*}
$$

Since we divided by $y$ to separate variables in (2.3.10), our derivation of (2.3.11) is legitimate only on open intervals where $y$ has no zeros. However, (2.3.11) actually defines $y$ for all $x$, and differentiating (2.3.11) shows that

$$
y^{\prime}=\frac{10}{3} x\left(x^{2}+c\right)^{2 / 3}=\frac{10}{3} x y^{2 / 5},-\infty<x<\infty .
$$

Therefore (2.3.11) satisfies (2.3.10) on $(-\infty, \infty)$ even if $c \leq 0$, so that $y(\sqrt{|c|})=y(-\sqrt{|c|})=0$. In particular, taking $c=0$ in (2.3.11) yields

$$
y=x^{10 / 3}
$$

as a second solution of (2.3.9). Both solutions are defined on $(-\infty, \infty)$, and they differ on every open interval that contains $x_{0}=0$ (see Figure 2.3.2.) In fact, there are four distinct solutions of (2.3.9) defined on $(-\infty, \infty)$ that differ from each other on every open interval that contains $x_{0}=0$. Can you identify the other two?


Figure 2.3.2 Two solutions ( $y=0$ and $y=x^{1 / 2}$ ) of (2.3.9) that differ on every interval containing $x_{0}=0$

Example 2.3.7 From Example 2.3.5, the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{10}{3} x y^{2 / 5}, \quad y(0)=-1 \tag{2.3.12}
\end{equation*}
$$

has a unique solution on some open interval that contains $x_{0}=0$. Find a solution and determine the largest open interval $(a, b)$ on which it's unique.

Solution Let $y$ be any solution of (2.3.12). Because of the initial condition $y(0)=-1$ and the continuity of $y$, there's an open interval $I$ that contains $x_{0}=0$ on which $y$ has no zeros, and is consequently of the form (2.3.11). Setting $x=0$ and $y=-1$ in (2.3.11) yields $c=-1$, so

$$
\begin{equation*}
y=\left(x^{2}-1\right)^{5 / 3} \tag{2.3.13}
\end{equation*}
$$

for $x$ in $I$. Therefore every solution of (2.3.12) differs from zero and is given by (2.3.13) on $(-1,1)$; that is, (2.3.13) is the unique solution of (2.3.12) on $(-1,1)$. This is the largest open interval on which (2.3.12) has a unique solution. To see this, note that (2.3.13) is a solution of (2.3.12) on $(-\infty, \infty)$. From Exercise 2.2.15, there are infinitely many other solutions of (2.3.12) that differ from (2.3.13) on every open interval larger than $(-1,1)$. One such solution is

$$
y=\left\{\begin{array}{cl}
\left(x^{2}-1\right)^{5 / 3}, & -1 \leq x \leq 1 \\
0, & |x|>1
\end{array}\right.
$$

(Figure 2.3.3).
Example 2.3.8 From Example 2.3.5, the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{10}{3} x y^{2 / 5}, \quad y(0)=1 \tag{2.3.14}
\end{equation*}
$$

has a unique solution on some open interval that contains $x_{0}=0$. Find the solution and determine the largest open interval on which it's unique.

Solution Let $y$ be any solution of (2.3.14). Because of the initial condition $y(0)=1$ and the continuity of $y$, there's an open interval $I$ that contains $x_{0}=0$ on which $y$ has no zeros, and is consequently of the form (2.3.11). Setting $x=0$ and $y=1$ in (2.3.11) yields $c=1$, so

$$
\begin{equation*}
y=\left(x^{2}+1\right)^{5 / 3} \tag{2.3.15}
\end{equation*}
$$



Figure 2.3.3 Two solutions of (2.3.12) on $(-\infty, \infty)$ that coincide on $(-1,1)$, but on no larger open interval


Figure 2.3.4 The unique solution of (2.3.14)
for $x$ in $I$. Therefore every solution of (2.3.14) differs from zero and is given by (2.3.15) on $(-\infty, \infty)$; that is, (2.3.15) is the unique solution of (2.3.14) on $(-\infty, \infty)$. Figure 2.3 . shows the graph of this solution.

### 2.3 Exercises

In Exercises 1-13 find all $\left(x_{0}, y_{0}\right)$ for which Theorem 2.3.1 implies that the initial value problem $y^{\prime}=$ $f(x, y), y\left(x_{0}\right)=y_{0}$ has (a) a solution (b) a unique solution on some open interval that contains $x_{0}$.

1. $y^{\prime}=\frac{x^{2}+y^{2}}{\sin x}$
2. $y^{\prime}=\frac{e^{x}+y}{x^{2}+y^{2}}$
3. $y^{\prime}=\tan x y$
4. $y^{\prime}=\frac{x^{2}+y^{2}}{\ln x y}$
5. $y^{\prime}=\left(x^{2}+y^{2}\right) y^{1 / 3}$
6. $y^{\prime}=2 x y$
7. $y^{\prime}=\ln \left(1+x^{2}+y^{2}\right)$
8. $y^{\prime}=\frac{2 x+3 y}{x-4 y}$
9. $y^{\prime}=\left(x^{2}+y^{2}\right)^{1 / 2}$
10. $y^{\prime}=x\left(y^{2}-1\right)^{2 / 3}$
11. $y^{\prime}=\left(x^{2}+y^{2}\right)^{2}$
12. $y^{\prime}=(x+y)^{1 / 2}$
13. $y^{\prime}=\frac{\tan y}{x-1}$
14. Apply Theorem 2.3.1 to the initial value problem

$$
y^{\prime}+p(x) y=q(x), \quad y\left(x_{0}\right)=y_{0}
$$

for a linear equation, and compare the conclusions that can be drawn from it to those that follow from Theorem 2.1.2.
15. (a) Verify that the function

$$
y=\left\{\begin{array}{cl}
\left(x^{2}-1\right)^{5 / 3}, & -1<x<1 \\
0, & |x| \geq 1
\end{array}\right.
$$

is a solution of the initial value problem

$$
y^{\prime}=\frac{10}{3} x y^{2 / 5}, \quad y(0)=-1
$$

on $(-\infty, \infty)$. Hint: You'll need the definition

$$
y^{\prime}(\bar{x})=\lim _{x \rightarrow \bar{x}} \frac{y(x)-y(\bar{x})}{x-\bar{x}}
$$

to verify that $y$ satisfies the differential equation at $\bar{x}= \pm 1$.
(b) Verify that if $\epsilon_{i}=0$ or 1 for $i=1,2$ and $a, b>1$, then the function

$$
y=\left\{\begin{array}{cl}
\epsilon_{1}\left(x^{2}-a^{2}\right)^{5 / 3}, & -\infty<x<-a, \\
0, & -a \leq x \leq-1, \\
\left(x^{2}-1\right)^{5 / 3}, & -1<x<1, \\
0, & 1 \leq x \leq b, \\
\epsilon_{2}\left(x^{2}-b^{2}\right)^{5 / 3}, & b<x<\infty,
\end{array}\right.
$$

is a solution of the initial value problem of (a) on $(-\infty, \infty)$.
16. Use the ideas developed in Exercise 15 to find infinitely many solutions of the initial value problem

$$
y^{\prime}=y^{2 / 5}, \quad y(0)=1
$$

on $(-\infty, \infty)$.
17. Consider the initial value problem

$$
\begin{equation*}
y^{\prime}=3 x(y-1)^{1 / 3}, \quad y\left(x_{0}\right)=y_{0} \tag{A}
\end{equation*}
$$

(a) For what points $\left(x_{0}, y_{0}\right)$ does Theorem 2.3.1 imply that (A) has a solution?
(b) For what points $\left(x_{0}, y_{0}\right)$ does Theorem 2.3.1 imply that (A) has a unique solution on some open interval that contains $x_{0}$ ?
18. Find nine solutions of the initial value problem

$$
y^{\prime}=3 x(y-1)^{1 / 3}, \quad y(0)=1
$$

that are all defined on $(-\infty, \infty)$ and differ from each other for values of $x$ in every open interval that contains $x_{0}=0$.
19. From Theorem 2.3.1, the initial value problem

$$
y^{\prime}=3 x(y-1)^{1 / 3}, \quad y(0)=9
$$

has a unique solution on an open interval that contains $x_{0}=0$. Find the solution and determine the largest open interval on which it's unique.
20. (a) From Theorem 2.3.1, the initial value problem

$$
\begin{equation*}
y^{\prime}=3 x(y-1)^{1 / 3}, \quad y(3)=-7 \tag{A}
\end{equation*}
$$

has a unique solution on some open interval that contains $x_{0}=3$. Determine the largest such open interval, and find the solution on this interval.
(b) Find infinitely many solutions of (A), all defined on $(-\infty, \infty)$.
21. Prove:
(a) If

$$
\begin{equation*}
f\left(x, y_{0}\right)=0, \quad a<x<b, \tag{A}
\end{equation*}
$$

and $x_{0}$ is in $(a, b)$, then $y \equiv y_{0}$ is a solution of

$$
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0}
$$

on $(a, b)$.
(b) If $f$ and $f_{y}$ are continuous on an open rectangle that contains $\left(x_{0}, y_{0}\right)$ and (A) holds, no solution of $y^{\prime}=f(x, y)$ other than $y \equiv y_{0}$ can equal $y_{0}$ at any point in $(a, b)$.

### 2.5 EXACT EQUATIONS

In this section it's convenient to write first order differential equations in the form

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 \tag{2.5.1}
\end{equation*}
$$

This equation can be interpreted as

$$
\begin{equation*}
M(x, y)+N(x, y) \frac{d y}{d x}=0 \tag{2.5.2}
\end{equation*}
$$

where $x$ is the independent variable and $y$ is the dependent variable, or as

$$
\begin{equation*}
M(x, y) \frac{d x}{d y}+N(x, y)=0 \tag{2.5.3}
\end{equation*}
$$

where $y$ is the independent variable and $x$ is the dependent variable. Since the solutions of (2.5.2) and (2.5.3) will often have to be left in implicit, form we'll say that $F(x, y)=c$ is an implicit solution of (2.5.1) if every differentiable function $y=y(x)$ that satisfies $F(x, y)=c$ is a solution of (2.5.2) and every differentiable function $x=x(y)$ that satisfies $F(x, y)=c$ is a solution of (2.5.3).

Here are some examples:

| Equation (2.5.1) | Equation (2.5.2) | Equation (2.5.3) |
| :---: | :---: | :---: |
| $3 x^{2} y^{2} d x+2 x^{3} y d y=0$ | $3 x^{2} y^{2}+2 x^{3} y \frac{d y}{d x}=0$ | $3 x^{2} y^{2} \frac{d x}{d y}+2 x^{3} y=0$ |
| $\left(x^{2}+y^{2}\right) d x+2 x y d y=0$ | $\left(x^{2}+y^{2}\right)+2 x y \frac{d y}{d x}=0$ | $\left(x^{2}+y^{2}\right) \frac{d x}{d y}+2 x y=0$ |
| $3 y \sin x d x-2 x y \cos x d y=0$ | $3 y \sin x-2 x y \cos x \frac{d y}{d x}=0$ | $3 y \sin x \frac{d x}{d y}-2 x y \cos x=0$ |

Note that a separable equation can be written as (2.5.1) as

$$
M(x) d x+N(y) d y=0
$$

We'll develop a method for solving (2.5.1) under appropriate assumptions on $M$ and $N$. This method is an extension of the method of separation of variables (Exercise 41). Before stating it we consider an example.

Example 2.5.1 Show that

$$
\begin{equation*}
x^{4} y^{3}+x^{2} y^{5}+2 x y=c \tag{2.5.4}
\end{equation*}
$$

is an implicit solution of

$$
\begin{equation*}
\left(4 x^{3} y^{3}+2 x y^{5}+2 y\right) d x+\left(3 x^{4} y^{2}+5 x^{2} y^{4}+2 x\right) d y=0 \tag{2.5.5}
\end{equation*}
$$

Solution Regarding $y$ as a function of $x$ and differentiating (2.5.4) implicitly with respect to $x$ yields

$$
\left(4 x^{3} y^{3}+2 x y^{5}+2 y\right)+\left(3 x^{4} y^{2}+5 x^{2} y^{4}+2 x\right) \frac{d y}{d x}=0
$$

Similarly, regarding $x$ as a function of $y$ and differentiating (2.5.4) implicitly with respect to $y$ yields

$$
\left(4 x^{3} y^{3}+2 x y^{5}+2 y\right) \frac{d x}{d y}+\left(3 x^{4} y^{2}+5 x^{2} y^{4}+2 x\right)=0
$$

Therefore (2.5.4) is an implicit solution of (2.5.5) in either of its two possible interpretations.
You may think this example is pointless, since concocting a differential equation that has a given implicit solution isn't particularly interesting. However, it illustrates the next important theorem, which we'll prove by using implicit differentiation, as in Example 2.5.1.

## 56 Chapter 2 First Order Equations

Theorem 2.5.1 If $F=F(x, y)$ has continuous partial derivatives $F_{x}$ and $F_{y}$, then

$$
\begin{equation*}
F(x, y)=c \quad(c=\text { constant }), \tag{2.5.6}
\end{equation*}
$$

is an implicit solution of the differential equation

$$
\begin{equation*}
F_{x}(x, y) d x+F_{y}(x, y) d y=0 \tag{2.5.7}
\end{equation*}
$$

Proof Regarding $y$ as a function of $x$ and differentiating (2.5.6) implicitly with respect to $x$ yields

$$
F_{x}(x, y)+F_{y}(x, y) \frac{d y}{d x}=0
$$

On the other hand, regarding $x$ as a function of $y$ and differentiating (2.5.6) implicitly with respect to $y$ yields

$$
F_{x}(x, y) \frac{d x}{d y}+F_{y}(x, y)=0
$$

Thus, (2.5.6) is an implicit solution of (2.5.7) in either of its two possible interpretations.
We'll say that the equation

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 \tag{2.5.8}
\end{equation*}
$$

is exact on an an open rectangle $R$ if there's a function $F=F(x, y)$ such $F_{x}$ and $F_{y}$ are continuous, and

$$
\begin{equation*}
F_{x}(x, y)=M(x, y) \quad \text { and } \quad F_{y}(x, y)=N(x, y) \tag{2.5.9}
\end{equation*}
$$

for all $(x, y)$ in $R$. This usage of "exact" is related to its usage in calculus, where the expression

$$
F_{x}(x, y) d x+F_{y}(x, y) d y
$$

(obtained by substituting (2.5.9) into the left side of (2.5.8)) is the exact differential of $F$.
Example 2.5.1 shows that it's easy to solve (2.5.8) if it's exact and we know a function $F$ that satisfies (2.5.9). The important questions are:

Question 1. Given an equation (2.5.8), how can we determine whether it's exact?
Question 2. If (2.5.8) is exact, how do we find a function $F$ satisfying (2.5.9)?
To discover the answer to Question 1, assume that there's a function $F$ that satisfies (2.5.9) on some open rectangle $R$, and in addition that $F$ has continuous mixed partial derivatives $F_{x y}$ and $F_{y x}$. Then a theorem from calculus implies that

$$
\begin{equation*}
F_{x y}=F_{y x} . \tag{2.5.10}
\end{equation*}
$$

If $F_{x}=M$ and $F_{y}=N$, differentiating the first of these equations with respect to $y$ and the second with respect to $x$ yields

$$
\begin{equation*}
F_{x y}=M_{y} \quad \text { and } \quad F_{y x}=N_{x} . \tag{2.5.11}
\end{equation*}
$$

From (2.5.10) and (2.5.11), we conclude that a necessary condition for exactness is that $M_{y}=N_{x}$. This motivates the next theorem, which we state without proof.

Theorem 2.5.2 [The Exactness Condition] Suppose $M$ and $N$ are continuous and have continuous partial derivatives $M_{y}$ and $N_{x}$ on an open rectangle $R$. Then

$$
M(x, y) d x+N(x, y) d y=0
$$

is exact on $R$ if and only if

$$
\begin{equation*}
M_{y}(x, y)=N_{x}(x, y) \tag{2.5.12}
\end{equation*}
$$

for all $(x, y)$ in $R$..
To help you remember the exactness condition, observe that the coefficients of $d x$ and $d y$ are differentiated in (2.5.12) with respect to the "opposite" variables; that is, the coefficient of $d x$ is differentiated with respect to $y$, while the coefficient of $d y$ is differentiated with respect to $x$.

Example 2.5.2 Show that the equation

$$
3 x^{2} y d x+4 x^{3} d y=0
$$

is not exact on any open rectangle.

## Solution Here

$$
M(x, y)=3 x^{2} y \quad \text { and } \quad N(x, y)=4 x^{3}
$$

so

$$
M_{y}(x, y)=3 x^{2} \quad \text { and } \quad N_{x}(x, y)=12 x^{2} .
$$

Therefore $M_{y}=N_{x}$ on the line $x=0$, but not on any open rectangle, so there's no function $F$ such that $F_{x}(x, y)=M(x, y)$ and $F_{y}(x, y)=N(x, y)$ for all $(x, y)$ on any open rectangle.

The next example illustrates two possible methods for finding a function $F$ that satisfies the condition $F_{x}=M$ and $F_{y}=N$ if $M d x+N d y=0$ is exact.

Example 2.5.3 Solve

$$
\begin{equation*}
\left(4 x^{3} y^{3}+3 x^{2}\right) d x+\left(3 x^{4} y^{2}+6 y^{2}\right) d y=0 \tag{2.5.13}
\end{equation*}
$$

## Solution (Method 1) Here

$$
M(x, y)=4 x^{3} y^{3}+3 x^{2}, \quad N(x, y)=3 x^{4} y^{2}+6 y^{2}
$$

and

$$
M_{y}(x, y)=N_{x}(x, y)=12 x^{3} y^{2}
$$

for all $(x, y)$. Therefore Theorem 2.5.2 implies that there's a function $F$ such that

$$
\begin{equation*}
F_{x}(x, y)=M(x, y)=4 x^{3} y^{3}+3 x^{2} \tag{2.5.14}
\end{equation*}
$$

and

$$
\begin{equation*}
F_{y}(x, y)=N(x, y)=3 x^{4} y^{2}+6 y^{2} \tag{2.5.15}
\end{equation*}
$$

for all $(x, y)$. To find $F$, we integrate (2.5.14) with respect to $x$ to obtain

$$
\begin{equation*}
F(x, y)=x^{4} y^{3}+x^{3}+\phi(y) \tag{2.5.16}
\end{equation*}
$$

where $\phi(y)$ is the "constant" of integration. (Here $\phi$ is "constant" in that it's independent of $x$, the variable of integration.) If $\phi$ is any differentiable function of $y$ then $F$ satisfies (2.5.14). To determine $\phi$ so that $F$ also satisfies (2.5.15), assume that $\phi$ is differentiable and differentiate $F$ with respect to $y$. This yields

$$
F_{y}(x, y)=3 x^{4} y^{2}+\phi^{\prime}(y) .
$$

Comparing this with (2.5.15) shows that

$$
\phi^{\prime}(y)=6 y^{2} .
$$

We integrate this with respect to $y$ and take the constant of integration to be zero because we're interested only in finding some $F$ that satisfies (2.5.14) and (2.5.15). This yields

$$
\phi(y)=2 y^{3} .
$$

Substituting this into (2.5.16) yields

$$
\begin{equation*}
F(x, y)=x^{4} y^{3}+x^{3}+2 y^{3} . \tag{2.5.17}
\end{equation*}
$$

Now Theorem 2.5.1 implies that

$$
x^{4} y^{3}+x^{3}+2 y^{3}=c
$$

is an implicit solution of (2.5.13). Solving this for $y$ yields the explicit solution

$$
y=\left(\frac{c-x^{3}}{2+x^{4}}\right)^{1 / 3}
$$

Solution (Method 2) Instead of first integrating (2.5.14) with respect to $x$, we could begin by integrating (2.5.15) with respect to $y$ to obtain

$$
\begin{equation*}
F(x, y)=x^{4} y^{3}+2 y^{3}+\psi(x) \tag{2.5.18}
\end{equation*}
$$



Figure 2.5.1 A direction field and integral curves for $\left(4 x^{3} y^{3}+3 x^{2}\right) d x+\left(3 x^{4} y^{2}+6 y^{2}\right) d y=0$
where $\psi$ is an arbitrary function of $x$. To determine $\psi$, we assume that $\psi$ is differentiable and differentiate $F$ with respect to $x$, which yields

$$
F_{x}(x, y)=4 x^{3} y^{3}+\psi^{\prime}(x)
$$

Comparing this with (2.5.14) shows that

$$
\psi^{\prime}(x)=3 x^{2} .
$$

Integrating this and again taking the constant of integration to be zero yields

$$
\psi(x)=x^{3} .
$$

Substituting this into (2.5.18) yields (2.5.17).
Figure 2.5.1 shows a direction field and some integral curves of (2.5.13),
Here's a summary of the procedure used in Method 1 of this example. You should summarize procedure used in Method 2.

## Procedure For Solving An Exact Equation

Step 1. Check that the equation

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 \tag{2.5.19}
\end{equation*}
$$

satisfies the exactness condition $M_{y}=N_{x}$. If not, don't go further with this procedure.
Step 2. Integrate

$$
\frac{\partial F(x, y)}{\partial x}=M(x, y)
$$

with respect to $x$ to obtain

$$
\begin{equation*}
F(x, y)=G(x, y)+\phi(y), \tag{2.5.20}
\end{equation*}
$$

where $G$ is an antiderivative of $M$ with respect to $x$, and $\phi$ is an unknown function of $y$.
Step 3. Differentiate (2.5.20) with respect to $y$ to obtain

$$
\frac{\partial F(x, y)}{\partial y}=\frac{\partial G(x, y)}{\partial y}+\phi^{\prime}(y)
$$

Step 4. Equate the right side of this equation to $N$ and solve for $\phi^{\prime}$; thus,

$$
\frac{\partial G(x, y)}{\partial y}+\phi^{\prime}(y)=N(x, y), \quad \text { so } \quad \phi^{\prime}(y)=N(x, y)-\frac{\partial G(x, y)}{\partial y}
$$

Step 5. Integrate $\phi^{\prime}$ with respect to $y$, taking the constant of integration to be zero, and substitute the result in (2.5.20) to obtain $F(x, y)$.

Step 6. Set $F(x, y)=c$ to obtain an implicit solution of (2.5.19). If possible, solve for $y$ explicitly as a function of $x$.

It's a common mistake to omit Step 6. However, it's important to include this step, since $F$ isn't itself a solution of (2.5.19).

Many equations can be conveniently solved by either of the two methods used in Example 2.5.3. However, sometimes the integration required in one approach is more difficult than in the other. In such cases we choose the approach that requires the easier integration.

Example 2.5.4 Solve the equation

$$
\begin{equation*}
\left(y e^{x y} \tan x+e^{x y} \sec ^{2} x\right) d x+x e^{x y} \tan x d y=0 \tag{2.5.21}
\end{equation*}
$$

Solution We leave it to you to check that $M_{y}=N_{x}$ on any open rectangle where $\tan x$ and $\sec x$ are defined. Here we must find a function $F$ such that

$$
\begin{equation*}
F_{x}(x, y)=y e^{x y} \tan x+e^{x y} \sec ^{2} x \tag{2.5.22}
\end{equation*}
$$

and

$$
\begin{equation*}
F_{y}(x, y)=x e^{x y} \tan x \tag{2.5.23}
\end{equation*}
$$

It's difficult to integrate (2.5.22) with respect to $x$, but easy to integrate (2.5.23) with respect to $y$. This yields

$$
\begin{equation*}
F(x, y)=e^{x y} \tan x+\psi(x) \tag{2.5.24}
\end{equation*}
$$

Differentiating this with respect to $x$ yields

$$
F_{x}(x, y)=y e^{x y} \tan x+e^{x y} \sec ^{2} x+\psi^{\prime}(x)
$$

Comparing this with (2.5.22) shows that $\psi^{\prime}(x)=0$. Hence, $\psi$ is a constant, which we can take to be zero in (2.5.24), and

$$
e^{x y} \tan x=c
$$

is an implicit solution of (2.5.21).
Attempting to apply our procedure to an equation that isn't exact will lead to failure in Step 4, since the function

$$
N-\frac{\partial G}{\partial y}
$$

won't be independent of $x$ if $M_{y} \neq N_{x}$ (Exercise 31), and therefore can't be the derivative of a function of $y$ alone. Here's an example that illustrates this.

Example 2.5.5 Verify that the equation

$$
\begin{equation*}
3 x^{2} y^{2} d x+6 x^{3} y d y=0 \tag{2.5.25}
\end{equation*}
$$

is not exact, and show that the procedure for solving exact equations fails when applied to (2.5.25).

## Solution Here

$$
M_{y}(x, y)=6 x^{2} y \quad \text { and } \quad N_{x}(x, y)=18 x^{2} y
$$

so (2.5.25) isn't exact. Nevertheless, let's try to find a function $F$ such that

$$
\begin{equation*}
F_{x}(x, y)=3 x^{2} y^{2} \tag{2.5.26}
\end{equation*}
$$

and

$$
\begin{equation*}
F_{y}(x, y)=6 x^{3} y \tag{2.5.27}
\end{equation*}
$$

Integrating (2.5.26) with respect to $x$ yields

$$
F(x, y)=x^{3} y^{2}+\phi(y)
$$

and differentiating this with respect to $y$ yields

$$
F_{y}(x, y)=2 x^{3} y+\phi^{\prime}(y)
$$

For this equation to be consistent with (2.5.27),

$$
6 x^{3} y=2 x^{3} y+\phi^{\prime}(y)
$$

or

$$
\phi^{\prime}(y)=4 x^{3} y .
$$

This is a contradiction, since $\phi^{\prime}$ must be independent of $x$. Therefore the procedure fails.

### 2.5 Exercises

In Exercises 1-17 determine which equations are exact and solve them.

1. $6 x^{2} y^{2} d x+4 x^{3} y d y=0$
2. $\left(3 y \cos x+4 x e^{x}+2 x^{2} e^{x}\right) d x+(3 \sin x+3) d y=0$
3. $14 x^{2} y^{3} d x+21 x^{2} y^{2} d y=0$
4. $\left(2 x-2 y^{2}\right) d x+\left(12 y^{2}-4 x y\right) d y=0$
5. $(x+y)^{2} d x+(x+y)^{2} d y=0$
6. $(4 x+7 y) d x+(3 x+4 y) d y=0$
7. $\left(-2 y^{2} \sin x+3 y^{3}-2 x\right) d x+\left(4 y \cos x+9 x y^{2}\right) d y=0$
8. $(2 x+y) d x+(2 y+2 x) d y=0$
9. $\left(3 x^{2}+2 x y+4 y^{2}\right) d x+\left(x^{2}+8 x y+18 y\right) d y=0$
10. $\left(2 x^{2}+8 x y+y^{2}\right) d x+\left(2 x^{2}+x y^{3} / 3\right) d y=0$
11. $\left(\frac{1}{x}+2 x\right) d x+\left(\frac{1}{y}+2 y\right) d y=0$
12. $\left(y \sin x y+x y^{2} \cos x y\right) d x+\left(x \sin x y+x y^{2} \cos x y\right) d y=0$
13. $\frac{x d x}{\left(x^{2}+y^{2}\right)^{3 / 2}}+\frac{y d y}{\left(x^{2}+y^{2}\right)^{3 / 2}}=0$
14. $\left(e^{x}\left(x^{2} y^{2}+2 x y^{2}\right)+6 x\right) d x+\left(2 x^{2} y e^{x}+2\right) d y=0$
15. $\left(x^{2} e^{x^{2}+y}\left(2 x^{2}+3\right)+4 x\right) d x+\left(x^{3} e^{x^{2}+y}-12 y^{2}\right) d y=0$
16. $\left(e^{x y}\left(x^{4} y+4 x^{3}\right)+3 y\right) d x+\left(x^{5} e^{x y}+3 x\right) d y=0$
17. $\left(3 x^{2} \cos x y-x^{3} y \sin x y+4 x\right) d x+\left(8 y-x^{4} \sin x y\right) d y=0$

In Exercises 18-22 solve the initial value problem.
18. $\left(4 x^{3} y^{2}-6 x^{2} y-2 x-3\right) d x+\left(2 x^{4} y-2 x^{3}\right) d y=0, \quad y(1)=3$
19. $\left(-4 y \cos x+4 \sin x \cos x+\sec ^{2} x\right) d x+(4 y-4 \sin x) d y=0, \quad y(\pi / 4)=0$
20. $\left(y^{3}-1\right) e^{x} d x+3 y^{2}\left(e^{x}+1\right) d y=0, \quad y(0)=0$
21. $\quad(\sin x-y \sin x-2 \cos x) d x+\cos x d y=0, \quad y(0)=1$
22. $(2 x-1)(y-1) d x+(x+2)(x-3) d y=0, \quad y(1)=-1$
23. $\mathrm{C} / \mathrm{G}$ Solve the exact equation

$$
(7 x+4 y) d x+(4 x+3 y) d y=0
$$

Plot a direction field and some integral curves for this equation on the rectangle

$$
\{-1 \leq x \leq 1,-1 \leq y \leq 1\} .
$$

24. C/G Solve the exact equation

$$
e^{x}\left(x^{4} y^{2}+4 x^{3} y^{2}+1\right) d x+\left(2 x^{4} y e^{x}+2 y\right) d y=0
$$

Plot a direction field and some integral curves for this equation on the rectangle

$$
\{-2 \leq x \leq 2,-1 \leq y \leq 1\}
$$

25. C/G Plot a direction field and some integral curves for the exact equation

$$
\left(x^{3} y^{4}+x\right) d x+\left(x^{4} y^{3}+y\right) d y=0
$$

on the rectangle $\{-1 \leq x \leq 1,-1 \leq y \leq 1\}$. (See Exercise 37(a)).
26. $\mathrm{C} / \mathrm{G}$ Plot a direction field and some integral curves for the exact equation

$$
\left(3 x^{2}+2 y\right) d x+(2 y+2 x) d y=0
$$

on the rectangle $\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$. (See Exercise 37(b)).
27. L
(a) Solve the exact equation

$$
\begin{equation*}
\left(x^{3} y^{4}+2 x\right) d x+\left(x^{4} y^{3}+3 y\right) d y=0 \tag{A}
\end{equation*}
$$

implicitly.
(b) For what choices of $\left(x_{0}, y_{0}\right)$ does Theorem 2.3.1 imply that the initial value problem

$$
\begin{equation*}
\left(x^{3} y^{4}+2 x\right) d x+\left(x^{4} y^{3}+3 y\right) d y=0, \quad y\left(x_{0}\right)=y_{0} \tag{B}
\end{equation*}
$$

has a unique solution on an open interval $(a, b)$ that contains $x_{0}$ ?
(c) Plot a direction field and some integral curves for (A) on a rectangular region centered at the origin. What is the interval of validity of the solution of (B)?
28. L
(a) Solve the exact equation

$$
\begin{equation*}
\left(x^{2}+y^{2}\right) d x+2 x y d y=0 \tag{A}
\end{equation*}
$$

implicitly.
(b) For what choices of $\left(x_{0}, y_{0}\right)$ does Theorem 2.3.1 imply that the initial value problem

$$
\begin{equation*}
\left(x^{2}+y^{2}\right) d x+2 x y d y=0, \quad y\left(x_{0}\right)=y_{0} \tag{B}
\end{equation*}
$$

has a unique solution $y=y(x)$ on some open interval $(a, b)$ that contains $x_{0}$ ?
(c) Plot a direction field and some integral curves for (A). From the plot determine, the interval $(a, b)$ of (b), the monotonicity properties (if any) of the solution of (B), and $\lim _{x \rightarrow a+} y(x)$ and $\lim _{x \rightarrow b-} y(x)$. Hint: Your answers will depend upon which quadrant contains $\left(x_{0}, y_{0}\right)$.
29. Find all functions $M$ such that the equation is exact.
(a) $M(x, y) d x+\left(x^{2}-y^{2}\right) d y=0$
(b) $M(x, y) d x+2 x y \sin x \cos y d y=0$
(c) $M(x, y) d x+\left(e^{x}-e^{y} \sin x\right) d y=0$
30. Find all functions $N$ such that the equation is exact.
(a) $\left(x^{3} y^{2}+2 x y+3 y^{2}\right) d x+N(x, y) d y=0$
(b) $(\ln x y+2 y \sin x) d x+N(x, y) d y=0$
(c) $(x \sin x+y \sin y) d x+N(x, y) d y=0$
31. Suppose $M, N$, and their partial derivatives are continuous on an open rectangle $R$, and $G$ is an antiderivative of $M$ with respect to $x$; that is,

$$
\frac{\partial G}{\partial x}=M
$$

Show that if $M_{y} \neq N_{x}$ in $R$ then the function

$$
N-\frac{\partial G}{\partial y}
$$

is not independent of $x$.
32. Prove: If the equations $M_{1} d x+N_{1} d y=0$ and $M_{2} d x+N_{2} d y=0$ are exact on an open rectangle $R$, so is the equation

$$
\left(M_{1}+M_{2}\right) d x+\left(N_{1}+N_{2}\right) d y=0 .
$$

33. Find conditions on the constants $A, B, C$, and $D$ such that the equation

$$
(A x+B y) d x+(C x+D y) d y=0
$$

is exact.
34. Find conditions on the constants $A, B, C, D, E$, and $F$ such that the equation

$$
\left(A x^{2}+B x y+C y^{2}\right) d x+\left(D x^{2}+E x y+F y^{2}\right) d y=0
$$

is exact.
35. Suppose $M$ and $N$ are continuous and have continuous partial derivatives $M_{y}$ and $N_{x}$ that satisfy the exactness condition $M_{y}=N_{x}$ on an open rectangle $R$. Show that if $(x, y)$ is in $R$ and

$$
F(x, y)=\int_{x_{0}}^{x} M\left(s, y_{0}\right) d s+\int_{y_{0}}^{y} N(x, t) d t,
$$

then $F_{x}=M$ and $F_{y}=N$.
36. Under the assumptions of Exercise 35, show that

$$
F(x, y)=\int_{y_{0}}^{y} N\left(x_{0}, s\right) d s+\int_{x_{0}}^{x} M(t, y) d t .
$$

37. Use the method suggested by Exercise 35, with $\left(x_{0}, y_{0}\right)=(0,0)$, to solve the these exact equations:
(a) $\left(x^{3} y^{4}+x\right) d x+\left(x^{4} y^{3}+y\right) d y=0$
(b) $\left(x^{2}+y^{2}\right) d x+2 x y d y=0$
(c) $\left(3 x^{2}+2 y\right) d x+(2 y+2 x) d y=0$
38. Solve the initial value problem

$$
y^{\prime}+\frac{2}{x} y=-\frac{2 x y}{x^{2}+2 x^{2} y+1}, \quad y(1)=-2 .
$$

39. Solve the initial value problem

$$
y^{\prime}-\frac{3}{x} y=\frac{2 x^{4}\left(4 x^{3}-3 y\right)}{3 x^{5}+3 x^{3}+2 y}, \quad y(1)=1 .
$$

40. Solve the initial value problem

$$
y^{\prime}+2 x y=-e^{-x^{2}}\left(\frac{3 x+2 y e^{x^{2}}}{2 x+3 y e^{x^{2}}}\right), \quad y(0)=-1
$$

41. Rewrite the separable equation

$$
\begin{equation*}
h(y) y^{\prime}=g(x) \tag{A}
\end{equation*}
$$

as an exact equation

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 \tag{B}
\end{equation*}
$$

Show that applying the method of this section to (B) yields the same solutions that would be obtained by applying the method of separation of variables to (A)
42. Suppose all second partial derivatives of $M=M(x, y)$ and $N=N(x, y)$ are continuous and $M d x+N d y=0$ and $-N d x+M d y=0$ are exact on an open rectangle $R$. Show that $M_{x x}+M_{y y}=N_{x x}+N_{y y}=0$ on $R$.
43. Suppose all second partial derivatives of $F=F(x, y)$ are continuous and $F_{x x}+F_{y y}=0$ on an open rectangle $R$. (A function with these properties is said to be harmonic; see also Exercise 42.) Show that $-F_{y} d x+F_{x} d y=0$ is exact on $R$, and therefore there's a function $G$ such that $G_{x}=-F_{y}$ and $G_{y}=F_{x}$ in $R$. (A function $G$ with this property is said to be a harmonic conjugate of $F$.)
44. Verify that the following functions are harmonic, and find all their harmonic conjugates. (See Exercise 43.)
(a) $x^{2}-y^{2}$
(b) $e^{x} \cos y$
(c) $x^{3}-3 x y^{2}$
(d) $\cos x \cosh y$
(e) $\sin x \cosh y$

### 2.6 INTEGRATING FACTORS

In Section 2.5 we saw that if $M, N, M_{y}$ and $N_{x}$ are continuous and $M_{y}=N_{x}$ on an open rectangle $R$ then

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 \tag{2.6.1}
\end{equation*}
$$

is exact on $R$. Sometimes an equation that isn't exact can be made exact by multiplying it by an appropriate function. For example,

$$
\begin{equation*}
\left(3 x+2 y^{2}\right) d x+2 x y d y=0 \tag{2.6.2}
\end{equation*}
$$

is not exact, since $M_{y}(x, y)=4 y \neq N_{x}(x, y)=2 y$ in (2.6.2). However, multiplying (2.6.2) by $x$ yields

$$
\begin{equation*}
\left(3 x^{2}+2 x y^{2}\right) d x+2 x^{2} y d y=0 \tag{2.6.3}
\end{equation*}
$$

which is exact, since $M_{y}(x, y)=N_{x}(x, y)=4 x y$ in (2.6.3). Solving (2.6.3) by the procedure given in Section 2.5 yields the implicit solution

$$
x^{3}+x^{2} y^{2}=c
$$

A function $\mu=\mu(x, y)$ is an integrating factor for (2.6.1) if

$$
\begin{equation*}
\mu(x, y) M(x, y) d x+\mu(x, y) N(x, y) d y=0 \tag{2.6.4}
\end{equation*}
$$

is exact. If we know an integrating factor $\mu$ for (2.6.1), we can solve the exact equation (2.6.4) by the method of Section 2.5. It would be nice if we could say that (2.6.1) and (2.6.4) always have the same solutions, but this isn't so. For example, a solution $y=y(x)$ of (2.6.4) such that $\mu(x, y(x))=0$ on some interval $a<x<b$ could fail to be a solution of (2.6.1) (Exercise 1), while (2.6.1) may have a solution $y=y(x)$ such that $\mu(x, y(x))$ isn't even defined (Exercise 2). Similar comments apply if $y$ is the independent variable and $x$ is the dependent variable in (2.6.1) and (2.6.4). However, if $\mu(x, y)$ is defined and nonzero for all $(x, y),(2.6 .1)$ and (2.6.4) are equivalent; that is, they have the same solutions.
Finding Integrating Factors
By applying Theorem 2.5.2 (with $M$ and $N$ replaced by $\mu M$ and $\mu N$ ), we see that (2.6.4) is exact on an open rectangle $R$ if $\mu M, \mu N,(\mu M)_{y}$, and $(\mu N)_{x}$ are continuous and

$$
\frac{\partial}{\partial y}(\mu M)=\frac{\partial}{\partial x}(\mu N) \quad \text { or, equivalently, } \quad \mu_{y} M+\mu M_{y}=\mu_{x} N+\mu N_{x}
$$

on $R$. It's better to rewrite the last equation as

$$
\begin{equation*}
\mu\left(M_{y}-N_{x}\right)=\mu_{x} N-\mu_{y} M, \tag{2.6.5}
\end{equation*}
$$

which reduces to the known result for exact equations; that is, if $M_{y}=N_{x}$ then (2.6.5) holds with $\mu=1$, so (2.6.1) is exact.

You may think (2.6.5) is of little value, since it involves partial derivatives of the unknown integrating factor $\mu$, and we haven't studied methods for solving such equations. However, we'll now show that (2.6.5) is useful if we restrict our search to integrating factors that are products of a function of $x$ and a function of $y$; that is, $\mu(x, y)=P(x) Q(y)$. We're not saying that every equation $M d x+N d y=0$ has an integrating factor of this form; rather, we're saying that some equations have such integrating factors. We'll now develop a way to determine whether a given equation has such an integrating factor, and a method for finding the integrating factor in this case.

If $\mu(x, y)=P(x) Q(y)$, then $\mu_{x}(x, y)=P^{\prime}(x) Q(y)$ and $\mu_{y}(x, y)=P(x) Q^{\prime}(y)$, so (2.6.5) becomes

$$
\begin{equation*}
P(x) Q(y)\left(M_{y}-N_{x}\right)=P^{\prime}(x) Q(y) N-P(x) Q^{\prime}(y) M \tag{2.6.6}
\end{equation*}
$$

or, after dividing through by $P(x) Q(y)$,

$$
\begin{equation*}
M_{y}-N_{x}=\frac{P^{\prime}(x)}{P(x)} N-\frac{Q^{\prime}(y)}{Q(y)} M \tag{2.6.7}
\end{equation*}
$$

Now let

$$
p(x)=\frac{P^{\prime}(x)}{P(x)} \quad \text { and } \quad q(y)=\frac{Q^{\prime}(y)}{Q(y)}
$$

so (2.6.7) becomes

$$
\begin{equation*}
M_{y}-N_{x}=p(x) N-q(y) M \tag{2.6.8}
\end{equation*}
$$

We obtained (2.6.8) by assuming that $M d x+N d y=0$ has an integrating factor $\mu(x, y)=P(x) Q(y)$. However, we can now view (2.6.7) differently: If there are functions $p=p(x)$ and $q=q(y)$ that satisfy (2.6.8) and we define

$$
\begin{equation*}
P(x)= \pm e^{\int p(x) d x} \quad \text { and } \quad Q(y)= \pm e^{\int q(y) d y} \tag{2.6.9}
\end{equation*}
$$

then reversing the steps that led from (2.6.6) to (2.6.8) shows that $\mu(x, y)=P(x) Q(y)$ is an integrating factor for $M d x+N d y=0$. In using this result, we take the constants of integration in (2.6.9) to be zero and choose the signs conveniently so the integrating factor has the simplest form.

There's no simple general method for ascertaining whether functions $p=p(x)$ and $q=q(y)$ satisfying (2.6.8) exist. However, the next theorem gives simple sufficient conditions for the given equation to have an integrating factor that depends on only one of the independent variables $x$ and $y$, and for finding an integrating factor in this case.

Theorem 2.6.1 Let $M, N, M_{y}$, and $N_{x}$ be continuous on an open rectangle $R$. Then: (a) If $\left(M_{y}-N_{x}\right) / N$ is independent of $y$ on $R$ and we define

$$
p(x)=\frac{M_{y}-N_{x}}{N}
$$

then

$$
\begin{equation*}
\mu(x)= \pm e^{\int p(x) d x} \tag{2.6.10}
\end{equation*}
$$

is an integrating factor for

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 \tag{2.6.11}
\end{equation*}
$$

on $R$.
(b) If $\left(N_{x}-M_{y}\right) / M$ is independent of $x$ on $R$ and we define

$$
q(y)=\frac{N_{x}-M_{y}}{M}
$$

then

$$
\begin{equation*}
\mu(y)= \pm e^{\int q(y) d y} \tag{2.6.12}
\end{equation*}
$$

is an integrating factor for (2.6.11) on $R$.

Proof (a) If $\left(M_{y}-N_{x}\right) / N$ is independent of $y$, then (2.6.8) holds with $p=\left(M_{y}-N_{x}\right) / N$ and $q \equiv 0$. Therefore

$$
P(x)= \pm e^{\int p(x) d x} \quad \text { and } \quad Q(y)= \pm e^{\int q(y) d y}= \pm e^{0}= \pm 1
$$

so (2.6.10) is an integrating factor for $(2.6 .11)$ on $R$.
(b) If $\left(N_{x}-M_{y}\right) / M$ is independent of $x$ then eqrefeq:2.6.8 holds with $p \equiv 0$ and $q=\left(N_{x}-M_{y}\right) / M$, and a similar argument shows that (2.6.12) is an integrating factor for (2.6.11) on $R$.

The next two examples show how to apply Theorem 2.6.1.
Example 2.6.1 Find an integrating factor for the equation

$$
\begin{equation*}
\left(2 x y^{3}-2 x^{3} y^{3}-4 x y^{2}+2 x\right) d x+\left(3 x^{2} y^{2}+4 y\right) d y=0 \tag{2.6.13}
\end{equation*}
$$

and solve the equation.

Solution In (2.6.13)

$$
M=2 x y^{3}-2 x^{3} y^{3}-4 x y^{2}+2 x, N=3 x^{2} y^{2}+4 y
$$

and

$$
M_{y}-N_{x}=\left(6 x y^{2}-6 x^{3} y^{2}-8 x y\right)-6 x y^{2}=-6 x^{3} y^{2}-8 x y
$$

so (2.6.13) isn't exact. However,

$$
\frac{M_{y}-N_{x}}{N}=-\frac{6 x^{3} y^{2}+8 x y}{3 x^{2} y^{2}+4 y}=-2 x
$$

is independent of $y$, so Theorem 2.6.1(a) applies with $p(x)=-2 x$. Since

$$
\int p(x) d x=-\int 2 x d x=-x^{2}
$$

$\mu(x)=e^{-x^{2}}$ is an integrating factor. Multiplying (2.6.13) by $\mu$ yields the exact equation

$$
\begin{equation*}
e^{-x^{2}}\left(2 x y^{3}-2 x^{3} y^{3}-4 x y^{2}+2 x\right) d x+e^{-x^{2}}\left(3 x^{2} y^{2}+4 y\right) d y=0 . \tag{2.6.14}
\end{equation*}
$$

To solve this equation, we must find a function $F$ such that

$$
\begin{equation*}
F_{x}(x, y)=e^{-x^{2}}\left(2 x y^{3}-2 x^{3} y^{3}-4 x y^{2}+2 x\right) \tag{2.6.15}
\end{equation*}
$$

and

$$
\begin{equation*}
F_{y}(x, y)=e^{-x^{2}}\left(3 x^{2} y^{2}+4 y\right) \tag{2.6.16}
\end{equation*}
$$

Integrating (2.6.16) with respect to $y$ yields

$$
\begin{equation*}
F(x, y)=e^{-x^{2}}\left(x^{2} y^{3}+2 y^{2}\right)+\psi(x) . \tag{2.6.17}
\end{equation*}
$$

Differentiating this with respect to $x$ yields

$$
F_{x}(x, y)=e^{-x^{2}}\left(2 x y^{3}-2 x^{3} y^{3}-4 x y^{2}\right)+\psi^{\prime}(x) .
$$

Comparing this with (2.6.15) shows that $\psi^{\prime}(x)=2 x e^{-x^{2}}$; therefore, we can let $\psi(x)=-e^{-x^{2}}$ in (2.6.17) and conclude that

$$
e^{-x^{2}}\left(y^{2}\left(x^{2} y+2\right)-1\right)=c
$$

is an implicit solution of (2.6.14). It is also an implicit solution of (2.6.13).
Figure 2.6.1 shows a direction field and some integal curves for (2.6.13)
Example 2.6.2 Find an integrating factor for

$$
\begin{equation*}
2 x y^{3} d x+\left(3 x^{2} y^{2}+x^{2} y^{3}+1\right) d y=0 \tag{2.6.18}
\end{equation*}
$$

and solve the equation.


Figure 2.6.1 A direction field and integral curves for $\left(2 x y^{3}-2 x^{3} y^{3}-4 x y^{2}+2 x\right) d x+\left(3 x^{2} y^{2}+4 y\right) d y=0$

Solution In (2.6.18),

$$
M=2 x y^{3}, \quad N=3 x^{2} y^{2}+x^{2} y^{3}+1,
$$

and

$$
M_{y}-N_{x}=6 x y^{2}-\left(6 x y^{2}+2 x y^{3}\right)=-2 x y^{3}
$$

so (2.6.18) isn't exact. Moreover,

$$
\frac{M_{y}-N_{x}}{N}=-\frac{2 x y^{3}}{3 x^{2} y^{2}+x^{2} y^{2}+1}
$$

is not independent of $y$, so Theorem 2.6.1(a) does not apply. However, Theorem 2.6.1(b) does apply, since

$$
\frac{N_{x}-M_{y}}{M}=\frac{2 x y^{3}}{2 x y^{3}}=1
$$

is independent of $x$, so we can take $q(y)=1$. Since

$$
\int q(y) d y=\int d y=y
$$

$\mu(y)=e^{y}$ is an integrating factor. Multiplying (2.6.18) by $\mu$ yields the exact equation

$$
\begin{equation*}
2 x y^{3} e^{y} d x+\left(3 x^{2} y^{2}+x^{2} y^{3}+1\right) e^{y} d y=0 \tag{2.6.19}
\end{equation*}
$$

To solve this equation, we must find a function $F$ such that

$$
\begin{equation*}
F_{x}(x, y)=2 x y^{3} e^{y} \tag{2.6.20}
\end{equation*}
$$

and

$$
\begin{equation*}
F_{y}(x, y)=\left(3 x^{2} y^{2}+x^{2} y^{3}+1\right) e^{y} \tag{2.6.21}
\end{equation*}
$$

Integrating (2.6.20) with respect to $x$ yields

$$
\begin{equation*}
F(x, y)=x^{2} y^{3} e^{y}+\phi(y) \tag{2.6.22}
\end{equation*}
$$

Differentiating this with respect to $y$ yields

$$
F_{y}=\left(3 x^{2} y^{2}+x^{2} y^{3}\right) e^{y}+\phi^{\prime}(y)
$$

and comparing this with (2.6.21) shows that $\phi^{\prime}(y)=e^{y}$. Therefore we set $\phi(y)=e^{y}$ in (2.6.22) and conclude that

$$
\left(x^{2} y^{3}+1\right) e^{y}=c
$$

is an implicit solution of (2.6.19). It is also an implicit solution of (2.6.18). Figure 2.6 .2 shows a direction field and some integral curves for (2.6.18).


Figure 2.6.2 A direction field and integral curves for $2 x y^{3} e^{y} d x+\left(3 x^{2} y^{2}+x^{2} y^{3}+1\right) e^{y} d y=0$
Theorem 2.6.1 does not apply in the next example, but the more general argument that led to Theorem 2.6.1 provides an integrating factor.

Example 2.6.3 Find an integrating factor for

$$
\begin{equation*}
\left(3 x y+6 y^{2}\right) d x+\left(2 x^{2}+9 x y\right) d y=0 \tag{2.6.23}
\end{equation*}
$$

and solve the equation.

Solution In (2.6.23)

$$
M=3 x y+6 y^{2}, N=2 x^{2}+9 x y
$$

and

$$
M_{y}-N_{x}=(3 x+12 y)-(4 x+9 y)=-x+3 y
$$

Therefore

$$
\frac{M_{y}-N_{x}}{M}=\frac{-x+3 y}{3 x y+6 y^{2}} \quad \text { and } \quad \frac{N_{x}-M_{y}}{N}=\frac{x-3 y}{2 x^{2}+9 x y}
$$

so Theorem 2.6.1 does not apply. Following the more general argument that led to Theorem 2.6.1, we look for functions $p=p(x)$ and $q=q(y)$ such that

$$
M_{y}-N_{x}=p(x) N-q(y) M
$$

that is,

$$
-x+3 y=p(x)\left(2 x^{2}+9 x y\right)-q(y)\left(3 x y+6 y^{2}\right)
$$

Since the left side contains only first degree terms in $x$ and $y$, we rewrite this equation as

$$
x p(x)(2 x+9 y)-y q(y)(3 x+6 y)=-x+3 y
$$

This will be an identity if

$$
\begin{equation*}
x p(x)=A \quad \text { and } \quad y q(y)=B \tag{2.6.24}
\end{equation*}
$$



Figure 2.6.3 A direction field and integral curves for $\left(3 x y+6 y^{2}\right) d x+\left(2 x^{2}+9 x y\right) d y=0$
where $A$ and $B$ are constants such that

$$
-x+3 y=A(2 x+9 y)-B(3 x+6 y)
$$

or, equivalently,

$$
-x+3 y=(2 A-3 B) x+(9 A-6 B) y
$$

Equating the coefficients of $x$ and $y$ on both sides shows that the last equation holds for all $(x, y)$ if

$$
\begin{aligned}
& 2 A-3 B=-1 \\
& 9 A-6 B=3,
\end{aligned}
$$

which has the solution $A=1, B=1$. Therefore (2.6.24) implies that

$$
p(x)=\frac{1}{x} \quad \text { and } \quad q(y)=\frac{1}{y} .
$$

Since

$$
\int p(x) d x=\ln |x| \quad \text { and } \quad \int q(y) d y=\ln |y|
$$

we can let $P(x)=x$ and $Q(y)=y$; hence, $\mu(x, y)=x y$ is an integrating factor. Multiplying (2.6.23) by $\mu$ yields the exact equation

$$
\left(3 x^{2} y^{2}+6 x y^{3}\right) d x+\left(2 x^{3} y+9 x^{2} y^{2}\right) d y=0
$$

We leave it to you to use the method of Section 2.5 to show that this equation has the implicit solution

$$
\begin{equation*}
x^{3} y^{2}+3 x^{2} y^{3}=c \tag{2.6.25}
\end{equation*}
$$

This is also an implicit solution of (2.6.23). Since $x \equiv 0$ and $y \equiv 0$ satisfy (2.6.25), you should check to see that $x \equiv 0$ and $y \equiv 0$ are also solutions of (2.6.23). (Why is it necesary to check this?)

Figure 2.6.3 shows a direction field and integral curves for (2.6.23).
See Exercise 28 for a general discussion of equations like (2.6.23).
Example 2.6.4 The separable equation

$$
\begin{equation*}
-y d x+\left(x+x^{6}\right) d y=0 \tag{2.6.26}
\end{equation*}
$$

can be converted to the exact equation

$$
\begin{equation*}
-\frac{d x}{x+x^{6}}+\frac{d y}{y}=0 \tag{2.6.27}
\end{equation*}
$$

by multiplying through by the integrating factor

$$
\mu(x, y)=\frac{1}{y\left(x+x^{6}\right)} .
$$

However, to solve (2.6.27) by the method of Section 2.5 we would have to evaluate the nasty integral

$$
\int \frac{d x}{x+x^{6}}
$$

Instead, we solve (2.6.26) explicitly for $y$ by finding an integrating factor of the form $\mu(x, y)=x^{a} y^{b}$.


Figure 2.6.4 A direction field and integral curves for $-y d x+\left(x+x^{6}\right) d y=0$

Solution In (2.6.26)

$$
M=-y, N=x+x^{6}
$$

and

$$
M_{y}-N_{x}=-1-\left(1+6 x^{5}\right)=-2-6 x^{5} .
$$

We look for functions $p=p(x)$ and $q=q(y)$ such that

$$
M_{y}-N_{x}=p(x) N-q(y) M
$$

that is,

$$
\begin{equation*}
-2-6 x^{5}=p(x)\left(x+x^{6}\right)+q(y) y . \tag{2.6.28}
\end{equation*}
$$

The right side will contain the term $-6 x^{5}$ if $p(x)=-6 / x$. Then (2.6.28) becomes

$$
-2-6 x^{5}=-6-6 x^{5}+q(y) y,
$$

so $q(y)=4 / y$. Since

$$
\int p(x) d x=-\int \frac{6}{x} d x=-6 \ln |x|=\ln \frac{1}{x^{6}},
$$

and

$$
\int q(y) d y=\int \frac{4}{y} d y=4 \ln |y|=\ln y^{4}
$$

we can take $P(x)=x^{-6}$ and $Q(y)=y^{4}$, which yields the integrating factor $\mu(x, y)=x^{-6} y^{4}$. Multiplying (2.6.26) by $\mu$ yields the exact equation

$$
-\frac{y^{5}}{x^{6}} d x+\left(\frac{y^{4}}{x^{5}}+y^{4}\right) d y=0
$$

We leave it to you to use the method of the Section 2.5 to show that this equation has the implicit solution

$$
\left(\frac{y}{x}\right)^{5}+y^{5}=k
$$

Solving for $y$ yields

$$
y=k^{1 / 5} x\left(1+x^{5}\right)^{-1 / 5}
$$

which we rewrite as

$$
y=c x\left(1+x^{5}\right)^{-1 / 5}
$$

by renaming the arbitrary constant. This is also a solution of (2.6.26).
Figure 2.6.4 shows a direction field and some integral curves for (2.6.26).

### 2.6 Exercises

1. (a) Verify that $\mu(x, y)=y$ is an integrating factor for

$$
\begin{equation*}
y d x+\left(2 x+\frac{1}{y}\right) d y=0 \tag{A}
\end{equation*}
$$

on any open rectangle that does not intersect the $x$ axis or, equivalently, that

$$
\begin{equation*}
y^{2} d x+(2 x y+1) d y=0 \tag{B}
\end{equation*}
$$

is exact on any such rectangle.
(b) Verify that $y \equiv 0$ is a solution of (B), but not of (A).
(c) Show that

$$
\begin{equation*}
y(x y+1)=c \tag{C}
\end{equation*}
$$

is an implicit solution of (B), and explain why every differentiable function $y=y(x)$ other than $y \equiv 0$ that satisfies $(\mathrm{C})$ is also a solution of (A).
2. (a) Verify that $\mu(x, y)=1 /(x-y)^{2}$ is an integrating factor for

$$
\begin{equation*}
-y^{2} d x+x^{2} d y=0 \tag{A}
\end{equation*}
$$

on any open rectangle that does not intersect the line $y=x$ or, equivalently, that

$$
\begin{equation*}
-\frac{y^{2}}{(x-y)^{2}} d x+\frac{x^{2}}{(x-y)^{2}} d y=0 \tag{B}
\end{equation*}
$$

is exact on any such rectangle.
(b) Use Theorem 2.2.1 to show that

$$
\begin{equation*}
\frac{x y}{(x-y)}=c \tag{C}
\end{equation*}
$$

is an implicit solution of (B), and explain why it's also an implicit solution of (A)
(c) Verify that $y=x$ is a solution of (A), even though it can't be obtained from (C).

In Exercises 3-16 find an integrating factor; that is a function of only one variable, and solve the given equation.
3. $y d x-x d y=0$
4. $3 x^{2} y d x+2 x^{3} d y=0$
5. $2 y^{3} d x+3 y^{2} d y=0$
6. $(5 x y+2 y+5) d x+2 x d y=0$
7. $(x y+x+2 y+1) d x+(x+1) d y=0$
8. $\left(27 x y^{2}+8 y^{3}\right) d x+\left(18 x^{2} y+12 x y^{2}\right) d y=0$
9. $\left(6 x y^{2}+2 y\right) d x+\left(12 x^{2} y+6 x+3\right) d y=0$
10. $y^{2} d x+\left(x y^{2}+3 x y+\frac{1}{y}\right) d y=0$
11. $\left(12 x^{3} y+24 x^{2} y^{2}\right) d x+\left(9 x^{4}+32 x^{3} y+4 y\right) d y=0$
12. $\left(x^{2} y+4 x y+2 y\right) d x+\left(x^{2}+x\right) d y=0$
13. $-y d x+\left(x^{4}-x\right) d y=0$
14. $\cos x \cos y d x+(\sin x \cos y-\sin x \sin y+y) d y=0$
15. $\left(2 x y+y^{2}\right) d x+\left(2 x y+x^{2}-2 x^{2} y^{2}-2 x y^{3}\right) d y=0$
16. $y \sin y d x+x(\sin y-y \cos y) d y=0$

In Exercises 17-23 find an integrating factor of the form $\mu(x, y)=P(x) Q(y)$ and solve the given equation.
17. $y(1+5 \ln |x|) d x+4 x \ln |x| d y=0$
18. $(\alpha y+\gamma x y) d x+(\beta x+\delta x y) d y=0$
19. $\left(3 x^{2} y^{3}-y^{2}+y\right) d x+(-x y+2 x) d y=0$
20. $2 y d x+3\left(x^{2}+x^{2} y^{3}\right) d y=0$
21. $(a \cos x y-y \sin x y) d x+(b \cos x y-x \sin x y) d y=0$
22. $x^{4} y^{4} d x+x^{5} y^{3} d y=0$
23. $y(x \cos x+2 \sin x) d x+x(y+1) \sin x d y=0$

In Exercises 24-27 find an integrating factor and solve the equation. Plot a direction field and some integral curves for the equation in the indicated rectangular region.
24. $\mathrm{C} / \mathrm{G}\left(x^{4} y^{3}+y\right) d x+\left(x^{5} y^{2}-x\right) d y=0 ; \quad\{-1 \leq x \leq 1,-1 \leq y \leq 1\}$
25. C/G $\left(3 x y+2 y^{2}+y\right) d x+\left(x^{2}+2 x y+x+2 y\right) d y=0 ; \quad\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$
26. C/G $\left(12 x y+6 y^{3}\right) d x+\left(9 x^{2}+10 x y^{2}\right) d y=0 ; \quad\{-2 \leq x \leq 2,-2 \leq y \leq 2\}$
27. $\mathrm{C} / \mathrm{G}\left(3 x^{2} y^{2}+2 y\right) d x+2 x d y=0 ; \quad\{-4 \leq x \leq 4,-4 \leq y \leq 4\}$
28. Suppose $a, b, c$, and $d$ are constants such that $a d-b c \neq 0$, and let $m$ and $n$ be arbitrary real numbers. Show that

$$
\left(a x^{m} y+b y^{n+1}\right) d x+\left(c x^{m+1}+d x y^{n}\right) d y=0
$$

has an integrating factor $\mu(x, y)=x^{\alpha} y^{\beta}$.
29. Suppose $M, N, M_{x}$, and $N_{y}$ are continuous for all $(x, y)$, and $\mu=\mu(x, y)$ is an integrating factor for

$$
\begin{equation*}
M(x, y) d x+N(x, y) d y=0 . \tag{A}
\end{equation*}
$$

Assume that $\mu_{x}$ and $\mu_{y}$ are continuous for all $(x, y)$, and suppose $y=y(x)$ is a differentiable function such that $\mu(x, y(x))=0$ and $\mu_{x}(x, y(x)) \neq 0$ for all $x$ in some interval $I$. Show that $y$ is a solution of (A) on $I$.
30. According to Theorem 2.1.2, the general solution of the linear nonhomogeneous equation

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{A}
\end{equation*}
$$

is

$$
\begin{equation*}
y=y_{1}(x)\left(c+\int f(x) / y_{1}(x) d x\right) \tag{B}
\end{equation*}
$$

where $y_{1}$ is any nontrivial solution of the complementary equation $y^{\prime}+p(x) y=0$. In this exercise we obtain this conclusion in a different way. You may find it instructive to apply the method suggested here to solve some of the exercises in Section 2.1.
(a) Rewrite (A) as

$$
\begin{equation*}
[p(x) y-f(x)] d x+d y=0 \tag{C}
\end{equation*}
$$

and show that $\mu= \pm e^{\int p(x) d x}$ is an integrating factor for (C).
(b) Multiply (A) through by $\mu= \pm e^{\int p(x) d x}$ and verify that the resulting equation can be rewritten as

$$
(\mu(x) y)^{\prime}=\mu(x) f(x)
$$

Then integrate both sides of this equation and solve for $y$ to show that the general solution of (A) is

$$
y=\frac{1}{\mu(x)}\left(c+\int f(x) \mu(x) d x\right) .
$$

Why is this form of the general solution equivalent to (B)?

# CHAPTER 3 Numerical Methods 

In this chapter we study numerical methods for solving a first order differential equation

$$
y^{\prime}=f(x, y)
$$

SECTION 3.1 deals with Euler's method, which is really too crude to be of much use in practical applications. However, its simplicity allows for an introduction to the ideas required to understand the better methods discussed in the other two sections.

SECTION 3.2 discusses improvements on Euler's method.
SECTION 3.3 deals with the Runge-Kutta method, perhaps the most widely used method for numerical solution of differential equations.

### 3.1 EULER'S METHOD

If an initial value problem

$$
\begin{equation*}
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0} \tag{3.1.1}
\end{equation*}
$$

can't be solved analytically, it's necessary to resort to numerical methods to obtain useful approximations to a solution of (3.1.1). We'll consider such methods in this chapter.

We're interested in computing approximate values of the solution of (3.1.1) at equally spaced points $x_{0}, x_{1}, \ldots, x_{n}=b$ in an interval $\left[x_{0}, b\right]$. Thus,

$$
x_{i}=x_{0}+i h, \quad i=0,1, \ldots, n,
$$

where

$$
h=\frac{b-x_{0}}{n} .
$$

We'll denote the approximate values of the solution at these points by $y_{0}, y_{1}, \ldots, y_{n}$; thus, $y_{i}$ is an approximation to $y\left(x_{i}\right)$. We'll call

$$
e_{i}=y\left(x_{i}\right)-y_{i}
$$

the error at the ith step. Because of the initial condition $y\left(x_{0}\right)=y_{0}$, we'll always have $e_{0}=0$. However, in general $e_{i} \neq 0$ if $i>0$.

We encounter two sources of error in applying a numerical method to solve an initial value problem:

- The formulas defining the method are based on some sort of approximation. Errors due to the inaccuracy of the approximation are called truncation errors.
- Computers do arithmetic with a fixed number of digits, and therefore make errors in evaluating the formulas defining the numerical methods. Errors due to the computer's inability to do exact arithmetic are called roundoff errors.

Since a careful analysis of roundoff error is beyond the scope of this book, we'll consider only truncation errors.

## Euler's Method

The simplest numerical method for solving (3.1.1) is Euler's method. This method is so crude that it is seldom used in practice; however, its simplicity makes it useful for illustrative purposes.

Euler's method is based on the assumption that the tangent line to the integral curve of (3.1.1) at $\left(x_{i}, y\left(x_{i}\right)\right)$ approximates the integral curve over the interval $\left[x_{i}, x_{i+1}\right]$. Since the slope of the integral curve of (3.1.1) at $\left(x_{i}, y\left(x_{i}\right)\right)$ is $y^{\prime}\left(x_{i}\right)=f\left(x_{i}, y\left(x_{i}\right)\right)$, the equation of the tangent line to the integral curve at $\left(x_{i}, y\left(x_{i}\right)\right)$ is

$$
\begin{equation*}
y=y\left(x_{i}\right)+f\left(x_{i}, y\left(x_{i}\right)\right)\left(x-x_{i}\right) . \tag{3.1.2}
\end{equation*}
$$

Setting $x=x_{i+1}=x_{i}+h$ in (3.1.2) yields

$$
\begin{equation*}
y_{i+1}=y\left(x_{i}\right)+h f\left(x_{i}, y\left(x_{i}\right)\right) \tag{3.1.3}
\end{equation*}
$$

as an approximation to $y\left(x_{i+1}\right)$. Since $y\left(x_{0}\right)=y_{0}$ is known, we can use (3.1.3) with $i=0$ to compute

$$
y_{1}=y_{0}+h f\left(x_{0}, y_{0}\right) .
$$

However, setting $i=1$ in (3.1.3) yields

$$
y_{2}=y\left(x_{1}\right)+h f\left(x_{1}, y\left(x_{1}\right)\right),
$$

which isn't useful, since we don't know $y\left(x_{1}\right)$. Therefore we replace $y\left(x_{1}\right)$ by its approximate value $y_{1}$ and redefine

$$
y_{2}=y_{1}+h f\left(x_{1}, y_{1}\right) .
$$

Having computed $y_{2}$, we can compute

$$
y_{3}=y_{2}+h f\left(x_{2}, y_{2}\right)
$$

In general, Euler's method starts with the known value $y\left(x_{0}\right)=y_{0}$ and computes $y_{1}, y_{2}, \ldots, y_{n}$ successively by with the formula

$$
\begin{equation*}
y_{i+1}=y_{i}+h f\left(x_{i}, y_{i}\right), \quad 0 \leq i \leq n-1 \tag{3.1.4}
\end{equation*}
$$

The next example illustrates the computational procedure indicated in Euler's method.

Example 3.1.1 Use Euler's method with $h=0.1$ to find approximate values for the solution of the initial value problem

$$
\begin{equation*}
y^{\prime}+2 y=x^{3} e^{-2 x}, \quad y(0)=1 \tag{3.1.5}
\end{equation*}
$$

at $x=0.1,0.2,0.3$.

Solution We rewrite (3.1.5) as

$$
y^{\prime}=-2 y+x^{3} e^{-2 x}, \quad y(0)=1,
$$

which is of the form (3.1.1), with

$$
f(x, y)=-2 y+x^{3} e^{-2 x}, x_{0}=0, \text { and } y_{0}=1
$$

Euler's method yields

$$
\begin{aligned}
y_{1} & =y_{0}+h f\left(x_{0}, y_{0}\right) \\
& =1+(.1) f(0,1)=1+(.1)(-2)=.8 \\
y_{2} & =y_{1}+h f\left(x_{1}, y_{1}\right) \\
& =.8+(.1) f(.1, .8)=.8+(.1)\left(-2(.8)+(.1)^{3} e^{-.2}\right)=.640081873, \\
y_{3} & =y_{2}+h f\left(x_{2}, y_{2}\right) \\
& =.640081873+(.1)\left(-2(.640081873)+(.2)^{3} e^{-.4}\right)=.512601754 .
\end{aligned}
$$

We've written the details of these computations to ensure that you understand the procedure. However, in the rest of the examples as well as the exercises in this chapter, we'll assume that you can use a programmable calculator or a computer to carry out the necessary computations.

Examples Illustrating The Error in Euler's Method
Example 3.1.2 Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+2 y=x^{3} e^{-2 x}, \quad y(0)=1
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$. Compare these approximate values with the values of the exact solution

$$
\begin{equation*}
y=\frac{e^{-2 x}}{4}\left(x^{4}+4\right) \tag{3.1.6}
\end{equation*}
$$

which can be obtained by the method of Section 2.1. (Verify.)

Solution Table 3.1.1 shows the values of the exact solution (3.1.6) at the specified points, and the approximate values of the solution at these points obtained by Euler's method with step sizes $h=0.1$, $h=0.05$, and $h=0.025$. In examining this table, keep in mind that the approximate values in the column corresponding to $h=.05$ are actually the results of 20 steps with Euler's method. We haven't listed the estimates of the solution obtained for $x=0.05,0.15, \ldots$, since there's nothing to compare them with in the column corresponding to $h=0.1$. Similarly, the approximate values in the column corresponding to $h=0.025$ are actually the results of 40 steps with Euler's method.

Table 3.1.1. Numerical solution of $y^{\prime}+2 y=x^{3} e^{-2 x}, y(0)=1$, by Euler's method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 0.0 | 1.000000000 | 1.000000000 | 1.000000000 | 1.000000000 |
| 0.1 | 0.800000000 | 0.810005655 | 0.814518349 | 0.818751221 |
| 0.2 | 0.640081873 | 0.656266437 | 0.663635953 | 0.670588174 |
| 0.3 | 0.512601754 | 0.532290981 | 0.541339495 | 0.549922980 |
| 0.4 | 0.411563195 | 0.432887056 | 0.442774766 | 0.452204669 |
| 0.5 | 0.332126261 | 0.353785015 | 0.363915597 | 0.373627557 |
| 0.6 | 0.270299502 | 0.291404256 | 0.301359885 | 0.310952904 |
| 0.7 | 0.222745397 | 0.242707257 | 0.252202935 | 0.261398947 |
| 0.8 | 0.186654593 | 0.205105754 | 0.213956311 | 0.222570721 |
| 0.9 | 0.159660776 | 0.176396883 | 0.184492463 | 0.192412038 |
| 1.0 | 0.139778910 | 0.154715925 | 0.162003293 | 0.169169104 |

You can see from Table 3.1.1 that decreasing the step size improves the accuracy of Euler's method. For example,

$$
y_{\text {exact }}(1)-y_{\text {approx }}(1) \approx\left\{\begin{array}{l}
.0293 \text { with } h=0.1 \\
.0144 \text { with } h=0.05 \\
.0071 \text { with } h=0.025
\end{array}\right.
$$

Based on this scanty evidence, you might guess that the error in approximating the exact solution at a fixed value of $x$ by Euler's method is roughly halved when the step size is halved. You can find more evidence to support this conjecture by examining Table 3.1.2, which lists the approximate values of $y_{\text {exact }}-y_{\text {approx }}$ at $x=0.1,0.2, \ldots, 1.0$.

Table 3.1.2. Errors in approximate solutions of $y^{\prime}+2 y=x^{3} e^{-2 x}, y(0)=1$, obtained by Euler's method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ |
| :---: | ---: | ---: | ---: |
| 0.1 | 0.0187 | 0.0087 | 0.0042 |
| 0.2 | 0.0305 | 0.0143 | 0.0069 |
| 0.3 | 0.0373 | 0.0176 | 0.0085 |
| 0.4 | 0.0406 | 0.0193 | 0.0094 |
| 0.5 | 0.0415 | 0.0198 | 0.0097 |
| 0.6 | 0.0406 | 0.0195 | 0.0095 |
| 0.7 | 0.0386 | 0.0186 | 0.0091 |
| 0.8 | 0.0359 | 0.0174 | 0.0086 |
| 0.9 | 0.0327 | 0.0160 | 0.0079 |
| 1.0 | 0.0293 | 0.0144 | 0.0071 |

Example 3.1.3 Tables 3.1.3 and 3.1.4 show analogous results for the nonlinear initial value problem

$$
\begin{equation*}
y^{\prime}=-2 y^{2}+x y+x^{2}, y(0)=1 \tag{3.1.7}
\end{equation*}
$$

except in this case we can't solve (3.1.7) exactly. The results in the "Exact" column were obtained by using a more accurate numerical method known as the Runge-Kutta method with a small step size. They are exact to eight decimal places.

Since we think it's important in evaluating the accuracy of the numerical methods that we'll be studying in this chapter, we often include a column listing values of the exact solution of the initial value problem, even if the directions in the example or exercise don't specifically call for it. If quotation marks are included in the heading, the values were obtained by applying the Runge-Kutta method in a way that's explained in Section 3.3. If quotation marks are not included, the values were obtained from a known formula for the solution. In either case, the values are exact to eight places to the right of the decimal point.

Table 3.1.3. Numerical solution of $y^{\prime}=-2 y^{2}+x y+x^{2}, y(0)=1$, by Euler's method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 0.0 | 1.000000000 | 1.000000000 | 1.000000000 | 1.000000000 |
| 0.1 | 0.800000000 | 0.821375000 | 0.829977007 | 0.837584494 |
| 0.2 | 0.681000000 | 0.707795377 | 0.719226253 | 0.729641890 |
| 0.3 | 0.605867800 | 0.633776590 | 0.646115227 | 0.657580377 |
| 0.4 | 0.559628676 | 0.587454526 | 0.600045701 | 0.611901791 |
| 0.5 | 0.535376972 | 0.562906169 | 0.575556391 | 0.587575491 |
| 0.6 | 0.529820120 | 0.557143535 | 0.569824171 | 0.581942225 |
| 0.7 | 0.541467455 | 0.568716935 | 0.581435423 | 0.593629526 |
| 0.8 | 0.569732776 | 0.596951988 | 0.609684903 | 0.621907458 |
| 0.9 | 0.614392311 | 0.641457729 | 0.654110862 | 0.666250842 |
| 1.0 | 0.675192037 | 0.701764495 | 0.714151626 | 0.726015790 |

Table 3.1.4. Errors in approximate solutions of $y^{\prime}=-2 y^{2}+x y+x^{2}, y(0)=1$, obtained by Euler's method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ |
| :---: | ---: | ---: | ---: |
| 0.1 | 0.0376 | 0.0162 | 0.0076 |
| 0.2 | 0.0486 | 0.0218 | 0.0104 |
| 0.3 | 0.0517 | 0.0238 | 0.0115 |
| 0.4 | 0.0523 | 0.0244 | 0.0119 |
| 0.5 | 0.0522 | 0.0247 | 0.0121 |
| 0.6 | 0.0521 | 0.0248 | 0.0121 |
| 0.7 | 0.0522 | 0.0249 | 0.0122 |
| 0.8 | 0.0522 | 0.0250 | 0.0122 |
| 0.9 | 0.0519 | 0.0248 | 0.0121 |
| 1.0 | 0.0508 | 0.0243 | 0.0119 |

Truncation Error in Euler's Method
Consistent with the results indicated in Tables 3.1.1-3.1.4, we'll now show that under reasonable assumptions on $f$ there's a constant $K$ such that the error in approximating the solution of the initial value problem

$$
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0},
$$

at a given point $b>x_{0}$ by Euler's method with step size $h=\left(b-x_{0}\right) / n$ satisfies the inequality

$$
\left|y(b)-y_{n}\right| \leq K h,
$$

where $K$ is a constant independent of $n$.
There are two sources of error (not counting roundoff) in Euler's method:

1. The error committed in approximating the integral curve by the tangent line (3.1.2) over the interval $\left[x_{i}, x_{i+1}\right]$.
2. The error committed in replacing $y\left(x_{i}\right)$ by $y_{i}$ in (3.1.2) and using (3.1.4) rather than (3.1.2) to compute $y_{i+1}$.

Euler's method assumes that $y_{i+1}$ defined in (3.1.2) is an approximation to $y\left(x_{i+1}\right)$. We call the error in this approximation the local truncation error at the ith step, and denote it by $T_{i}$; thus,

$$
\begin{equation*}
T_{i}=y\left(x_{i+1}\right)-y\left(x_{i}\right)-h f\left(x_{i}, y\left(x_{i}\right)\right) . \tag{3.1.8}
\end{equation*}
$$

We'll now use Taylor's theorem to estimate $T_{i}$, assuming for simplicity that $f, f_{x}$, and $f_{y}$ are continuous and bounded for all $(x, y)$. Then $y^{\prime \prime}$ exists and is bounded on $\left[x_{0}, b\right]$. To see this, we differentiate

$$
y^{\prime}(x)=f(x, y(x))
$$

to obtain

$$
\begin{aligned}
y^{\prime \prime}(x) & =f_{x}(x, y(x))+f_{y}(x, y(x)) y^{\prime}(x) \\
& =f_{x}(x, y(x))+f_{y}(x, y(x)) f(x, y(x))
\end{aligned}
$$

Since we assumed that $f, f_{x}$ and $f_{y}$ are bounded, there's a constant $M$ such that

$$
\left|f_{x}(x, y(x))+f_{y}(x, y(x)) f(x, y(x))\right| \leq M, \quad x_{0}<x<b,
$$

which implies that

$$
\begin{equation*}
\left|y^{\prime \prime}(x)\right| \leq M, \quad x_{0}<x<b \tag{3.1.9}
\end{equation*}
$$

Since $x_{i+1}=x_{i}+h$, Taylor's theorem implies that

$$
y\left(x_{i+1}\right)=y\left(x_{i}\right)+h y^{\prime}\left(x_{i}\right)+\frac{h^{2}}{2} y^{\prime \prime}\left(\tilde{x}_{i}\right),
$$

where $\tilde{x}_{i}$ is some number between $x_{i}$ and $x_{i+1}$. Since $y^{\prime}\left(x_{i}\right)=f\left(x_{i}, y\left(x_{i}\right)\right)$ this can be written as

$$
y\left(x_{i+1}\right)=y\left(x_{i}\right)+h f\left(x_{i}, y\left(x_{i}\right)\right)+\frac{h^{2}}{2} y^{\prime \prime}\left(\tilde{x}_{i}\right),
$$

or, equivalently,

$$
y\left(x_{i+1}\right)-y\left(x_{i}\right)-h f\left(x_{i}, y\left(x_{i}\right)\right)=\frac{h^{2}}{2} y^{\prime \prime}\left(\tilde{x}_{i}\right) .
$$

Comparing this with (3.1.8) shows that

$$
T_{i}=\frac{h^{2}}{2} y^{\prime \prime}\left(\tilde{x}_{i}\right)
$$

Recalling (3.1.9), we can establish the bound

$$
\begin{equation*}
\left|T_{i}\right| \leq \frac{M h^{2}}{2}, \quad 1 \leq i \leq n \tag{3.1.10}
\end{equation*}
$$

Although it may be difficult to determine the constant $M$, what is important is that there's an $M$ such that (3.1.10) holds. We say that the local truncation error of Euler's method is of order $h^{2}$, which we write as $O\left(h^{2}\right)$.

Note that the magnitude of the local truncation error in Euler's method is determined by the second derivative $y^{\prime \prime}$ of the solution of the initial value problem. Therefore the local truncation error will be larger where $\left|y^{\prime \prime}\right|$ is large, or smaller where $\left|y^{\prime \prime}\right|$ is small.

Since the local truncation error for Euler's method is $O\left(h^{2}\right)$, it's reasonable to expect that halving $h$ reduces the local truncation error by a factor of 4 . This is true, but halving the step size also requires twice as many steps to approximate the solution at a given point. To analyze the overall effect of truncation error in Euler's method, it's useful to derive an equation relating the errors

$$
e_{i+1}=y\left(x_{i+1}\right)-y_{i+1} \quad \text { and } \quad e_{i}=y\left(x_{i}\right)-y_{i}
$$

To this end, recall that

$$
\begin{equation*}
y\left(x_{i+1}\right)=y\left(x_{i}\right)+h f\left(x_{i}, y\left(x_{i}\right)\right)+T_{i} \tag{3.1.11}
\end{equation*}
$$

and

$$
\begin{equation*}
y_{i+1}=y_{i}+h f\left(x_{i}, y_{i}\right) \tag{3.1.12}
\end{equation*}
$$

Subtracting (3.1.12) from (3.1.11) yields

$$
\begin{equation*}
e_{i+1}=e_{i}+h\left[f\left(x_{i}, y\left(x_{i}\right)\right)-f\left(x_{i}, y_{i}\right)\right]+T_{i} . \tag{3.1.13}
\end{equation*}
$$

The last term on the right is the local truncation error at the $i$ th step. The other terms reflect the way errors made at previous steps affect $e_{i+1}$. Since $\left|T_{i}\right| \leq M h^{2} / 2$, we see from (3.1.13) that

$$
\begin{equation*}
\left|e_{i+1}\right| \leq\left|e_{i}\right|+h\left|f\left(x_{i}, y\left(x_{i}\right)\right)-f\left(x_{i}, y_{i}\right)\right|+\frac{M h^{2}}{2} \tag{3.1.14}
\end{equation*}
$$

Since we assumed that $f_{y}$ is continuous and bounded, the mean value theorem implies that

$$
f\left(x_{i}, y\left(x_{i}\right)\right)-f\left(x_{i}, y_{i}\right)=f_{y}\left(x_{i}, y_{i}^{*}\right)\left(y\left(x_{i}\right)-y_{i}\right)=f_{y}\left(x_{i}, y_{i}^{*}\right) e_{i}
$$

where $y_{i}^{*}$ is between $y_{i}$ and $y\left(x_{i}\right)$. Therefore

$$
\left|f\left(x_{i}, y\left(x_{i}\right)\right)-f\left(x_{i}, y_{i}\right)\right| \leq R\left|e_{i}\right|
$$

for some constant $R$. From this and (3.1.14),

$$
\begin{equation*}
\left|e_{i+1}\right| \leq(1+R h)\left|e_{i}\right|+\frac{M h^{2}}{2}, \quad 0 \leq i \leq n-1 \tag{3.1.15}
\end{equation*}
$$

For convenience, let $C=1+R h$. Since $e_{0}=y\left(x_{0}\right)-y_{0}=0$, applying (3.1.15) repeatedly yields

$$
\begin{align*}
\left|e_{1}\right| & \leq \frac{M h^{2}}{2} \\
\left|e_{2}\right| & \leq C\left|e_{1}\right|+\frac{M h^{2}}{2} \leq(1+C) \frac{M h^{2}}{2} \\
\left|e_{3}\right| & \leq C\left|e_{2}\right|+\frac{M h^{2}}{2} \leq\left(1+C+C^{2}\right) \frac{M h^{2}}{2} \\
& \vdots  \tag{3.1.16}\\
\left|e_{n}\right| & \leq C\left|e_{n-1}\right|+\frac{M h^{2}}{2} \leq\left(1+C+\cdots+C^{n-1}\right) \frac{M h^{2}}{2}
\end{align*}
$$

Recalling the formula for the sum of a geometric series, we see that

$$
1+C+\cdots+C^{n-1}=\frac{1-C^{n}}{1-C}=\frac{(1+R h)^{n}-1}{R h}
$$

(since $C=1+R h$ ). From this and (3.1.16),

$$
\begin{equation*}
\left|y(b)-y_{n}\right|=\left|e_{n}\right| \leq \frac{(1+R h)^{n}-1}{R} \frac{M h}{2} . \tag{3.1.17}
\end{equation*}
$$

Since Taylor's theorem implies that

$$
1+R h<e^{R h}
$$

(verify),

$$
(1+R h)^{n}<e^{n R h}=e^{R\left(b-x_{0}\right)} \quad\left(\text { since } n h=b-x_{0}\right)
$$

This and (3.1.17) imply that

$$
\begin{equation*}
\left|y(b)-y_{n}\right| \leq K h, \tag{3.1.18}
\end{equation*}
$$

with

$$
K=M \frac{e^{R\left(b-x_{0}\right)}-1}{2 R}
$$

Because of (3.1.18) we say that the global truncation error of Euler's method is of order $h$, which we write as $O(h)$.
Semilinear Equations and Variation of Parameters
An equation that can be written in the form

$$
\begin{equation*}
y^{\prime}+p(x) y=h(x, y) \tag{3.1.19}
\end{equation*}
$$

with $p \not \equiv 0$ is said to be semilinear. (Of course, (3.1.19) is linear if $h$ is independent of $y$.) One way to apply Euler's method to an initial value problem

$$
\begin{equation*}
y^{\prime}+p(x) y=h(x, y), \quad y\left(x_{0}\right)=y_{0} \tag{3.1.20}
\end{equation*}
$$

for (3.1.19) is to think of it as

$$
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0}
$$

where

$$
f(x, y)=-p(x) y+h(x, y)
$$

However, we can also start by applying variation of parameters to (3.1.20), as in Sections 2.1 and 2.4; thus, we write the solution of (3.1.20) as $y=u y_{1}$, where $y_{1}$ is a nontrivial solution of the complementary equation $y^{\prime}+p(x) y=0$. Then $y=u y_{1}$ is a solution of (3.1.20) if and only if $u$ is a solution of the initial value problem

$$
\begin{equation*}
u^{\prime}=h\left(x, u y_{1}(x)\right) / y_{1}(x), \quad u\left(x_{0}\right)=y\left(x_{0}\right) / y_{1}\left(x_{0}\right) \tag{3.1.21}
\end{equation*}
$$

We can apply Euler's method to obtain approximate values $u_{0}, u_{1}, \ldots, u_{n}$ of this initial value problem, and then take

$$
y_{i}=u_{i} y_{1}\left(x_{i}\right)
$$

as approximate values of the solution of (3.1.20). We'll call this procedure the Euler semilinear method.
The next two examples show that the Euler and Euler semilinear methods may yield drastically different results.

Example 3.1.4 In Example 2.1.7 we had to leave the solution of the initial value problem

$$
\begin{equation*}
y^{\prime}-2 x y=1, \quad y(0)=3 \tag{3.1.22}
\end{equation*}
$$

in the form

$$
\begin{equation*}
y=e^{x^{2}}\left(3+\int_{0}^{x} e^{-t^{2}} d t\right) \tag{3.1.23}
\end{equation*}
$$

because it was impossible to evaluate this integral exactly in terms of elementary functions. Use step sizes $h=0.2, h=0.1$, and $h=0.05$ to find approximate values of the solution of (3.1.22) at $x=0,0.2$, $0.4,0.6, \ldots, 2.0$ by (a) Euler's method; (b) the Euler semilinear method.

SOLUTION(a) Rewriting (3.1.22) as

$$
\begin{equation*}
y^{\prime}=1+2 x y, \quad y(0)=3 \tag{3.1.24}
\end{equation*}
$$

and applying Euler's method with $f(x, y)=1+2 x y$ yields the results shown in Table 3.1.5. Because of the large differences between the estimates obtained for the three values of $h$, it would be clear that these results are useless even if the "exact" values were not included in the table.

Table 3.1.5. Numerical solution of $y^{\prime}-2 x y=1, y(0)=3$, with Euler's method.

| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| :---: | ---: | ---: | ---: | ---: |
| 0.0 | 3.000000000 | 3.000000000 | 3.000000000 | 3.000000000 |
| 0.2 | 3.200000000 | 3.262000000 | 3.294348537 | 3.327851973 |
| 0.4 | 3.656000000 | 3.802028800 | 3.881421103 | 3.966059348 |
| 0.6 | 4.440960000 | 4.726810214 | 4.888870783 | 5.067039535 |
| 0.8 | 5.706790400 | 6.249191282 | 6.570796235 | 6.936700945 |
| 1.0 | 7.732963328 | 8.771893026 | 9.419105620 | 10.184923955 |
| 1.2 | 11.026148659 | 13.064051391 | 14.405772067 | 16.067111677 |
| 1.4 | 16.518700016 | 20.637273893 | 23.522935872 | 27.289392347 |
| 1.6 | 25.969172024 | 34.570423758 | 41.033441257 | 50.000377775 |
| 1.8 | 42.789442120 | 61.382165543 | 76.491018246 | 98.982969504 |
| 2.0 | 73.797840446 | 115.440048291 | 152.363866569 | 211.954462214 |

It's easy to see why Euler's method yields such poor results. Recall that the constant $M$ in (3.1.10) which plays an important role in determining the local truncation error in Euler's method - must be an upper bound for the values of the second derivative $y^{\prime \prime}$ of the solution of the initial value problem (3.1.22) on $(0,2)$. The problem is that $y^{\prime \prime}$ assumes very large values on this interval. To see this, we differentiate (3.1.24) to obtain

$$
y^{\prime \prime}(x)=2 y(x)+2 x y^{\prime}(x)=2 y(x)+2 x(1+2 x y(x))=2\left(1+2 x^{2}\right) y(x)+2 x,
$$

where the second equality follows again from (3.1.24). Since (3.1.23) implies that $y(x)>3 e^{x^{2}}$ if $x>0$,

$$
y^{\prime \prime}(x)>6\left(1+2 x^{2}\right) e^{x^{2}}+2 x, \quad x>0
$$

For example, letting $x=2$ shows that $y^{\prime \prime}(2)>2952$.
 the Euler semilinear method to (3.1.22), with

$$
y=u e^{x^{2}} \quad \text { and } \quad u^{\prime}=e^{-x^{2}}, \quad u(0)=3
$$

The results listed in Table 3.1.6 are clearly better than those obtained by Euler's method.
Table 3.1.6. Numerical solution of $y^{\prime}-2 x y=1, y(0)=3$, by the Euler semilinear method.

| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| :---: | ---: | ---: | ---: | ---: |
| 0.0 | 3.000000000 | 3.000000000 | 3.000000000 | 3.000000000 |
| 0.2 | 3.330594477 | 3.329558853 | 3.328788889 | 3.327851973 |
| 0.4 | 3.980734157 | 3.974067628 | 3.970230415 | 3.966059348 |
| 0.6 | 5.106360231 | 5.087705244 | 5.077622723 | 5.067039535 |
| 0.8 | 7.021003417 | 6.980190891 | 6.958779586 | 6.936700945 |
| 1.0 | 10.350076600 | 10.269170824 | 10.227464299 | 10.184923955 |
| 1.2 | 16.381180092 | 16.226146390 | 16.147129067 | 16.067111677 |
| 1.4 | 27.890003380 | 27.592026085 | 27.441292235 | 27.289392347 |
| 1.6 | 51.183323262 | 50.594503863 | 50.298106659 | 50.000377775 |
| 1.8 | 101.424397595 | 100.206659076 | 99.595562766 | 98.982969504 |
| 2.0 | 217.301032800 | 214.631041938 | 213.293582978 | 211.954462214 |

We can't give a general procedure for determining in advance whether Euler's method or the semilinear Euler method will produce better results for a given semilinear initial value problem (3.1.19). As a rule of thumb, the Euler semilinear method will yield better results than Euler's method if $\left|u^{\prime \prime}\right|$ is small on $\left[x_{0}, b\right]$, while Euler's method yields better results if $\left|u^{\prime \prime}\right|$ is large on $\left[x_{0}, b\right]$. In many cases the results obtained by the two methods don't differ appreciably. However, we propose the an intuitive way to decide which is the better method: Try both methods with multiple step sizes, as we did in Example 3.1.4, and accept the results obtained by the method for which the approximations change less as the step size decreases.

Example 3.1.5 Applying Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to the initial value problem

$$
\begin{equation*}
y^{\prime}-2 y=\frac{x}{1+y^{2}}, \quad y(1)=7 \tag{3.1.25}
\end{equation*}
$$

on $[1,2]$ yields the results in Table 3.1.7. Applying the Euler semilinear method with

$$
y=u e^{2 x} \quad \text { and } \quad u^{\prime}=\frac{x e^{-2 x}}{1+u^{2} e^{4 x}}, \quad u(1)=7 e^{-2}
$$

yields the results in Table 3.1.8. Since the latter are clearly less dependent on step size than the former, we conclude that the Euler semilinear method is better than Euler's method for (3.1.25). This conclusion is supported by comparing the approximate results obtained by the two methods with the "exact" values of the solution.

Table 3.1.7. Numerical solution of $y^{\prime}-2 y=x /\left(1+y^{2}\right), y(1)=7$, by Euler's method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | ---: | ---: | ---: | ---: |
| 1.0 | 7.000000000 | 7.000000000 | 7.000000000 | 7.000000000 |
| 1.1 | 8.402000000 | 8.471970569 | 8.510493955 | 8.551744786 |
| 1.2 | 10.083936450 | 10.252570169 | 10.346014101 | 10.446546230 |
| 1.3 | 12.101892354 | 12.406719381 | 12.576720827 | 12.760480158 |
| 1.4 | 14.523152445 | 15.012952416 | 15.287872104 | 15.586440425 |
| 1.5 | 17.428443554 | 18.166277405 | 18.583079406 | 19.037865752 |
| 1.6 | 20.914624471 | 21.981638487 | 22.588266217 | 23.253292359 |
| 1.7 | 25.097914310 | 26.598105180 | 27.456479695 | 28.401914416 |
| 1.8 | 30.117766627 | 32.183941340 | 33.373738944 | 34.690375086 |
| 1.9 | 36.141518172 | 38.942738252 | 40.566143158 | 42.371060528 |
| 2.0 | 43.369967155 | 47.120835251 | 49.308511126 | 51.752229656 |

Table 3.1.8. Numerical solution of $y^{\prime}-2 y=x /\left(1+y^{2}\right), y(1)=7$, by the Euler semilinear method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | ---: | ---: | ---: | ---: |
| 1.0 | 7.000000000 | 7.000000000 | 7.000000000 | 7.000000000 |
| 1.1 | 8.552262113 | 8.551993978 | 8.551867007 | 8.551744786 |
| 1.2 | 10.447568674 | 10.447038547 | 10.446787646 | 10.446546230 |
| 1.3 | 12.762019799 | 12.761221313 | 12.760843543 | 12.760480158 |
| 1.4 | 15.588535141 | 15.587448600 | 15.586934680 | 15.586440425 |
| 1.5 | 19.040580614 | 19.039172241 | 19.038506211 | 19.037865752 |
| 1.6 | 23.256721636 | 23.254942517 | 23.254101253 | 23.253292359 |
| 1.7 | 28.406184597 | 28.403969107 | 28.402921581 | 28.401914416 |
| 1.8 | 34.695649222 | 34.692912768 | 34.691618979 | 34.690375086 |
| 1.9 | 42.377544138 | 42.374180090 | 42.372589624 | 42.371060528 |
| 2.0 | 51.760178446 | 51.756054133 | 51.754104262 | 51.752229656 |

Example 3.1.6 Applying Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to the initial value problem

$$
\begin{equation*}
y^{\prime}+3 x^{2} y=1+y^{2}, \quad y(2)=2 \tag{3.1.26}
\end{equation*}
$$

on $[2,3]$ yields the results in Table 3.1.9. Applying the Euler semilinear method with

$$
y=u e^{-x^{3}} \quad \text { and } \quad u^{\prime}=e^{x^{3}}\left(1+u^{2} e^{-2 x^{3}}\right), \quad u(2)=2 e^{8}
$$

yields the results in Table 3.1.10. Noting the close agreement among the three columns of Table 3.1.9 (at least for larger values of $x$ ) and the lack of any such agreement among the columns of Table 3.1.10, we conclude that Euler's method is better than the Euler semilinear method for (3.1.26). Comparing the results with the exact values supports this conclusion.

Table 3.1.9. Numerical solution of $y^{\prime}+3 x^{2} y=1+y^{2}, \quad y(2)=2$, by Euler's method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | 2.000000000 | 2.000000000 | 2.000000000 | 2.000000000 |
| 2.1 | 0.100000000 | 0.493231250 | 0.609611171 | 0.701162906 |
| 2.2 | 0.068700000 | 0.122879586 | 0.180113445 | 0.236986800 |
| 2.3 | 0.069419569 | 0.070670890 | 0.083934459 | 0.103815729 |
| 2.4 | 0.059732621 | 0.061338956 | 0.063337561 | 0.068390786 |
| 2.5 | 0.056871451 | 0.056002363 | 0.056249670 | 0.057281091 |
| 2.6 | 0.050560917 | 0.051465256 | 0.051517501 | 0.051711676 |
| 2.7 | 0.048279018 | 0.047484716 | 0.047514202 | 0.047564141 |
| 2.8 | 0.042925892 | 0.043967002 | 0.043989239 | 0.044014438 |
| 2.9 | 0.042148458 | 0.040839683 | 0.040857109 | 0.040875333 |
| 3.0 | 0.035985548 | 0.038044692 | 0.038058536 | 0.038072838 |

Table 3.1.10. Numerical solution of $y^{\prime}+3 x^{2} y=1+y^{2}, \quad y(2)=2$, by the Euler semilinear method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | $h=.0125$ |
| 2.0 | 2.000000000 | 2.000000000 | 2.000000000 | 2.000000000 |
| 2.1 | 0.708426286 | 0.702568171 | 0.701214274 | 0.701162906 |
| 2.2 | 0.214501852 | 0.222599468 | 0.228942240 | 0.236986800 |
| 2.3 | 0.069861436 | 0.083620494 | 0.092852806 | 0.103815729 |
| 2.4 | 0.032487396 | 0.047079261 | 0.056825805 | 0.068390786 |
| 2.5 | 0.021895559 | 0.036030018 | 0.045683801 | 0.057281091 |
| 2.6 | 0.017332058 | 0.030750181 | 0.040189920 | 0.051711676 |
| 2.7 | 0.014271492 | 0.026931911 | 0.036134674 | 0.047564141 |
| 2.8 | 0.011819555 | 0.023720670 | 0.032679767 | 0.044014438 |
| 2.9 | 0.009776792 | 0.020925522 | 0.029636506 | 0.040875333 |
| 3.0 | 0.008065020 | 0.018472302 | 0.026931099 | 0.038072838 |

In the next two sections we'll study other numerical methods for solving initial value problems, called the improved Euler method, the midpoint method, Heun's method and the Runge-Kutta method. If the initial value problem is semilinear as in (3.1.19), we also have the option of using variation of parameters and then applying the given numerical method to the initial value problem (3.1.21) for $u$. By analogy with the terminology used here, we'll call the resulting procedure the improved Euler semilinear method, the midpoint semilinear method, Heun's semilinear method or the Runge-Kutta semilinear method, as the case may be.

### 3.1 Exercises

You may want to save the results of these exercises, since we'll revisit in the next two sections. In Exercises 1-5 use Euler's method to find approximate values of the solution of the given initial value problem at the points $x_{i}=x_{0}+i h$, where $x_{0}$ is the point where the initial condition is imposed and $i=1,2,3$. The purpose of these exercises is to familiarize you with the computational procedure of Euler's method.

1. C $y^{\prime}=2 x^{2}+3 y^{2}-2, \quad y(2)=1 ; \quad h=0.05$
2. C $y^{\prime}=y+\sqrt{x^{2}+y^{2}}, \quad y(0)=1 ; \quad h=0.1$
3. C $y^{\prime}+3 y=x^{2}-3 x y+y^{2}, \quad y(0)=2 ; \quad h=0.05$
4. $\mathrm{C} y^{\prime}=\frac{1+x}{1-y^{2}}, \quad y(2)=3 ; \quad h=0.1$
5. C $y^{\prime}+x^{2} y=\sin x y, \quad y(1)=\pi ; \quad h=0.2$
6. C Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+3 y=7 e^{4 x}, \quad y(0)=2
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$. Compare these approximate values with the values of the exact solution $y=e^{4 x}+e^{-3 x}$, which can be obtained by the method of Section 2.1. Present your results in a table like Table 3.1.1.
7. C Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+\frac{2}{x} y=\frac{3}{x^{3}}+1, \quad y(1)=1
$$

at $x=1.0,1.1,1.2,1.3, \ldots, 2.0$. Compare these approximate values with the values of the exact solution

$$
y=\frac{1}{3 x^{2}}\left(9 \ln x+x^{3}+2\right),
$$

which can be obtained by the method of Section 2.1. Present your results in a table like Table 3.1.1.
8. C Use Euler's method with step sizes $h=0.05, h=0.025$, and $h=0.0125$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}=\frac{y^{2}+x y-x^{2}}{x^{2}}, \quad y(1)=2
$$

at $x=1.0,1.05,1.10,1.15, \ldots, 1.5$. Compare these approximate values with the values of the exact solution

$$
y=\frac{x\left(1+x^{2} / 3\right)}{1-x^{2} / 3}
$$

obtained in Example ??. Present your results in a table like Table 3.1.1.
9. C In Example 2.2.3 it was shown that

$$
y^{5}+y=x^{2}+x-4
$$

is an implicit solution of the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{2 x+1}{5 y^{4}+1}, \quad y(2)=1 . \tag{A}
\end{equation*}
$$

Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of (A) at $x=2.0,2.1,2.2,2.3, \ldots, 3.0$. Present your results in tabular form. To check the error in these approximate values, construct another table of values of the residual

$$
R(x, y)=y^{5}+y-x^{2}-x+4
$$

for each value of $(x, y)$ appearing in the first table.
10. C You can see from Example 2.5.1 that

$$
x^{4} y^{3}+x^{2} y^{5}+2 x y=4
$$

is an implicit solution of the initial value problem

$$
\begin{equation*}
y^{\prime}=-\frac{4 x^{3} y^{3}+2 x y^{5}+2 y}{3 x^{4} y^{2}+5 x^{2} y^{4}+2 x}, \quad y(1)=1 . \tag{A}
\end{equation*}
$$

Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of (A) at $x=1.0,1.1,1.2,1.3, \ldots, 2.0$. Present your results in tabular form. To check the error in these approximate values, construct another table of values of the residual

$$
R(x, y)=x^{4} y^{3}+x^{2} y^{5}+2 x y-4
$$

for each value of $(x, y)$ appearing in the first table.
11. C Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
\left(3 y^{2}+4 y\right) y^{\prime}+2 x+\cos x=0, \quad y(0)=1 ; \quad(\text { Exercise 2.2.13) }
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$.
12. C Use Euler's method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+\frac{(y+1)(y-1)(y-2)}{x+1}=0, \quad y(1)=0(\text { Exercise 2.2.14 })
$$

at $x=1.0,1.1,1.2,1.3, \ldots, 2.0$.
13. C Use Euler's method and the Euler semilinear method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+3 y=7 e^{-3 x}, \quad y(0)=6
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$. Compare these approximate values with the values of the exact solution $y=e^{-3 x}(7 x+6)$, which can be obtained by the method of Section 2.1. Do you notice anything special about the results? Explain.

The linear initial value problems in Exercises 14-19 can't be solved exactly in terms of known elementary functions. In each exercise, use Euler's method and the Euler semilinear methods with the indicated step sizes to find approximate values of the solution of the given initial value problem at 11 equally spaced points (including the endpoints) in the interval.
14. $\mathrm{C} y^{\prime}-2 y=\frac{1}{1+x^{2}}, \quad y(2)=2 ; \quad h=0.1,0.05,0.025$ on $[2,3]$
15. C $y^{\prime}+2 x y=x^{2}, \quad y(0)=3$ (Exercise 2.1.38); $h=0.2,0.1,0.05$ on $[0,2]$
16. C $y^{\prime}+\frac{1}{x} y=\frac{\sin x}{x^{2}}, \quad y(1)=2$; (Exercise 2.1.39); $h=0.2,0.1,0.05$ on $[1,3]$
17. $\mathrm{C} y^{\prime}+y=\frac{e^{-x} \tan x}{x}, \quad y(1)=0$; (Exercise 2.1.40); $\quad h=0.05,0.025,0.0125$ on $[1,1.5]$
18. $\mathrm{C} y^{\prime}+\frac{2 x}{1+x^{2}} y=\frac{e^{x}}{\left(1+x^{2}\right)^{2}}, \quad y(0)=1$; (Exercise 2.1.41); $\quad h=0.2,0.1,0.05$ on $[0,2]$
19. $\mathrm{C} x y^{\prime}+(x+1) y=e^{x^{2}}, \quad y(1)=2$; (Exercise 2.1.42); $h=0.05,0.025,0.0125$ on $[1,1.5]$

In Exercises 20-22, use Euler's method and the Euler semilinear method with the indicated step sizes to find approximate values of the solution of the given initial value problem at 11 equally spaced points (including the endpoints) in the interval.
20. C $y^{\prime}+3 y=x y^{2}(y+1), \quad y(0)=1 ; ~ h=0.1,0.05,0.025$ on $[0,1]$
21. C $y^{\prime}-4 y=\frac{x}{y^{2}(y+1)}, \quad y(0)=1 ; \quad h=0.1,0.05,0.025$ on $[0,1]$
22. $\mathrm{C} y^{\prime}+2 y=\frac{x^{2}}{1+y^{2}}, \quad y(2)=1 ; \quad h=0.1,0.05,0.025$ on $[2,3]$
23. Numerical Quadrature. The fundamental theorem of calculus says that if $f$ is continuous on a closed interval $[a, b]$ then it has an antiderivative $F$ such that $F^{\prime}(x)=f(x)$ on $[a, b]$ and

$$
\begin{equation*}
\int_{a}^{b} f(x) d x=F(b)-F(a) \tag{A}
\end{equation*}
$$

This solves the problem of evaluating a definite integral if the integrand $f$ has an antiderivative that can be found and evaluated easily. However, if $f$ doesn't have this property, (A) doesn't provide a useful way to evaluate the definite integral. In this case we must resort to approximate methods. There's a class of such methods called numerical quadrature, where the approximation takes the form

$$
\begin{equation*}
\int_{a}^{b} f(x) d x \approx \sum_{i=0}^{n} c_{i} f\left(x_{i}\right) \tag{B}
\end{equation*}
$$

where $a=x_{0}<x_{1}<\cdots<x_{n}=b$ are suitably chosen points and $c_{0}, c_{1}, \ldots, c_{n}$ are suitably chosen constants. We call (B) a quadrature formula.
(a) Derive the quadrature formula

$$
\begin{equation*}
\int_{a}^{b} f(x) d x \approx h \sum_{i=0}^{n-1} f(a+i h) \quad(\text { where } h=(b-a) / n) \tag{C}
\end{equation*}
$$

by applying Euler's method to the initial value problem

$$
y^{\prime}=f(x), \quad y(a)=0
$$

(b) The quadrature formula (C) is sometimes called the left rectangle rule. Draw a figure that justifies this terminology.
(c) L For several choices of $a, b$, and $A$, apply (C) to $f(x)=A$ with $n=10,20,40,80,160,320$. Compare your results with the exact answers and explain what you find.
(d) L For several choices of $a, b, A$, and $B$, apply (C) to $f(x)=A+B x$ with $n=10,20,40$, $80,160,320$. Compare your results with the exact answers and explain what you find.

### 3.2 THE IMPROVED EULER METHOD AND RELATED METHODS

In Section 3.1 we saw that the global truncation error of Euler's method is $O(h)$, which would seem to imply that we can achieve arbitrarily accurate results with Euler's method by simply choosing the step size sufficiently small. However, this isn't a good idea, for two reasons. First, after a certain point decreasing the step size will increase roundoff errors to the point where the accuracy will deteriorate rather than improve. The second and more important reason is that in most applications of numerical methods to an initial value problem

$$
\begin{equation*}
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0} \tag{3.2.1}
\end{equation*}
$$

the expensive part of the computation is the evaluation of $f$. Therefore we want methods that give good results for a given number of such evaluations. This is what motivates us to look for numerical methods better than Euler's.

To clarify this point, suppose we want to approximate the value of $e$ by applying Euler's method to the initial value problem

$$
y^{\prime}=y, \quad y(0)=1, \quad\left(\text { with solution } y=e^{x}\right)
$$

on $[0,1]$, with $h=1 / 12,1 / 24$, and $1 / 48$, respectively. Since each step in Euler's method requires one evaluation of $f$, the number of evaluations of $f$ in each of these attempts is $n=12,24$, and 48 , respectively. In each case we accept $y_{n}$ as an approximation to $e$. The second column of Table 3.2.1 shows the results. The first column of the table indicates the number of evaluations of $f$ required to obtain the approximation, and the last column contains the value of $e$ rounded to ten significant figures.

In this section we'll study the improved Euler method, which requires two evaluations of $f$ at each step. We've used this method with $h=1 / 6,1 / 12$, and $1 / 24$. The required number of evaluations of $f$ were 12,24 , and 48 , as in the three applications of Euler's method; however, you can see from the third column of Table 3.2.1 that the approximation to $e$ obtained by the improved Euler method with only 12 evaluations of $f$ is better than the approximation obtained by Euler's method with 48 evaluations.

In Section 3.1 we'll study the Runge-Kutta method, which requires four evaluations of $f$ at each step. We've used this method with $h=1 / 3,1 / 6$, and $1 / 12$. The required number of evaluations of $f$ were again 12, 24, and 48, as in the three applications of Euler's method and the improved Euler method; however, you can see from the fourth column of Table 3.2.1 that the approximation to $e$ obtained by the Runge-Kutta method with only 12 evaluations of $f$ is better than the approximation obtained by the improved Euler method with 48 evaluations.

Table 3.2.1. Approximations to $e$ obtained by three numerical methods.

| $n$ | Euler | Improved Euler | Runge-Kutta | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 12 | 2.613035290 | 2.707188994 | 2.718069764 | 2.718281828 |
| 24 | 2.663731258 | 2.715327371 | 2.718266612 | 2.718281828 |
| 48 | 2.690496599 | 2.717519565 | 2.718280809 | 2.718281828 |

The Improved Euler Method
The improved Euler method for solving the initial value problem (3.2.1) is based on approximating the integral curve of (3.2.1) at $\left(x_{i}, y\left(x_{i}\right)\right)$ by the line through $\left(x_{i}, y\left(x_{i}\right)\right)$ with slope

$$
m_{i}=\frac{f\left(x_{i}, y\left(x_{i}\right)\right)+f\left(x_{i+1}, y\left(x_{i+1}\right)\right)}{2} ;
$$

that is, $m_{i}$ is the average of the slopes of the tangents to the integral curve at the endpoints of $\left[x_{i}, x_{i+1}\right]$. The equation of the approximating line is therefore

$$
\begin{equation*}
y=y\left(x_{i}\right)+\frac{f\left(x_{i}, y\left(x_{i}\right)\right)+f\left(x_{i+1}, y\left(x_{i+1}\right)\right)}{2}\left(x-x_{i}\right) . \tag{3.2.2}
\end{equation*}
$$

Setting $x=x_{i+1}=x_{i}+h$ in (3.2.2) yields

$$
\begin{equation*}
y_{i+1}=y\left(x_{i}\right)+\frac{h}{2}\left(f\left(x_{i}, y\left(x_{i}\right)\right)+f\left(x_{i+1}, y\left(x_{i+1}\right)\right)\right) \tag{3.2.3}
\end{equation*}
$$

as an approximation to $y\left(x_{i+1}\right)$. As in our derivation of Euler's method, we replace $y\left(x_{i}\right)$ (unknown if $i>0$ ) by its approximate value $y_{i}$; then (3.2.3) becomes

$$
y_{i+1}=y_{i}+\frac{h}{2}\left(f\left(x_{i}, y_{i}\right)+f\left(x_{i+1}, y\left(x_{i+1}\right)\right) .\right.
$$

However, this still won't work, because we don't know $y\left(x_{i+1}\right)$, which appears on the right. We overcome this by replacing $y\left(x_{i+1}\right)$ by $y_{i}+h f\left(x_{i}, y_{i}\right)$, the value that the Euler method would assign to $y_{i+1}$. Thus, the improved Euler method starts with the known value $y\left(x_{0}\right)=y_{0}$ and computes $y_{1}, y_{2}, \ldots, y_{n}$ successively with the formula

$$
\begin{equation*}
y_{i+1}=y_{i}+\frac{h}{2}\left(f\left(x_{i}, y_{i}\right)+f\left(x_{i+1}, y_{i}+h f\left(x_{i}, y_{i}\right)\right)\right) . \tag{3.2.4}
\end{equation*}
$$

The computation indicated here can be conveniently organized as follows: given $y_{i}$, compute

$$
\begin{aligned}
k_{1 i} & =f\left(x_{i}, y_{i}\right), \\
k_{2 i} & =f\left(x_{i}+h, y_{i}+h k_{1 i}\right), \\
y_{i+1} & =y_{i}+\frac{h}{2}\left(k_{1 i}+k_{2 i}\right) .
\end{aligned}
$$

The improved Euler method requires two evaluations of $f(x, y)$ per step, while Euler's method requires only one. However, we'll see at the end of this section that if $f$ satisfies appropriate assumptions, the local truncation error with the improved Euler method is $O\left(h^{3}\right)$, rather than $O\left(h^{2}\right)$ as with Euler's method. Therefore the global truncation error with the improved Euler method is $O\left(h^{2}\right)$; however, we won't prove this.

We note that the magnitude of the local truncation error in the improved Euler method and other methods discussed in this section is determined by the third derivative $y^{\prime \prime \prime}$ of the solution of the initial value problem. Therefore the local truncation error will be larger where $\left|y^{\prime \prime \prime}\right|$ is large, or smaller where $\left|y^{\prime \prime \prime}\right|$ is small.

The next example, which deals with the initial value problem considered in Example 3.1.1, illustrates the computational procedure indicated in the improved Euler method.

Example 3.2.1 Use the improved Euler method with $h=0.1$ to find approximate values of the solution of the initial value problem

$$
\begin{equation*}
y^{\prime}+2 y=x^{3} e^{-2 x}, \quad y(0)=1 \tag{3.2.5}
\end{equation*}
$$

at $x=0.1,0.2,0.3$.

Solution As in Example 3.1.1, we rewrite (3.2.5) as

$$
y^{\prime}=-2 y+x^{3} e^{-2 x}, \quad y(0)=1,
$$

which is of the form (3.2.1), with

$$
f(x, y)=-2 y+x^{3} e^{-2 x}, x_{0}=0, \quad \text { and } y_{0}=1
$$

The improved Euler method yields

$$
\begin{aligned}
k_{10} & =f\left(x_{0}, y_{0}\right)=f(0,1)=-2, \\
k_{20} & =f\left(x_{1}, y_{0}+h k_{10}\right)=f(.1,1+(.1)(-2)) \\
& =f(.1, .8)=-2(.8)+(.1)^{3} e^{-.2}=-1.599181269, \\
y_{1} & =y_{0}+\frac{h}{2}\left(k_{10}+k_{20}\right), \\
& =1+(.05)(-2-1.599181269)=.820040937, \\
k_{11} & =f\left(x_{1}, y_{1}\right)=f(.1, .820040937)=-2(.820040937)+(.1)^{3} e^{-.2}=-1.639263142, \\
k_{21} & =f\left(x_{2}, y_{1}+h k_{11}\right)=f(.2, .820040937+.1(-1.639263142)), \\
& =f(.2, .656114622)=-2(.656114622)+(.2)^{3} e^{-.4}=-1.306866684, \\
y_{2} & =y_{1}+\frac{h}{2}\left(k_{11}+k_{21}\right), \\
& =.820040937+(.05)(-1.639263142-1.306866684)=.672734445, \\
k_{12} & =f\left(x_{2}, y_{2}\right)=f(.2, .672734445)=-2(.672734445)+(.2)^{3} e^{-.4}=-1.340106330, \\
k_{22} & =f\left(x_{3}, y_{2}+h k_{12}\right)=f(.3, .672734445+.1(-1.340106330)), \\
& =f(.3, .538723812)=-2(.538723812)+(.3)^{3} e^{-.6}=-1.062629710, \\
y_{3} & =y_{2}+\frac{h}{2}\left(k_{12}+k_{22}\right) \\
& =.672734445+(.05)(-1.340106330-1.062629710)=.552597643 .
\end{aligned}
$$

Example 3.2.2 Table 3.2.2 shows results of using the improved Euler method with step sizes $h=0.1$ and $h=0.05$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+2 y=x^{3} e^{-2 x}, \quad y(0)=1
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$. For comparison, it also shows the corresponding approximate values obtained with Euler's method in 3.1.2, and the values of the exact solution

$$
y=\frac{e^{-2 x}}{4}\left(x^{4}+4\right)
$$

The results obtained by the improved Euler method with $h=0.1$ are better than those obtained by Euler's method with $h=0.05$.

Table 3.2.2. Numerical solution of $y^{\prime}+2 y=x^{3} e^{-2 x}, y(0)=1$, by Euler's method and the improved Euler method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.1$ | $h=0.05$ | Exact |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 0.0 | 1.000000000 | 1.000000000 | 1.000000000 | 1.000000000 | 1.000000000 |  |  |  |  |
| 0.1 | 0.800000000 | 0.810005655 | 0.820040937 | 0.819050572 | 0.818751221 |  |  |  |  |
| 0.2 | 0.640081873 | 0.656266437 | 0.672734445 | 0.671086455 | 0.670588174 |  |  |  |  |
| 0.3 | 0.512601754 | 0.532290981 | 0.552597643 | 0.550543878 | 0.549922980 |  |  |  |  |
| 0.4 | 0.411563195 | 0.432887056 | 0.455160637 | 0.452890616 | 0.452204669 |  |  |  |  |
| 0.5 | 0.332126261 | 0.353785015 | 0.376681251 | 0.374335747 | 0.373627557 |  |  |  |  |
| 0.6 | 0.270299502 | 0.291404256 | 0.313970920 | 0.311652239 | 0.310952904 |  |  |  |  |
| 0.7 | 0.222745397 | 0.242707257 | 0.264287611 | 0.262067624 | 0.261398947 |  |  |  |  |
| 0.8 | 0.186654593 | 0.205105754 | 0.225267702 | 0.223194281 | 0.222570721 |  |  |  |  |
| 0.9 | 0.159660776 | 0.176396883 | 0.194879501 | 0.192981757 | 0.192412038 |  |  |  |  |
| 1.0 | 0.139778910 | 0.154715925 | 0.171388070 | 0.169680673 | 0.169169104 |  |  |  |  |
|  | Euler |  |  |  | Improved Euler |  |  |  | Exact |

Example 3.2.3 Table 3.2.3 shows analogous results for the nonlinear initial value problem

$$
y^{\prime}=-2 y^{2}+x y+x^{2}, y(0)=1 .
$$

We applied Euler's method to this problem in Example 3.1.3.

Table 3.2.3. Numerical solution of $y^{\prime}=-2 y^{2}+x y+x^{2}, y(0)=1$, by Euler's method and the improved Euler method.

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.1$ | $h=0.05$ | "Exact" |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 0.0 | 1.000000000 | 1.000000000 | 1.000000000 | 1.000000000 | 1.000000000 |  |  |  |  |
| 0.1 | 0.800000000 | 0.821375000 | 0.840500000 | 0.838288371 | 0.837584494 |  |  |  |  |
| 0.2 | 0.681000000 | 0.707795377 | 0.733430846 | 0.730556677 | 0.729641890 |  |  |  |  |
| 0.3 | 0.605867800 | 0.633776590 | 0.661600806 | 0.658552190 | 0.657580377 |  |  |  |  |
| 0.4 | 0.559628676 | 0.587454526 | 0.615961841 | 0.612884493 | 0.611901791 |  |  |  |  |
| 0.5 | 0.535376972 | 0.562906169 | 0.591634742 | 0.588558952 | 0.587575491 |  |  |  |  |
| 0.6 | 0.529820120 | 0.557143535 | 0.586006935 | 0.582927224 | 0.581942225 |  |  |  |  |
| 0.7 | 0.541467455 | 0.568716935 | 0.597712120 | 0.594618012 | 0.593629526 |  |  |  |  |
| 0.8 | 0.569732776 | 0.596951988 | 0.626008824 | 0.622898279 | 0.621907458 |  |  |  |  |
| 0.9 | 0.614392311 | 0.641457729 | 0.670351225 | 0.667237617 | 0.666250842 |  |  |  |  |
| 1.0 | 0.675192037 | 0.701764495 | 0.730069610 | 0.726985837 | 0.726015790 |  |  |  |  |
|  | Euler |  |  |  | Improved Euler |  |  |  | "Exact" |

Example 3.2.4 Use step sizes $h=0.2, h=0.1$, and $h=0.05$ to find approximate values of the solution of

$$
\begin{equation*}
y^{\prime}-2 x y=1, \quad y(0)=3 \tag{3.2.6}
\end{equation*}
$$

at $x=0,0.2,0.4,0.6, \ldots, 2.0$ by (a) the improved Euler method; (b) the improved Euler semilinear method. (We used Euler's method and the Euler semilinear method on this problem in 3.1.4.)
$\underline{\text { SOLUTION(a) Rewriting (3.2.6) as }}$

$$
y^{\prime}=1+2 x y, \quad y(0)=3
$$

and applying the improved Euler method with $f(x, y)=1+2 x y$ yields the results shown in Table 3.2.4.
SOLUTION(b) Since $y_{1}=e^{x^{2}}$ is a solution of the complementary equation $y^{\prime}-2 x y=0$, we can apply the improved Euler semilinear method to (3.2.6), with

$$
y=u e^{x^{2}} \quad \text { and } \quad u^{\prime}=e^{-x^{2}}, \quad u(0)=3
$$

The results listed in Table 3.2.5 are clearly better than those obtained by the improved Euler method.
Table 3.2.4. Numerical solution of $y^{\prime}-2 x y=1, y(0)=3$, by the improved Euler method.

| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| :---: | ---: | ---: | ---: | ---: |
| 0.0 | 3.000000000 | 3.000000000 | 3.000000000 | 3.000000000 |
| 0.2 | 3.328000000 | 3.328182400 | 3.327973600 | 3.327851973 |
| 0.4 | 3.964659200 | 3.966340117 | 3.966216690 | 3.966059348 |
| 0.6 | 5.057712497 | 5.065700515 | 5.066848381 | 5.067039535 |
| 0.8 | 6.900088156 | 6.928648973 | 6.934862367 | 6.936700945 |
| 1.0 | 10.065725534 | 10.154872547 | 10.177430736 | 10.184923955 |
| 1.2 | 15.708954420 | 15.970033261 | 16.041904862 | 16.067111677 |
| 1.4 | 26.244894192 | 26.991620960 | 27.210001715 | 27.289392347 |
| 1.6 | 46.958915746 | 49.096125524 | 49.754131060 | 50.000377775 |
| 1.8 | 89.982312641 | 96.200506218 | 98.210577385 | 98.982969504 |
| 2.0 | 184.563776288 | 203.151922739 | 209.464744495 | 211.954462214 |

Table 3.2.5. Numerical solution of $y^{\prime}-2 x y=1, y(0)=3$, by the improved Euler semilinear method.

| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| :---: | ---: | ---: | ---: | ---: |
| 0.0 | 3.000000000 | 3.000000000 | 3.000000000 | 3.000000000 |
| 0.2 | 3.326513400 | 3.327518315 | 3.327768620 | 3.327851973 |
| 0.4 | 3.963383070 | 3.965392084 | 3.965892644 | 3.966059348 |
| 0.6 | 5.063027290 | 5.066038774 | 5.066789487 | 5.067039535 |
| 0.8 | 6.931355329 | 6.935366847 | 6.936367564 | 6.936700945 |
| 1.0 | 10.178248417 | 10.183256733 | 10.184507253 | 10.184923955 |
| 1.2 | 16.059110511 | 16.065111599 | 16.066611672 | 16.067111677 |
| 1.4 | 27.280070674 | 27.287059732 | 27.288809058 | 27.289392347 |
| 1.6 | 49.989741531 | 49.997712997 | 49.999711226 | 50.000377775 |
| 1.8 | 98.971025420 | 98.979972988 | 98.982219722 | 98.982969504 |
| 2.0 | 211.941217796 | 211.951134436 | 211.953629228 | 211.954462214 |

A Family of Methods with $O\left(h^{3}\right)$ Local Truncation Error
We'll now derive a class of methods with $O\left(h^{3}\right)$ local truncation error for solving (3.2.1). For simplicity, we assume that $f, f_{x}, f_{y}, f_{x x}, f_{y y}$, and $f_{x y}$ are continuous and bounded for all $(x, y)$. This implies that if $y$ is the solution of (3.2.1 then $y^{\prime \prime}$ and $y^{\prime \prime \prime}$ are bounded (Exercise 31).

We begin by approximating the integral curve of (3.2.1) at $\left(x_{i}, y\left(x_{i}\right)\right)$ by the line through $\left(x_{i}, y\left(x_{i}\right)\right)$ with slope

$$
m_{i}=\sigma y^{\prime}\left(x_{i}\right)+\rho y^{\prime}\left(x_{i}+\theta h\right),
$$

where $\sigma, \rho$, and $\theta$ are constants that we'll soon specify; however, we insist at the outset that $0<\theta \leq 1$, so that

$$
x_{i}<x_{i}+\theta h \leq x_{i+1}
$$

The equation of the approximating line is

$$
\begin{align*}
y & =y\left(x_{i}\right)+m_{i}\left(x-x_{i}\right)  \tag{3.2.7}\\
& =y\left(x_{i}\right)+\left[\sigma y^{\prime}\left(x_{i}\right)+\rho y^{\prime}\left(x_{i}+\theta h\right)\right]\left(x-x_{i}\right) .
\end{align*}
$$

Setting $x=x_{i+1}=x_{i}+h$ in (3.2.7) yields

$$
\hat{y}_{i+1}=y\left(x_{i}\right)+h\left[\sigma y^{\prime}\left(x_{i}\right)+\rho y^{\prime}\left(x_{i}+\theta h\right)\right]
$$

as an approximation to $y\left(x_{i+1}\right)$.
To determine $\sigma, \rho$, and $\theta$ so that the error

$$
\begin{align*}
E_{i} & =y\left(x_{i+1}\right)-\hat{y}_{i+1} \\
& =y\left(x_{i+1}\right)-y\left(x_{i}\right)-h\left[\sigma y^{\prime}\left(x_{i}\right)+\rho y^{\prime}\left(x_{i}+\theta h\right)\right] \tag{3.2.8}
\end{align*}
$$

in this approximation is $O\left(h^{3}\right)$, we begin by recalling from Taylor's theorem that

$$
y\left(x_{i+1}\right)=y\left(x_{i}\right)+h y^{\prime}\left(x_{i}\right)+\frac{h^{2}}{2} y^{\prime \prime}\left(x_{i}\right)+\frac{h^{3}}{6} y^{\prime \prime \prime}\left(\hat{x}_{i}\right)
$$

where $\hat{x}_{i}$ is in $\left(x_{i}, x_{i+1}\right)$. Since $y^{\prime \prime \prime}$ is bounded this implies that

$$
y\left(x_{i+1}\right)-y\left(x_{i}\right)-h y^{\prime}\left(x_{i}\right)-\frac{h^{2}}{2} y^{\prime \prime}\left(x_{i}\right)=O\left(h^{3}\right) .
$$

Comparing this with (3.2.8) shows that $E_{i}=O\left(h^{3}\right)$ if

$$
\begin{equation*}
\sigma y^{\prime}\left(x_{i}\right)+\rho y^{\prime}\left(x_{i}+\theta h\right)=y^{\prime}\left(x_{i}\right)+\frac{h}{2} y^{\prime \prime}\left(x_{i}\right)+O\left(h^{2}\right) . \tag{3.2.9}
\end{equation*}
$$

However, applying Taylor's theorem to $y^{\prime}$ shows that

$$
y^{\prime}\left(x_{i}+\theta h\right)=y^{\prime}\left(x_{i}\right)+\theta h y^{\prime \prime}\left(x_{i}\right)+\frac{(\theta h)^{2}}{2} y^{\prime \prime \prime}\left(\bar{x}_{i}\right)
$$

where $\bar{x}_{i}$ is in $\left(x_{i}, x_{i}+\theta h\right)$. Since $y^{\prime \prime \prime}$ is bounded, this implies that

$$
y^{\prime}\left(x_{i}+\theta h\right)=y^{\prime}\left(x_{i}\right)+\theta h y^{\prime \prime}\left(x_{i}\right)+O\left(h^{2}\right) .
$$

Substituting this into (3.2.9) and noting that the sum of two $O\left(h^{2}\right)$ terms is again $O\left(h^{2}\right)$ shows that $E_{i}=O\left(h^{3}\right)$ if

$$
(\sigma+\rho) y^{\prime}\left(x_{i}\right)+\rho \theta h y^{\prime \prime}\left(x_{i}\right)=y^{\prime}\left(x_{i}\right)+\frac{h}{2} y^{\prime \prime}\left(x_{i}\right)
$$

which is true if

$$
\begin{equation*}
\sigma+\rho=1 \quad \text { and } \quad \rho \theta=\frac{1}{2} \tag{3.2.10}
\end{equation*}
$$

Since $y^{\prime}=f(x, y)$, we can now conclude from (3.2.8) that

$$
\begin{equation*}
y\left(x_{i+1}\right)=y\left(x_{i}\right)+h\left[\sigma f\left(x_{i}, y_{i}\right)+\rho f\left(x_{i}+\theta h, y\left(x_{i}+\theta h\right)\right)\right]+O\left(h^{3}\right) \tag{3.2.11}
\end{equation*}
$$

if $\sigma, \rho$, and $\theta$ satisfy (3.2.10). However, this formula would not be useful even if we knew $y\left(x_{i}\right)$ exactly (as we would for $i=0$ ), since we still wouldn't know $y\left(x_{i}+\theta h\right)$ exactly. To overcome this difficulty, we again use Taylor's theorem to write

$$
y\left(x_{i}+\theta h\right)=y\left(x_{i}\right)+\theta h y^{\prime}\left(x_{i}\right)+\frac{h^{2}}{2} y^{\prime \prime}\left(\tilde{x}_{i}\right),
$$

where $\tilde{x}_{i}$ is in $\left(x_{i}, x_{i}+\theta h\right)$. Since $y^{\prime}\left(x_{i}\right)=f\left(x_{i}, y\left(x_{i}\right)\right)$ and $y^{\prime \prime}$ is bounded, this implies that

$$
\begin{equation*}
\left|y\left(x_{i}+\theta h\right)-y\left(x_{i}\right)-\theta h f\left(x_{i}, y\left(x_{i}\right)\right)\right| \leq K h^{2} \tag{3.2.12}
\end{equation*}
$$

for some constant $K$. Since $f_{y}$ is bounded, the mean value theorem implies that

$$
\left|f\left(x_{i}+\theta h, u\right)-f\left(x_{i}+\theta h, v\right)\right| \leq M|u-v|
$$

for some constant $M$. Letting

$$
u=y\left(x_{i}+\theta h\right) \quad \text { and } \quad v=y\left(x_{i}\right)+\theta h f\left(x_{i}, y\left(x_{i}\right)\right)
$$

and recalling (3.2.12) shows that

$$
f\left(x_{i}+\theta h, y\left(x_{i}+\theta h\right)\right)=f\left(x_{i}+\theta h, y\left(x_{i}\right)+\theta h f\left(x_{i}, y\left(x_{i}\right)\right)\right)+O\left(h^{2}\right) .
$$

Substituting this into (3.2.11) yields

$$
\begin{aligned}
y\left(x_{i+1}\right)= & y\left(x_{i}\right)+h\left[\sigma f\left(x_{i}, y\left(x_{i}\right)\right)+\right. \\
& \left.\rho f\left(x_{i}+\theta h, y\left(x_{i}\right)+\theta h f\left(x_{i}, y\left(x_{i}\right)\right)\right)\right]+O\left(h^{3}\right) .
\end{aligned}
$$

This implies that the formula

$$
y_{i+1}=y_{i}+h\left[\sigma f\left(x_{i}, y_{i}\right)+\rho f\left(x_{i}+\theta h, y_{i}+\theta h f\left(x_{i}, y_{i}\right)\right)\right]
$$

has $O\left(h^{3}\right)$ local truncation error if $\sigma, \rho$, and $\theta$ satisfy (3.2.10). Substituting $\sigma=1-\rho$ and $\theta=1 / 2 \rho$ here yields

$$
\begin{equation*}
y_{i+1}=y_{i}+h\left[(1-\rho) f\left(x_{i}, y_{i}\right)+\rho f\left(x_{i}+\frac{h}{2 \rho}, y_{i}+\frac{h}{2 \rho} f\left(x_{i}, y_{i}\right)\right)\right] . \tag{3.2.13}
\end{equation*}
$$

The computation indicated here can be conveniently organized as follows: given $y_{i}$, compute

$$
\begin{aligned}
k_{1 i} & =f\left(x_{i}, y_{i}\right) \\
k_{2 i} & =f\left(x_{i}+\frac{h}{2 \rho}, y_{i}+\frac{h}{2 \rho} k_{1 i}\right), \\
y_{i+1} & =y_{i}+h\left[(1-\rho) k_{1 i}+\rho k_{2 i}\right] .
\end{aligned}
$$

Consistent with our requirement that $0<\theta<1$, we require that $\rho \geq 1 / 2$. Letting $\rho=1 / 2$ in (3.2.13) yields the improved Euler method (3.2.4). Letting $\rho=3 / 4$ yields Heun's method,

$$
y_{i+1}=y_{i}+h\left[\frac{1}{4} f\left(x_{i}, y_{i}\right)+\frac{3}{4} f\left(x_{i}+\frac{2}{3} h, y_{i}+\frac{2}{3} h f\left(x_{i}, y_{i}\right)\right)\right],
$$

which can be organized as

$$
\begin{aligned}
k_{1 i} & =f\left(x_{i}, y_{i}\right) \\
k_{2 i} & =f\left(x_{i}+\frac{2 h}{3}, y_{i}+\frac{2 h}{3} k_{1 i}\right), \\
y_{i+1} & =y_{i}+\frac{h}{4}\left(k_{1 i}+3 k_{2 i}\right)
\end{aligned}
$$

Letting $\rho=1$ yields the midpoint method,

$$
y_{i+1}=y_{i}+h f\left(x_{i}+\frac{h}{2}, y_{i}+\frac{h}{2} f\left(x_{i}, y_{i}\right)\right),
$$

which can be organized as

$$
\begin{aligned}
k_{1 i} & =f\left(x_{i}, y_{i}\right) \\
k_{2 i} & =f\left(x_{i}+\frac{h}{2}, y_{i}+\frac{h}{2} k_{1 i}\right), \\
y_{i+1} & =y_{i}+h k_{2 i} .
\end{aligned}
$$

Examples involving the midpoint method and Heun's method are given in Exercises 23-30.

### 3.2 Exercises

Most of the following numerical exercises involve initial value problems considered in the exercises in Section 3.1. You'll find it instructive to compare the results that you obtain here with the corresponding results that you obtained in Section 3.1.

In Exercises 1-5 use the improved Euler method to find approximate values of the solution of the given initial value problem at the points $x_{i}=x_{0}+i h$, where $x_{0}$ is the point where the initial condition is imposed and $i=1,2,3$.

1. C $y^{\prime}=2 x^{2}+3 y^{2}-2, \quad y(2)=1 ; \quad h=0.05$
2. $\quad \mathrm{C} y^{\prime}=y+\sqrt{x^{2}+y^{2}}, \quad y(0)=1 ; \quad h=0.1$
3. $\mathrm{C} y^{\prime}+3 y=x^{2}-3 x y+y^{2}, \quad y(0)=2 ; \quad h=0.05$
4. $\mathrm{C} y^{\prime}=\frac{1+x}{1-y^{2}}, \quad y(2)=3 ; \quad h=0.1$
5. C $y^{\prime}+x^{2} y=\sin x y, \quad y(1)=\pi ; \quad h=0.2$
6. $\quad$ C Use the improved Euler method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+3 y=7 e^{4 x}, \quad y(0)=2
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$. Compare these approximate values with the values of the exact solution $y=e^{4 x}+e^{-3 x}$, which can be obtained by the method of Section 2.1. Present your results in a table like Table 3.2.2.
7. C Use the improved Euler method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+\frac{2}{x} y=\frac{3}{x^{3}}+1, \quad y(1)=1
$$

at $x=1.0,1.1,1.2,1.3, \ldots, 2.0$. Compare these approximate values with the values of the exact solution

$$
y=\frac{1}{3 x^{2}}\left(9 \ln x+x^{3}+2\right)
$$

which can be obtained by the method of Section 2.1. Present your results in a table like Table 3.2.2.
8. C Use the improved Euler method with step sizes $h=0.05, h=0.025$, and $h=0.0125$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}=\frac{y^{2}+x y-x^{2}}{x^{2}}, \quad y(1)=2,
$$

at $x=1.0,1.05,1.10,1.15, \ldots, 1.5$. Compare these approximate values with the values of the exact solution

$$
y=\frac{x\left(1+x^{2} / 3\right)}{1-x^{2} / 3}
$$

obtained in Example ??. Present your results in a table like Table 3.2.2.
9. C In Example 3.2.2 it was shown that

$$
y^{5}+y=x^{2}+x-4
$$

is an implicit solution of the initial value problem

$$
\begin{equation*}
y^{\prime}=\frac{2 x+1}{5 y^{4}+1}, \quad y(2)=1 \tag{A}
\end{equation*}
$$

Use the improved Euler method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of (A) at $x=2.0,2.1,2.2,2.3, \ldots, 3.0$. Present your results in tabular form. To check the error in these approximate values, construct another table of values of the residual

$$
R(x, y)=y^{5}+y-x^{2}-x+4
$$

for each value of $(x, y)$ appearing in the first table.
10. C You can see from Example 2.5 .1 that

$$
x^{4} y^{3}+x^{2} y^{5}+2 x y=4
$$

is an implicit solution of the initial value problem

$$
\begin{equation*}
y^{\prime}=-\frac{4 x^{3} y^{3}+2 x y^{5}+2 y}{3 x^{4} y^{2}+5 x^{2} y^{4}+2 x}, \quad y(1)=1 \tag{A}
\end{equation*}
$$

Use the improved Euler method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of (A) at $x=1.0,1.14,1.2,1.3, \ldots, 2.0$. Present your results in tabular form. To check the error in these approximate values, construct another table of values of the residual

$$
R(x, y)=x^{4} y^{3}+x^{2} y^{5}+2 x y-4
$$

for each value of $(x, y)$ appearing in the first table.
11. C Use the improved Euler method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
\left(3 y^{2}+4 y\right) y^{\prime}+2 x+\cos x=0, \quad y(0)=1 \text { (Exercise 2.2.13) }
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$.
12. C Use the improved Euler method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+\frac{(y+1)(y-1)(y-2)}{x+1}=0, \quad y(1)=0(\text { Exercise 2.2.14 })
$$

at $x=1.0,1.1,1.2,1.3, \ldots, 2.0$.
13. C Use the improved Euler method and the improved Euler semilinear method with step sizes $h=0.1, h=0.05$, and $h=0.025$ to find approximate values of the solution of the initial value problem

$$
y^{\prime}+3 y=e^{-3 x}(1-2 x), \quad y(0)=2
$$

at $x=0,0.1,0.2,0.3, \ldots, 1.0$. Compare these approximate values with the values of the exact solution $y=e^{-3 x}\left(2+x-x^{2}\right)$, which can be obtained by the method of Section 2.1. Do you notice anything special about the results? Explain.

The linear initial value problems in Exercises 14-19 can't be solved exactly in terms of known elementary functions. In each exercise use the improved Euler and improved Euler semilinear methods with the indicated step sizes to find approximate values of the solution of the given initial value problem at 11 equally spaced points (including the endpoints) in the interval.
14. $\mathrm{C} y^{\prime}-2 y=\frac{1}{1+x^{2}}, \quad y(2)=2 ; \quad h=0.1,0.05,0.025$ on $[2,3]$
15. C $y^{\prime}+2 x y=x^{2}, \quad y(0)=3 ; \quad h=0.2,0.1,0.05$ on $[0,2] \quad$ (Exercise 2.1.38)
16. $\mathrm{C} y^{\prime}+\frac{1}{x} y=\frac{\sin x}{x^{2}}, \quad y(1)=2, \quad h=0.2,0.1,0.05$ on $[1,3] \quad$ (Exercise 2.139)
17. $\mathrm{C} y^{\prime}+y=\frac{e^{-x} \tan x}{x}, \quad y(1)=0 ; \quad h=0.05,0.025,0.0125$ on $[1,1.5]$ (Exercise 2.1.40),
18. $\mathrm{C} y^{\prime}+\frac{2 x}{1+x^{2}} y=\frac{e^{x}}{\left(1+x^{2}\right)^{2}}, \quad y(0)=1 ; h=0.2,0.1,0.05$ on $[0,2] \quad$ (Exercise 2.1.41)
19. $\mathrm{C} x y^{\prime}+(x+1) y=e^{x^{2}}, \quad y(1)=2 ; \quad h=0.05,0.025,0.0125$ on $[1,1.5]$ (Exercise 2.1.42)

In Exercises 20-22 use the improved Euler method and the improved Euler semilinear method with the indicated step sizes to find approximate values of the solution of the given initial value problem at 11 equally spaced points (including the endpoints) in the interval.
20. C $y^{\prime}+3 y=x y^{2}(y+1), \quad y(0)=1 ; \quad h=0.1,0.05,0.025$ on $[0,1]$
21. $\mathrm{C} y^{\prime}-4 y=\frac{x}{y^{2}(y+1)}, \quad y(0)=1 ; \quad h=0.1,0.05,0.025$ on $[0,1]$
22. $\mathrm{C} y^{\prime}+2 y=\frac{x^{2}}{1+y^{2}}, \quad y(2)=1 ; \quad h=0.1,0.05,0.025$ on $[2,3]$
23. C Do Exercise 7 with "improved Euler method" replaced by "midpoint method."
24. C Do Exercise 7 with "improved Euler method" replaced by "Heun's method."
25. C Do Exercise 8 with "improved Euler method" replaced by "midpoint method."
26. C Do Exercise 8 with "improved Euler method" replaced by "Heun's method."
27. C Do Exercise 11 with "improved Euler method" replaced by "midpoint method."
28. C Do Exercise 11 with "improved Euler method" replaced by "Heun's method."
29. C Do Exercise 12 with "improved Euler method" replaced by "midpoint method."
30. C Do Exercise 12 with "improved Euler method" replaced by "Heun's method."
31. Show that if $f, f_{x}, f_{y}, f_{x x}, f_{y y}$, and $f_{x y}$ are continuous and bounded for all $(x, y)$ and $y$ is the solution of the initial value problem

$$
y^{\prime}=f(x, y), \quad y\left(x_{0}\right)=y_{0}
$$

then $y^{\prime \prime}$ and $y^{\prime \prime \prime}$ are bounded.
32. Numerical Quadrature (see Exercise 3.1.23).
(a) Derive the quadrature formula

$$
\begin{equation*}
\int_{a}^{b} f(x) d x \approx .5 h(f(a)+f(b))+h \sum_{i=1}^{n-1} f(a+i h) \quad(\text { where } h=(b-a) / n) \tag{A}
\end{equation*}
$$

by applying the improved Euler method to the initial value problem

$$
y^{\prime}=f(x), \quad y(a)=0
$$

(b) The quadrature formula (A) is called the trapezoid rule. Draw a figure that justifies this terminology.
(c) L For several choices of $a, b, A$, and $B$, apply (A) to $f(x)=A+B x$, with $n=$ $10,20,40,80,160,320$. Compare your results with the exact answers and explain what you find.
(d) L For several choices of $a, b, A, B$, and $C$, apply (A) to $f(x)=A+B x+C x^{2}$, with $n=10,20,40,80,160,320$. Compare your results with the exact answers and explain what you find.

## CHAPTER 4 Applications of First Order Equations

IN THIS CHAPTER we consider applications of first order differential equations.
SECTION 4.1 begins with a discussion of exponential growth and decay, which you have probably already seen in calculus. We consider applications to radioactive decay, carbon dating, and compound interest. We also consider more complicated problems where the rate of change of a quantity is in part proportional to the magnitude of the quantity, but is also influenced by other other factors for example, a radioactive susbstance is manufactured at a certain rate, but decays at a rate proportional to its mass, or a saver makes regular deposits in a savings account that draws compound interest.

SECTION 4.2 deals with applications of Newton's law of cooling and with mixing problems.
SECTION 4.3 discusses applications to elementary mechanics involving Newton's second law of motion. The problems considered include motion under the influence of gravity in a resistive medium, and determining the initial velocity required to launch a satellite.

SECTION 4.4 deals with methods for dealing with a type of second order equation that often arises in applications of Newton's second law of motion, by reformulating it as first order equation with a different independent variable. Although the method doesn't usually lead to an explicit solution of the given equation, it does provide valuable insights into the behavior of the solutions.
SECTION 4.5 deals with applications of differential equations to curves.

### 4.2 COOLING AND MIXING

## Newton's Law of Cooling

Newton's law of cooling states that if an object with temperature $T(t)$ at time $t$ is in a medium with temperature $T_{m}(t)$, the rate of change of $T$ at time $t$ is proportional to $T(t)-T_{m}(t)$; thus, $T$ satisfies a differential equation of the form

$$
\begin{equation*}
T^{\prime}=-k\left(T-T_{m}\right) . \tag{4.2.1}
\end{equation*}
$$

Here $k>0$, since the temperature of the object must decrease if $T>T_{m}$, or increase if $T<T_{m}$. We'll call $k$ the temperature decay constant of the medium.

For simplicity, in this section we'll assume that the medium is maintained at a constant temperature $T_{m}$. This is another example of building a simple mathematical model for a physical phenomenon. Like most mathematical models it has its limitations. For example, it's reasonable to assume that the temperature of a room remains approximately constant if the cooling object is a cup of coffee, but perhaps not if it's a huge cauldron of molten metal. (For more on this see Exercise 17.)

To solve (4.2.1), we rewrite it as

$$
T^{\prime}+k T=k T_{m}
$$

Since $e^{-k t}$ is a solution of the complementary equation, the solutions of this equation are of the form $T=u e^{-k t}$, where $u^{\prime} e^{-k t}=k T_{m}$, so $u^{\prime}=k T_{m} e^{k t}$. Hence,

$$
u=T_{m} e^{k t}+c
$$

so

$$
T=u e^{-k t}=T_{m}+c e^{-k t}
$$

If $T(0)=T_{0}$, setting $t=0$ here yields $c=T_{0}-T_{m}$, so

$$
\begin{equation*}
T=T_{m}+\left(T_{0}-T_{m}\right) e^{-k t} \tag{4.2.2}
\end{equation*}
$$

Note that $T-T_{m}$ decays exponentially, with decay constant $k$.
Example 4.2.1 A ceramic insulator is baked at $400^{\circ} \mathrm{C}$ and cooled in a room in which the temperature is $25^{\circ} \mathrm{C}$. After 4 minutes the temperature of the insulator is $200^{\circ} \mathrm{C}$. What is its temperature after 8 minutes?

Solution Here $T_{0}=400$ and $T_{m}=25$, so (4.2.2) becomes

$$
\begin{equation*}
T=25+375 e^{-k t} \tag{4.2.3}
\end{equation*}
$$

We determine $k$ from the stated condition that $T(4)=200$; that is,

$$
200=25+375 e^{-4 k}
$$

hence,

$$
e^{-4 k}=\frac{175}{375}=\frac{7}{15} .
$$

Taking logarithms and solving for $k$ yields

$$
k=-\frac{1}{4} \ln \frac{7}{15}=\frac{1}{4} \ln \frac{15}{7} .
$$

Substituting this into (4.2.3) yields

$$
T=25+375 e^{-\frac{t}{4} \ln \frac{15}{7}}
$$

(Figure 4.2.1). Therefore the temperature of the insulator after 8 minutes is

$$
\begin{aligned}
T(8) & =25+375 e^{-2 \ln \frac{15}{7}} \\
& =25+375\left(\frac{7}{15}\right)^{2} \approx 107^{\circ} \mathrm{C}
\end{aligned}
$$



Figure 4.2.1 $T=25+375 e^{-(t / 4) \ln 15 / 7}$

Example 4.2.2 An object with temperature $72^{\circ} \mathrm{F}$ is placed outside, where the temperature is $-20^{\circ} \mathrm{F}$. At $11: 05$ the temperature of the object is $60^{\circ} \mathrm{F}$ and at $11: 07$ its temperature is $50^{\circ} \mathrm{F}$. At what time was the object placed outside?

Solution Let $T(t)$ be the temperature of the object at time $t$. For convenience, we choose the origin $t_{0}=0$ of the time scale to be 11:05 so that $T_{0}=60$. We must determine the time $\tau$ when $T(\tau)=72$. Substituting $T_{0}=60$ and $T_{m}=-20$ into (4.2.2) yields

$$
T=-20+(60-(-20)) e^{-k t}
$$

or

$$
\begin{equation*}
T=-20+80 e^{-k t} \tag{4.2.4}
\end{equation*}
$$

We obtain $k$ from the stated condition that the temperature of the object is $50^{\circ} \mathrm{F}$ at 11:07. Since 11:07 is $t=2$ on our time scale, we can determine $k$ by substituting $T=50$ and $t=2$ into (4.2.4) to obtain

$$
50=-20+80 e^{-2 k}
$$

(Figure 4.2.2); hence,

$$
e^{-2 k}=\frac{70}{80}=\frac{7}{8}
$$

Taking logarithms and solving for $k$ yields

$$
k=-\frac{1}{2} \ln \frac{7}{8}=\frac{1}{2} \ln \frac{8}{7} .
$$

Substituting this into (4.2.4) yields

$$
T=-20+80 e^{-\frac{t}{2} \ln \frac{8}{7}},
$$

and the condition $T(\tau)=72$ implies that

$$
72=-20+80 e^{-\frac{7}{2} \ln \frac{8}{7}} ;
$$

hence,

$$
e^{-\frac{\tau}{2} \ln \frac{8}{7}}=\frac{92}{80}=\frac{23}{20} .
$$



Figure 4.2.2 $T=-20+80 e^{-\frac{t}{2} \ln \frac{8}{7}}$

Taking logarithms and solving for $\tau$ yields

$$
\tau=-\frac{2 \ln \frac{23}{20}}{\ln \frac{8}{7}} \approx-2.09 \mathrm{~min}
$$

Therefore the object was placed outside about 2 minutes and 5 seconds before 11:05; that is, at 11:02:55.

## Mixing Problems

In the next two examples a saltwater solution with a given concentration (weight of salt per unit volume of solution) is added at a specified rate to a tank that initially contains saltwater with a different concentration. The problem is to determine the quantity of salt in the tank as a function of time. This is an example of a mixing problem. To construct a tractable mathematical model for mixing problems we assume in our examples (and most exercises) that the mixture is stirred instantly so that the salt is always uniformly distributed throughout the mixture. Exercises 22 and 23 deal with situations where this isn't so, but the distribution of salt becomes approximately uniform as $t \rightarrow \infty$.

Example 4.2.3 A tank initially contains 40 pounds of salt dissolved in 600 gallons of water. Starting at $t_{0}=0$, water that contains $1 / 2$ pound of salt per gallon is poured into the tank at the rate of $4 \mathrm{gal} / \mathrm{min}$ and the mixture is drained from the tank at the same rate (Figure 4.2.3).
(a) Find a differential equation for the quantity $Q(t)$ of salt in the tank at time $t>0$, and solve the equation to determine $Q(t)$.
(b) Find $\lim _{t \rightarrow \infty} Q(t)$.
$\underline{\operatorname{Solution}(a)}$ To find a differential equation for $Q$, we must use the given information to derive an expression for $Q^{\prime}$. But $Q^{\prime}$ is the rate of change of the quantity of salt in the tank changes with respect to time; thus, if rate in denotes the rate at which salt enters the tank and rate out denotes the rate by which it leaves, then

$$
\begin{equation*}
Q^{\prime}=\text { rate in - rate out. } \tag{4.2.5}
\end{equation*}
$$

The rate in is

$$
\left(\frac{1}{2} \mathrm{lb} / \mathrm{gal}\right) \times(4 \mathrm{gal} / \mathrm{min})=2 \mathrm{lb} / \mathrm{min} .
$$

Determining the rate out requires a little more thought. We're removing 4 gallons of the mixture per minute, and there are always 600 gallons in the tank; that is, we're removing $1 / 150$ of the mixture per


Figure 4.2.3 A mixing problem
minute. Since the salt is evenly distributed in the mixture, we are also removing $1 / 150$ of the salt per minute. Therefore, if there are $Q(t)$ pounds of salt in the tank at time $t$, the rate out at any time $t$ is $Q(t) / 150$. Alternatively, we can arrive at this conclusion by arguing that

$$
\begin{aligned}
\text { rate out } & =(\text { concentration }) \times(\text { rate of flow out }) \\
& =(\mathrm{lb} / \mathrm{gal}) \times(\mathrm{gal} / \mathrm{min}) \\
& =\frac{Q(t)}{600} \times 4=\frac{Q(t)}{150} .
\end{aligned}
$$

We can now write (4.2.5) as

$$
Q^{\prime}=2-\frac{Q}{150} .
$$

This first order equation can be rewritten as

$$
Q^{\prime}+\frac{Q}{150}=2
$$

Since $e^{-t / 150}$ is a solution of the complementary equation, the solutions of this equation are of the form $Q=u e^{-t / 150}$, where $u^{\prime} e^{-t / 150}=2$, so $u^{\prime}=2 e^{t / 150}$. Hence,

$$
u=300 e^{t / 150}+c,
$$

so

$$
\begin{equation*}
Q=u e^{-t / 150}=300+c e^{-t / 150} \tag{4.2.6}
\end{equation*}
$$

(Figure 4.2.4). Since $Q(0)=40, c=-260$; therefore,

$$
Q=300-260 e^{-t / 150}
$$

$\underline{\text { SOLUTION(b) }}$ From (4.2.6), we see that that $\lim _{t \rightarrow \infty} Q(t)=300$ for any value of $Q(0)$. This is intuitively reasonable, since the incoming solution contains $1 / 2$ pound of salt per gallon and there are always 600 gallons of water in the tank.

Example 4.2.4 A 500-liter tank initially contains 10 g of salt dissolved in 200 liters of water. Starting at $t_{0}=0$, water that contains $1 / 4 \mathrm{~g}$ of salt per liter is poured into the tank at the rate of 4 liters $/ \mathrm{min}$ and the mixture is drained from the tank at the rate of 2 liters $/ \mathrm{min}$ (Figure 4.2.5). Find a differential equation for the quantity $Q(t)$ of salt in the tank at time $t$ prior to the time when the tank overflows and find the concentration $K(t)$ (g/liter) of salt in the tank at any such time.


Figure 4.2.4 $Q=300-260 e^{-t / 150}$


Figure 4.2.5 Another mixing problem

Solution We first determine the amount $W(t)$ of solution in the tank at any time $t$ prior to overflow. Since $W(0)=200$ and we're adding 4 liters/min while removing only 2 liters/min, there's a net gain of 2 liters/min in the tank; therefore,

$$
W(t)=2 t+200
$$

Since $W(150)=500$ liters (capacity of the tank), this formula is valid for $0 \leq t \leq 150$.
Now let $Q(t)$ be the number of grams of salt in the tank at time $t$, where $0 \leq t \leq 150$. As in Example 4.2.3,

$$
\begin{equation*}
Q^{\prime}=\text { rate in }- \text { rate out. } \tag{4.2.7}
\end{equation*}
$$

The rate in is

$$
\begin{equation*}
\left(\frac{1}{4} \mathrm{~g} / \text { liter }\right) \times(4 \text { liters } / \mathrm{min})=1 \mathrm{~g} / \mathrm{min} \tag{4.2.8}
\end{equation*}
$$

To determine the rate out, we observe that since the mixture is being removed from the tank at the constant rate of 2 liters $/ \mathrm{min}$ and there are $2 t+200$ liters in the tank at time $t$, the fraction of the mixture being removed per minute at time $t$ is

$$
\frac{2}{2 t+200}=\frac{1}{t+100}
$$

We're removing this same fraction of the salt per minute. Therefore, since there are $Q(t)$ grams of salt in the tank at time $t$,

$$
\begin{equation*}
\text { rate out }=\frac{Q(t)}{t+100} \tag{4.2.9}
\end{equation*}
$$

Alternatively, we can arrive at this conclusion by arguing that

$$
\begin{aligned}
\text { rate out } & =(\text { concentration }) \times(\text { rate of flow out })=(\mathrm{g} / \text { liter }) \times(\text { liters } / \mathrm{min}) \\
& =\frac{Q(t)}{2 t+200} \times 2=\frac{Q(t)}{t+100}
\end{aligned}
$$

Substituting (4.2.8) and (4.2.9) into (4.2.7) yields

$$
\begin{equation*}
Q^{\prime}=1-\frac{Q}{t+100}, \quad \text { so } \quad Q^{\prime}+\frac{1}{t+100} Q=1 \tag{4.2.10}
\end{equation*}
$$

By separation of variables, $1 /(t+100)$ is a solution of the complementary equation, so the solutions of (4.2.10) are of the form

$$
Q=\frac{u}{t+100}, \quad \text { where } \quad \frac{u^{\prime}}{t+100}=1, \quad \text { so } \quad u^{\prime}=t+100
$$

Hence,

$$
\begin{equation*}
u=\frac{(t+100)^{2}}{2}+c \tag{4.2.11}
\end{equation*}
$$

Since $Q(0)=10$ and $u=(t+100) Q$, (4.2.11) implies that

$$
(100)(10)=\frac{(100)^{2}}{2}+c
$$

so

$$
c=100(10)-\frac{(100)^{2}}{2}=-4000
$$

and therefore

$$
u=\frac{(t+100)^{2}}{2}-4000
$$

Hence,

$$
Q=\frac{u}{t+200}=\frac{t+100}{2}-\frac{4000}{t+100}
$$

Now let $K(t)$ be the concentration of salt at time $t$. Then

$$
K(t)=\frac{1}{4}-\frac{2000}{(t+100)^{2}}
$$

102 Chapter 4 Applications of First Order Equations
(Figure 4.2.6).


Figure 4.2.6 $K(t)=\frac{1}{4}-\frac{2000}{(t+100)^{2}}$

### 4.2 Exercises

1. A thermometer is moved from a room where the temperature is $70^{\circ} \mathrm{F}$ to a freezer where the temperature is $12^{\circ} \mathrm{F}$. After 30 seconds the thermometer reads $40^{\circ} \mathrm{F}$. What does it read after 2 minutes?
2. A fluid initially at $100^{\circ} \mathrm{C}$ is placed outside on a day when the temperature is $-10^{\circ} \mathrm{C}$, and the temperature of the fluid drops $20^{\circ} \mathrm{C}$ in one minute. Find the temperature $T(t)$ of the fluid for $t>0$.
3. At $12: 00 \mathrm{PM}$ a thermometer reading $10^{\circ} \mathrm{F}$ is placed in a room where the temperature is $70^{\circ} \mathrm{F}$. It reads $56^{\circ}$ when it's placed outside, where the temperature is $5^{\circ} \mathrm{F}$, at $12: 03$. What does it read at 12:05 PM?
4. A thermometer initially reading $212^{\circ} \mathrm{F}$ is placed in a room where the temperature is $70^{\circ} \mathrm{F}$. After 2 minutes the thermometer reads $125^{\circ} \mathrm{F}$.
(a) What does the thermometer read after 4 minutes?
(b) When will the thermometer read $72^{\circ} \mathrm{F}$ ?
(c) When will the thermometer read $69^{\circ} \mathrm{F}$ ?
5. An object with initial temperature $150^{\circ} \mathrm{C}$ is placed outside, where the temperature is $35^{\circ} \mathrm{C}$. Its temperatures at $12: 15$ and $12: 20$ are $120^{\circ} \mathrm{C}$ and $90^{\circ} \mathrm{C}$, respectively.
(a) At what time was the object placed outside?
(b) When will its temperature be $40^{\circ} \mathrm{C}$ ?
6. An object is placed in a room where the temperature is $20^{\circ} \mathrm{C}$. The temperature of the object drops by $5^{\circ} \mathrm{C}$ in 4 minutes and by $7^{\circ} \mathrm{C}$ in 8 minutes. What was the temperature of the object when it was initially placed in the room?
7. A cup of boiling water is placed outside at 1:00 PM. One minute later the temperature of the water is $152^{\circ} \mathrm{F}$. After another minute its temperature is $112^{\circ} \mathrm{F}$. What is the outside temperature?
8. A tank initially contains 40 gallons of pure water. A solution with 1 gram of salt per gallon of water is added to the tank at $3 \mathrm{gal} / \mathrm{min}$, and the resulting solution drains out at the same rate. Find the quantity $Q(t)$ of salt in the tank at time $t>0$.
9. A tank initially contains a solution of 10 pounds of salt in 60 gallons of water. Water with $1 / 2$ pound of salt per gallon is added to the tank at $6 \mathrm{gal} / \mathrm{min}$, and the resulting solution leaves at the same rate. Find the quantity $Q(t)$ of salt in the tank at time $t>0$.
10. A tank initially contains 100 liters of a salt solution with a concentration of $.1 \mathrm{~g} / \mathrm{liter}$. A solution with a salt concentration of $.3 \mathrm{~g} / \mathrm{liter}$ is added to the tank at 5 liters $/ \mathrm{min}$, and the resulting mixture is drained out at the same rate. Find the concentration $K(t)$ of salt in the tank as a function of $t$.
11. A 200 gallon tank initially contains 100 gallons of water with 20 pounds of salt. A salt solution with $1 / 4$ pound of salt per gallon is added to the tank at $4 \mathrm{gal} / \mathrm{min}$, and the resulting mixture is drained out at $2 \mathrm{gal} / \mathrm{min}$. Find the quantity of salt in the tank as it's about to overflow.
12. Suppose water is added to a tank at $10 \mathrm{gal} / \mathrm{min}$, but leaks out at the rate of $1 / 5 \mathrm{gal} / \mathrm{min}$ for each gallon in the tank. What is the smallest capacity the tank can have if the process is to continue indefinitely?
13. A chemical reaction in a laboratory with volume $V$ (in $\left.\mathrm{ft}^{3}\right)$ produces $q_{1} \mathrm{ft}^{3} / \mathrm{min}$ of a noxious gas as a byproduct. The gas is dangerous at concentrations greater than $\bar{c}$, but harmless at concentrations $\leq \bar{c}$. Intake fans at one end of the laboratory pull in fresh air at the rate of $q_{2} \mathrm{ft}^{3} / \mathrm{min}$ and exhaust fans at the other end exhaust the mixture of gas and air from the laboratory at the same rate. Assuming that the gas is always uniformly distributed in the room and its initial concentration $c_{0}$ is at a safe level, find the smallest value of $q_{2}$ required to maintain safe conditions in the laboratory for all time.
14. A 1200 -gallon tank initially contains 40 pounds of salt dissolved in 600 gallons of water. Starting at $t_{0}=0$, water that contains $1 / 2$ pound of salt per gallon is added to the tank at the rate of 6 $\mathrm{gal} / \mathrm{min}$ and the resulting mixture is drained from the tank at $4 \mathrm{gal} / \mathrm{min}$. Find the quantity $Q(t)$ of salt in the tank at any time $t>0$ prior to overflow.
15. Tank $T_{1}$ initially contain 50 gallons of pure water. Starting at $t_{0}=0$, water that contains 1 pound of salt per gallon is poured into $T_{1}$ at the rate of $2 \mathrm{gal} / \mathrm{min}$. The mixture is drained from $T_{1}$ at the same rate into a second tank $T_{2}$, which initially contains 50 gallons of pure water. Also starting at $t_{0}=0$, a mixture from another source that contains 2 pounds of salt per gallon is poured into $T_{2}$ at the rate of $2 \mathrm{gal} / \mathrm{min}$. The mixture is drained from $T_{2}$ at the rate of $4 \mathrm{gal} / \mathrm{min}$.
(a) Find a differential equation for the quantity $Q(t)$ of salt in $\operatorname{tank} T_{2}$ at time $t>0$.
(b) Solve the equation derived in (a) to determine $Q(t)$.
(c) Find $\lim _{t \rightarrow \infty} Q(t)$.
16. Suppose an object with initial temperature $T_{0}$ is placed in a sealed container, which is in turn placed in a medium with temperature $T_{m}$. Let the initial temperature of the container be $S_{0}$. Assume that the temperature of the object does not affect the temperature of the container, which in turn does not affect the temperature of the medium. (These assumptions are reasonable, for example, if the object is a cup of coffee, the container is a house, and the medium is the atmosphere.)
(a) Assuming that the container and the medium have distinct temperature decay constants $k$ and $k_{m}$ respectively, use Newton's law of cooling to find the temperatures $S(t)$ and $T(t)$ of the container and object at time $t$.
(b) Assuming that the container and the medium have the same temperature decay constant $k$, use Newton's law of cooling to find the temperatures $S(t)$ and $T(t)$ of the container and object at time $t$.
(c) Find $\lim { }_{\cdot t \rightarrow \infty} S(t)$ and $\lim _{t \rightarrow \infty} T(t)$.
17. In our previous examples and exercises concerning Newton's law of cooling we assumed that the temperature of the medium remains constant. This model is adequate if the heat lost or gained by the object is insignificant compared to the heat required to cause an appreciable change in the temperature of the medium. If this isn't so, we must use a model that accounts for the heat exchanged between the object and the medium. Let $T=T(t)$ and $T_{m}=T_{m}(t)$ be the temperatures of the object and the medium, respectively, and let $T_{0}$ and $T_{m 0}$ be their initial values. Again, we assume that $T$ and $T_{m}$ are related by Newton's law of cooling,

$$
\begin{equation*}
T^{\prime}=-k\left(T-T_{m}\right) . \tag{A}
\end{equation*}
$$

We also assume that the change in heat of the object as its temperature changes from $T_{0}$ to $T$ is $a\left(T-T_{0}\right)$ and that the change in heat of the medium as its temperature changes from $T_{m 0}$ to $T_{m}$ is $a_{m}\left(T_{m}-T_{m 0}\right)$, where $a$ and $a_{m}$ are positive constants depending upon the masses and thermal properties of the object and medium, respectively. If we assume that the total heat of the system consisting of the object and the medium remains constant (that is, energy is conserved), then

$$
\begin{equation*}
a\left(T-T_{0}\right)+a_{m}\left(T_{m}-T_{m 0}\right)=0 . \tag{B}
\end{equation*}
$$

(a) Equation (A) involves two unknown functions $T$ and $T_{m}$. Use (A) and (B) to derive a differential equation involving only $T$.
(b) Find $T(t)$ and $T_{m}(t)$ for $t>0$.
(c) Find $\lim _{t \rightarrow \infty} T(t)$ and $\lim _{t \rightarrow \infty} T_{m}(t)$.
18. Control mechanisms allow fluid to flow into a tank at a rate proportional to the volume $V$ of fluid in the tank, and to flow out at a rate proportional to $V^{2}$. Suppose $V(0)=V_{0}$ and the constants of proportionality are $a$ and $b$, respectively. Find $V(t)$ for $t>0$ and find $\lim _{t \rightarrow \infty} V(t)$.
19. Identical tanks $T_{1}$ and $T_{2}$ initially contain $W$ gallons each of pure water. Starting at $t_{0}=0$, a salt solution with constant concentration $c$ is pumped into $T_{1}$ at $r \mathrm{gal} / \mathrm{min}$ and drained from $T_{1}$ into $T_{2}$ at the same rate. The resulting mixture in $T_{2}$ is also drained at the same rate. Find the concentrations $c_{1}(t)$ and $c_{2}(t)$ in tanks $T_{1}$ and $T_{2}$ for $t>0$.
20. An infinite sequence of identical tanks $T_{1}, T_{2}, \ldots, T_{n}, \ldots$, initially contain $W$ gallons each of pure water. They are hooked together so that fluid drains from $T_{n}$ into $T_{n+1}(n=1,2, \cdots)$. A salt solution is circulated through the tanks so that it enters and leaves each tank at the constant rate of $r \mathrm{gal} / \mathrm{min}$. The solution has a concentration of $c$ pounds of salt per gallon when it enters $T_{1}$.
(a) Find the concentration $c_{n}(t)$ in $\operatorname{tank} T_{n}$ for $t>0$.
(b) Find $\lim _{t \rightarrow \infty} c_{n}(t)$ for each $n$.
21. Tanks $T_{1}$ and $T_{2}$ have capacities $W_{1}$ and $W_{2}$ liters, respectively. Initially they are both full of dye solutions with concentrations $c_{1}$ and $c_{2}$ grams per liter. Starting at $t_{0}=0$, the solution from $T_{1}$ is pumped into $T_{2}$ at a rate of $r$ liters per minute, and the solution from $T_{2}$ is pumped into $T_{1}$ at the same rate.
(a) Find the concentrations $c_{1}(t)$ and $c_{2}(t)$ of the dye in $T_{1}$ and $T_{2}$ for $t>0$.
(b) Find $\lim _{t \rightarrow \infty} c_{1}(t)$ and $\lim _{t \rightarrow \infty} c_{2}(t)$.
22. $L$ Consider the mixing problem of Example 4.2.3, but without the assumption that the mixture is stirred instantly so that the salt is always uniformly distributed throughout the mixture. Assume instead that the distribution approaches uniformity as $t \rightarrow \infty$. In this case the differential equation for $Q$ is of the form

$$
Q^{\prime}+\frac{a(t)}{150} Q=2
$$

where $\lim _{t \rightarrow \infty} a(t)=1$.
(a) Assuming that $Q(0)=Q_{0}$, can you guess the value of $\lim _{t \rightarrow \infty} Q(t)$ ?.
(b) Use numerical methods to confirm your guess in the these cases:

$$
\begin{array}{ll}
\text { (i) } a(t)=t /(1+t) & \text { (ii) } a(t)=1-e^{-t^{2}}
\end{array} \quad \text { (iii) } a(t)=1-\sin \left(e^{-t}\right)
$$

23. $L$ Consider the mixing problem of Example 4.2 .4 in a tank with infinite capacity, but without the assumption that the mixture is stirred instantly so that the salt is always uniformly distributed throughout the mixture. Assume instead that the distribution approaches uniformity as $t \rightarrow \infty$. In this case the differential equation for $Q$ is of the form

$$
Q^{\prime}+\frac{a(t)}{t+100} Q=1
$$

where $\lim _{t \rightarrow \infty} a(t)=1$.
(a) Let $K(t)$ be the concentration of salt at time $t$. Assuming that $Q(0)=Q_{0}$, can you guess the value of $\lim _{t \rightarrow \infty} K(t)$ ?
(b) Use numerical methods to confirm your guess in the these cases:

$$
\begin{array}{ll}
\text { (i) } a(t)=t /(1+t) & \text { (ii) } a(t)=1-e^{-t^{2}} \\
\text { (iii) } a(t)=1+\sin \left(e^{-t}\right) \text {. }
\end{array}
$$

### 4.3 ELEMENTARY MECHANICS

## Newton's Second Law of Motion

In this section we consider an object with constant mass $m$ moving along a line under a force $F$. Let $y=y(t)$ be the displacement of the object from a reference point on the line at time $t$, and let $v=v(t)$ and $a=a(t)$ be the velocity and acceleration of the object at time $t$. Thus, $v=y^{\prime}$ and $a=v^{\prime}=y^{\prime \prime}$, where the prime denotes differentiation with respect to $t$. Newton's second law of motion asserts that the force $F$ and the acceleration $a$ are related by the equation

$$
\begin{equation*}
F=m a \tag{4.3.1}
\end{equation*}
$$

Units
In applications there are three main sets of units in use for length, mass, force, and time: the cgs, mks, and British systems. All three use the second as the unit of time. Table 1 shows the other units. Consistent with (4.3.1), the unit of force in each system is defined to be the force required to impart an acceleration of (one unit of length) $/ s^{2}$ to one unit of mass.

|  | Length | Force | Mass |
| :---: | :---: | :---: | :---: |
| cgs | centimeter $(\mathrm{cm})$ | dyne $(\mathrm{d})$ | gram $(\mathrm{g})$ |
| mks | meter $(\mathrm{m})$ | newton $(\mathrm{N})$ | kilogram $(\mathrm{kg})$ |
| British | foot $(\mathrm{ft})$ | pound $(\mathrm{lb})$ | slug $(\mathrm{sl})$ |

Table 1.

If we assume that Earth is a perfect sphere with constant mass density, Newton's law of gravitation (discussed later in this section) asserts that the force exerted on an object by Earth's gravitational field is proportional to the mass of the object and inversely proportional to the square of its distance from the center of Earth. However, if the object remains sufficiently close to Earth's surface, we may assume that the gravitational force is constant and equal to its value at the surface. The magnitude of this force is $m g$, where $g$ is called the acceleration due to gravity. (To be completely accurate, $g$ should be called the magnitude of the acceleration due to gravity at Earth's surface.) This quantity has been determined experimentally. Approximate values of $g$ are

$$
\begin{array}{lll}
g & =980 \mathrm{~cm} / \mathrm{s}^{2} & \\
g & =9.8 \mathrm{~m} / \mathrm{s}^{2} & \\
g & =32 \mathrm{ft} / \mathrm{s}^{2} & \\
\text { (mks) } \\
\text { (British). }
\end{array}
$$

In general, the force $F$ in (4.3.1) may depend upon $t, y$, and $y^{\prime}$. Since $a=y^{\prime \prime}$, (4.3.1) can be written in the form

$$
\begin{equation*}
m y^{\prime \prime}=F\left(t, y, y^{\prime}\right) \tag{4.3.2}
\end{equation*}
$$

which is a second order equation. We'll consider this equation with restrictions on $F$ later; however, since Chapter 2 dealt only with first order equations, we consider here only problems in which (4.3.2) can be recast as a first order equation. This is possible if $F$ does not depend on $y$, so (4.3.2) is of the form

$$
m y^{\prime \prime}=F\left(t, y^{\prime}\right)
$$

Letting $v=y^{\prime}$ and $v^{\prime}=y^{\prime \prime}$ yields a first order equation for $v$ :

$$
\begin{equation*}
m v^{\prime}=F(t, v) \tag{4.3.3}
\end{equation*}
$$

Solving this equation yields $v$ as a function of $t$. If we know $y\left(t_{0}\right)$ for some time $t_{0}$, we can integrate $v$ to obtain $y$ as a function of $t$.

Equations of the form (4.3.3) occur in problems involving motion through a resisting medium.

## Motion Through a Resisting Medium Under Constant Gravitational Force

Now we consider an object moving vertically in some medium. We assume that the only forces acting on the object are gravity and resistance from the medium. We also assume that the motion takes place close to Earth's surface and take the upward direction to be positive, so the gravitational force can be assumed to have the constant value $-m g$. We'll see that, under reasonable assumptions on the resisting force, the velocity approaches a limit as $t \rightarrow \infty$. We call this limit the terminal velocity.

Example 4.3.1 An object with mass $m$ moves under constant gravitational force through a medium that exerts a resistance with magnitude proportional to the speed of the object. (Recall that the speed of an object is $|v|$, the absolute value of its velocity $v$.) Find the velocity of the object as a function of $t$, and find the terminal velocity. Assume that the initial velocity is $v_{0}$.

Solution The total force acting on the object is

$$
\begin{equation*}
F=-m g+F_{1} \tag{4.3.4}
\end{equation*}
$$

where $-m g$ is the force due to gravity and $F_{1}$ is the resisting force of the medium, which has magnitude $k|v|$, where $k$ is a positive constant. If the object is moving downward $(v \leq 0)$, the resisting force is upward (Figure 4.3.1(a)), so

$$
F_{1}=k|v|=k(-v)=-k v
$$

On the other hand, if the object is moving upward ( $v \geq 0$ ), the resisting force is downward (Figure 4.3.1(b)), so

$$
F_{1}=-k|v|=-k v .
$$

Thus, (4.3.4) can be written as

$$
\begin{equation*}
F=-m g-k v \tag{4.3.5}
\end{equation*}
$$

regardless of the sign of the velocity.
From Newton's second law of motion,

$$
F=m a=m v^{\prime}
$$



Figure 4.3.1 Resistive forces
so (4.3.5) yields

$$
m v^{\prime}=-m g-k v
$$

or

$$
\begin{equation*}
v^{\prime}+\frac{k}{m} v=-g \tag{4.3.6}
\end{equation*}
$$

Since $e^{-k t / m}$ is a solution of the complementary equation, the solutions of (4.3.6) are of the form $v=$ $u e^{-k t / m}$, where $u^{\prime} e^{-k t / m}=-g$, so $u^{\prime}=-g e^{k t / m}$. Hence,

$$
u=-\frac{m g}{k} e^{k t / m}+c
$$

so

$$
\begin{equation*}
v=u e^{-k t / m}=-\frac{m g}{k}+c e^{-k t / m} . \tag{4.3.7}
\end{equation*}
$$

Since $v(0)=v_{0}$,
so

$$
v_{0}=-\frac{m g}{k}+c
$$

$$
c=v_{0}+\frac{m g}{k}
$$

and (4.3.7) becomes

$$
v=-\frac{m g}{k}+\left(v_{0}+\frac{m g}{k}\right) e^{-k t / m}
$$

Letting $t \rightarrow \infty$ here shows that the terminal velocity is

$$
\lim _{t \rightarrow \infty} v(t)=-\frac{m g}{k}
$$

which is independent of the initial velocity $v_{0}$ (Figure 4.3.2).
Example 4.3.2 A 960-lb object is given an initial upward velocity of $60 \mathrm{ft} / \mathrm{s}$ near the surface of Earth. The atmosphere resists the motion with a force of 3 lb for each $\mathrm{ft} / \mathrm{s}$ of speed. Assuming that the only other force acting on the object is constant gravity, find its velocity $v$ as a function of $t$, and find its terminal velocity.


Figure 4.3.2 Solutions of $m v^{\prime}=-m g-k v$

Solution Since $m g=960$ and $g=32, m=960 / 32=30$. The atmospheric resistance is $-3 v \mathrm{lb}$ if $v$ is expressed in feet per second. Therefore

$$
30 v^{\prime}=-960-3 v
$$

which we rewrite as

$$
v^{\prime}+\frac{1}{10} v=-32 .
$$

Since $e^{-t / 10}$ is a solution of the complementary equation, the solutions of this equation are of the form $v=u e^{-t / 10}$, where $u^{\prime} e^{-t / 10}=-32$, so $u^{\prime}=-32 e^{t / 10}$. Hence,

$$
u=-320 e^{t / 10}+c
$$

so

$$
\begin{equation*}
v=u e^{-t / 10}=-320+c e^{-t / 10} \tag{4.3.8}
\end{equation*}
$$

The initial velocity is $60 \mathrm{ft} / \mathrm{s}$ in the upward (positive) direction; hence, $v_{0}=60$. Substituting $t=0$ and $v=60$ in (4.3.8) yields

$$
60=-320+c,
$$

so $c=380$, and (4.3.8) becomes

$$
v=-320+380 e^{-t / 10} \mathrm{ft} / \mathrm{s}
$$

The terminal velocity is

$$
\lim _{t \rightarrow \infty} v(t)=-320 \mathrm{ft} / \mathrm{s}
$$

Example 4.3.3 A 10 kg mass is given an initial velocity $v_{0} \leq 0$ near Earth's surface. The only forces acting on it are gravity and atmospheric resistance proportional to the square of the speed. Assuming that the resistance is 8 N if the speed is $2 \mathrm{~m} / \mathrm{s}$, find the velocity of the object as a function of $t$, and find the terminal velocity.

Solution Since the object is falling, the resistance is in the upward (positive) direction. Hence,

$$
\begin{equation*}
m v^{\prime}=-m g+k v^{2} \tag{4.3.9}
\end{equation*}
$$

where $k$ is a constant. Since the magnitude of the resistance is 8 N when $v=2 \mathrm{~m} / \mathrm{s}$,

$$
k\left(2^{2}\right)=8
$$

so $k=2 \mathrm{~N}-\mathrm{s}^{2} / \mathrm{m}^{2}$. Since $m=10$ and $g=9.8$, (4.3.9) becomes

$$
\begin{equation*}
10 v^{\prime}=-98+2 v^{2}=2\left(v^{2}-49\right) \tag{4.3.10}
\end{equation*}
$$

If $v_{0}=-7$, then $v \equiv-7$ for all $t \geq 0$. If $v_{0} \neq-7$, we separate variables to obtain

$$
\begin{equation*}
\frac{1}{v^{2}-49} v^{\prime}=\frac{1}{5} \tag{4.3.11}
\end{equation*}
$$

which is convenient for the required partial fraction expansion

$$
\begin{equation*}
\frac{1}{v^{2}-49}=\frac{1}{(v-7)(v+7)}=\frac{1}{14}\left[\frac{1}{v-7}-\frac{1}{v+7}\right] . \tag{4.3.12}
\end{equation*}
$$

Substituting (4.3.12) into (4.3.11) yields

$$
\frac{1}{14}\left[\frac{1}{v-7}-\frac{1}{v+7}\right] v^{\prime}=\frac{1}{5}
$$

so

$$
\left[\frac{1}{v-7}-\frac{1}{v+7}\right] v^{\prime}=\frac{14}{5}
$$

Integrating this yields

$$
\ln |v-7|-\ln |v+7|=14 t / 5+k
$$

Therefore

$$
\left|\frac{v-7}{v+7}\right|=e^{k} e^{14 t / 5}
$$

Since Theorem 2.3.1 implies that $(v-7) /(v+7)$ can't change sign (why?), we can rewrite the last equation as

$$
\begin{equation*}
\frac{v-7}{v+7}=c e^{14 t / 5} \tag{4.3.13}
\end{equation*}
$$

which is an implicit solution of (4.3.10). Solving this for $v$ yields

$$
\begin{equation*}
v=-7 \frac{c+e^{-14 t / 5}}{c-e^{-14 t / 5}} \tag{4.3.14}
\end{equation*}
$$

Since $v(0)=v_{0}$, it (4.3.13) implies that

$$
c=\frac{v_{0}-7}{v_{0}+7}
$$

Substituting this into (4.3.14) and simplifying yields

$$
v=-7 \frac{v_{0}\left(1+e^{-14 t / 5}\right)-7\left(1-e^{-14 t / 5}\right)}{v_{0}\left(1-e^{-14 t / 5}\right)-7\left(1+e^{-14 t / 5}\right.}
$$

Since $v_{0} \leq 0, v$ is defined and negative for all $t>0$. The terminal velocity is

$$
\lim _{t \rightarrow \infty} v(t)=-7 \mathrm{~m} / \mathrm{s}
$$

independent of $v_{0}$. More generally, it can be shown (Exercise 11) that if $v$ is any solution of (4.3.9) such that $v_{0} \leq 0$ then

$$
\lim _{t \rightarrow \infty} v(t)=-\sqrt{\frac{m g}{k}}
$$

(Figure 4.3.3).
Example 4.3.4 A 10-kg mass is launched vertically upward from Earth's surface with an initial velocity of $v_{0} \mathrm{~m} / \mathrm{s}$. The only forces acting on the mass are gravity and atmospheric resistance proportional to the square of the speed. Assuming that the atmospheric resistance is 8 N if the speed is $2 \mathrm{~m} / \mathrm{s}$, find the time $T$ required for the mass to reach maximum altitude.


Figure 4.3.3 Solutions of $m v^{\prime}=-m g+k v^{2}, v(0)=v_{0} \leq 0$

Solution The mass will climb while $v>0$ and reach its maximum altitude when $v=0$. Therefore $v>0$ for $0 \leq t<T$ and $v(T)=0$. Although the mass of the object and our assumptions concerning the forces acting on it are the same as those in Example 3, (4.3.10) does not apply here, since the resisting force is negative if $v>0$; therefore, we replace (4.3.10) by

$$
\begin{equation*}
10 v^{\prime}=-98-2 v^{2} \tag{4.3.15}
\end{equation*}
$$

Separating variables yields

$$
\frac{5}{v^{2}+49} v^{\prime}=-1
$$

and integrating this yields

$$
\frac{5}{7} \tan ^{-1} \frac{v}{7}=-t+c
$$

(Recall that $\tan ^{-1} u$ is the number $\theta$ such that $-\pi / 2<\theta<\pi / 2$ and $\tan \theta=u$.) Since $v(0)=v_{0}$,

$$
c=\frac{5}{7} \tan ^{-1} \frac{v_{0}}{7}
$$

so $v$ is defined implicitly by

$$
\begin{equation*}
\frac{5}{7} \tan ^{-1} \frac{v}{7}=-t+\frac{5}{7} \tan ^{-1} \frac{v_{0}}{7}, \quad 0 \leq t \leq T \tag{4.3.16}
\end{equation*}
$$

Solving this for $v$ yields

$$
\begin{equation*}
v=7 \tan \left(-\frac{7 t}{5}+\tan ^{-1} \frac{v_{0}}{7}\right) \tag{4.3.17}
\end{equation*}
$$

Using the identity

$$
\tan (A-B)=\frac{\tan A-\tan B}{1+\tan A \tan B}
$$

with $A=\tan ^{-1}\left(v_{0} / 7\right)$ and $B=7 t / 5$, and noting that $\tan \left(\tan ^{-1} \theta\right)=\theta$, we can simplify (4.3.17) to

$$
v=7 \frac{v_{0}-7 \tan (7 t / 5)}{7+v_{0} \tan (7 t / 5)}
$$

Since $v(T)=0$ and $\tan ^{-1}(0)=0$, (4.3.16) implies that

$$
-T+\frac{5}{7} \tan ^{-1} \frac{v_{0}}{7}=0
$$



Figure 4.3.4 Solutions of (4.3.15) for various $v_{0}>0$

Therefore

$$
T=\frac{5}{7} \tan ^{-1} \frac{v_{0}}{7} .
$$

Since $\tan ^{-1}\left(v_{0} / 7\right)<\pi / 2$ for all $v_{0}$, the time required for the mass to reach its maximum altitude is less than

$$
\frac{5 \pi}{14} \approx 1.122 \mathrm{~s}
$$

regardless of the initial velocity. Figure 4.3.4 shows graphs of $v$ over $[0, T]$ for various values of $v_{0}$.


Figure 4.3.5 Escape velocity

Escape Velocity
Suppose a space vehicle is launched vertically and its fuel is exhausted when the vehicle reaches an altitude $h$ above Earth, where $h$ is sufficiently large so that resistance due to Earth's atmosphere can be neglected. Let $t=0$ be the time when burnout occurs. Assuming that the gravitational forces of all other celestial bodies can be neglected, the motion of the vehicle for $t>0$ is that of an object with constant mass $m$ under the influence of Earth's gravitational force, which we now assume to vary inversely with the square of the distance from Earth's center; thus, if we take the upward direction to be positive then gravitational force on the vehicle at an altitude $y$ above Earth is

$$
\begin{equation*}
F=-\frac{K}{(y+R)^{2}} \tag{4.3.18}
\end{equation*}
$$

where $R$ is Earth's radius (Figure 4.3.5).
Since $F=-m g$ when $y=0$, setting $y=0$ in (4.3.18) yields

$$
-m g=-\frac{K}{R^{2}}
$$

therefore $K=m g R^{2}$ and (4.3.18) can be written more specifically as

$$
\begin{equation*}
F=-\frac{m g R^{2}}{(y+R)^{2}} \tag{4.3.19}
\end{equation*}
$$

From Newton's second law of motion,

$$
F=m \frac{d^{2} y}{d t^{2}}
$$

so (4.3.19) implies that

$$
\begin{equation*}
\frac{d^{2} y}{d t^{2}}=-\frac{g R^{2}}{(y+R)^{2}} \tag{4.3.20}
\end{equation*}
$$

We'll show that there's a number $v_{e}$, called the escape velocity, with these properties:

1. If $v_{0} \geq v_{e}$ then $v(t)>0$ for all $t>0$, and the vehicle continues to climb for all $t>0$; that is, it "escapes" Earth. (Is it really so obvious that $\lim _{t \rightarrow \infty} y(t)=\infty$ in this case? For a proof, see Exercise 20.)
2. If $v_{0}<v_{e}$ then $v(t)$ decreases to zero and becomes negative. Therefore the vehicle attains a maximum altitude $y_{m}$ and falls back to Earth.

Since (4.3.20) is second order, we can't solve it by methods discussed so far. However, we're concerned with $v$ rather than $y$, and $v$ is easier to find. Since $v=y^{\prime}$ the chain rule implies that

$$
\frac{d^{2} y}{d t^{2}}=\frac{d v}{d t}=\frac{d v}{d y} \frac{d y}{d t}=v \frac{d v}{d y}
$$

Substituting this into (4.3.20) yields the first order separable equation

$$
\begin{equation*}
v \frac{d v}{d y}=-\frac{g R^{2}}{(y+R)^{2}} \tag{4.3.21}
\end{equation*}
$$

When $t=0$, the velocity is $v_{0}$ and the altitude is $h$. Therefore we can obtain $v$ as a function of $y$ by solving the initial value problem

$$
v \frac{d v}{d y}=-\frac{g R^{2}}{(y+R)^{2}}, \quad v(h)=v_{0} .
$$

Integrating (4.3.21) with respect to $y$ yields

$$
\begin{equation*}
\frac{v^{2}}{2}=\frac{g R^{2}}{y+R}+c . \tag{4.3.22}
\end{equation*}
$$

Since $v(h)=v_{0}$,

$$
c=\frac{v_{0}^{2}}{2}-\frac{g R^{2}}{h+R},
$$

so (4.3.22) becomes

$$
\begin{equation*}
\frac{v^{2}}{2}=\frac{g R^{2}}{y+R}+\left(\frac{v_{0}^{2}}{2}-\frac{g R^{2}}{h+R}\right) \tag{4.3.23}
\end{equation*}
$$

If

$$
v_{0} \geq\left(\frac{2 g R^{2}}{h+R}\right)^{1 / 2}
$$

the parenthetical expression in (4.3.23) is nonnegative, so $v(y)>0$ for $y>h$. This proves that there's an escape velocity $v_{e}$. We'll now prove that

$$
v_{e}=\left(\frac{2 g R^{2}}{h+R}\right)^{1 / 2}
$$

by showing that the vehicle falls back to Earth if

$$
\begin{equation*}
v_{0}<\left(\frac{2 g R^{2}}{h+R}\right)^{1 / 2} \tag{4.3.24}
\end{equation*}
$$

If (4.3.24) holds then the parenthetical expression in (4.3.23) is negative and the vehicle will attain a maximum altitude $y_{m}>h$ that satisfies the equation

$$
0=\frac{g R^{2}}{y_{m}+R}+\left(\frac{v_{0}^{2}}{2}-\frac{g R^{2}}{h+R}\right)
$$

The velocity will be zero at the maximum altitude, and the object will then fall to Earth under the influence of gravity.

### 4.3 Exercises

Except where directed otherwise, assume that the magnitude of the gravitational force on an object with mass $m$ is constant and equal to $m g$. In exercises involving vertical motion take the upward direction to be positive.

1. A firefighter who weighs 192 lb slides down an infinitely long fire pole that exerts a frictional resistive force with magnitude proportional to his speed, with $k=2.5 \mathrm{lb}-\mathrm{s} / \mathrm{ft}$. Assuming that he starts from rest, find his velocity as a function of time and find his terminal velocity.
2. A firefighter who weighs 192 lb slides down an infinitely long fire pole that exerts a frictional resistive force with magnitude proportional to her speed, with constant of proportionality $k$. Find $k$, given that her terminal velocity is $-16 \mathrm{ft} / \mathrm{s}$, and then find her velocity $v$ as a function of $t$. Assume that she starts from rest.
3. A boat weighs $64,000 \mathrm{lb}$. Its propellor produces a constant thrust of $50,000 \mathrm{lb}$ and the water exerts a resistive force with magnitude proportional to the speed, with $k=2000 \mathrm{lb}-\mathrm{s} / \mathrm{ft}$. Assuming that the boat starts from rest, find its velocity as a function of time, and find its terminal velocity.
4. A constant horizontal force of 10 N pushes a 20 kg -mass through a medium that resists its motion with .5 N for every $\mathrm{m} / \mathrm{s}$ of speed. The initial velocity of the mass is $7 \mathrm{~m} / \mathrm{s}$ in the direction opposite to the direction of the applied force. Find the velocity of the mass for $t>0$.
5. A stone weighing $1 / 2 \mathrm{lb}$ is thrown upward from an initial height of 5 ft with an initial speed of 32 $\mathrm{ft} / \mathrm{s}$. Air resistance is proportional to speed, with $k=1 / 128 \mathrm{lb}-\mathrm{s} / \mathrm{ft}$. Find the maximum height attained by the stone.
6. A $3200-\mathrm{lb}$ car is moving at $64 \mathrm{ft} / \mathrm{s}$ down a 30 -degree grade when it runs out of fuel. Find its velocity after that if friction exerts a resistive force with magnitude proportional to the square of the speed, with $k=1 \mathrm{lb}-\mathrm{s}^{2} / \mathrm{ft}{ }^{2}$. Also find its terminal velocity.
7. A 96 lb weight is dropped from rest in a medium that exerts a resistive force with magnitude proportional to the speed. Find its velocity as a function of time if its terminal velocity is $-128 \mathrm{ft} / \mathrm{s}$.
8. An object with mass $m$ moves vertically through a medium that exerts a resistive force with magnitude proportional to the speed. Let $y=y(t)$ be the altitude of the object at time $t$, with $y(0)=y_{0}$. Use the results of Example 4.3.1 to show that

$$
y(t)=y_{0}+\frac{m}{k}\left(v_{0}-v-g t\right) .
$$

9. An object with mass $m$ is launched vertically upward with initial velocity $v_{0}$ from Earth's surface $\left(y_{0}=0\right)$ in a medium that exerts a resistive force with magnitude proportional to the speed. Find the time $T$ when the object attains its maximum altitude $y_{m}$. Then use the result of Exercise 8 to find $y_{m}$.
10. An object weighing 256 lb is dropped from rest in a medium that exerts a resistive force with magnitude proportional to the square of the speed. The magnitude of the resisting force is 1 lb when $|v|=4 \mathrm{ft} / \mathrm{s}$. Find $v$ for $t>0$, and find its terminal velocity.
11. An object with mass $m$ is given an initial velocity $v_{0} \leq 0$ in a medium that exerts a resistive force with magnitude proportional to the square of the speed. Find the velocity of the object for $t>0$, and find its terminal velocity.
12. An object with mass $m$ is launched vertically upward with initial velocity $v_{0}$ in a medium that exerts a resistive force with magnitude proportional to the square of the speed.
(a) Find the time $T$ when the object reaches its maximum altitude.
(b) Use the result of Exercise 11 to find the velocity of the object for $t>T$.
13. $L$ An object with mass $m$ is given an initial velocity $v_{0} \leq 0$ in a medium that exerts a resistive force of the form $a|v| /(1+|v|)$, where $a$ is positive constant.
(a) Set up a differential equation for the speed of the object.
(b) Use your favorite numerical method to solve the equation you found in (a), to convince yourself that there's a unique number $a_{0}$ such that $\lim _{t \rightarrow \infty} s(t)=\infty$ if $a \leq a_{0}$ and $\lim _{t \rightarrow \infty} s(t)$ exists (finite) if $a>a_{0}$. (We say that $a_{0}$ is the bifurcation value of $a$.) Try to find $a_{0}$ and $\lim _{t \rightarrow \infty} s(t)$ in the case where $a>a_{0}$. Hint: See Exercise 14.
14. An object of mass $m$ falls in a medium that exerts a resistive force $f=f(s)$, where $s=|v|$ is the speed of the object. Assume that $f(0)=0$ and $f$ is strictly increasing and differentiable on $(0, \infty)$.
(a) Write a differential equation for the speed $s=s(t)$ of the object. Take it as given that all solutions of this equation with $s(0) \geq 0$ are defined for all $t>0$ (which makes good sense on physical grounds).
(b) Show that if $\lim _{s \rightarrow \infty} f(s) \leq m g$ then $\lim _{t \rightarrow \infty} s(t)=\infty$.
(c) Show that if $\lim _{s \rightarrow \infty} f(s)>m g$ then $\lim _{t \rightarrow \infty} s(t)=s_{T}$ (terminal speed), where $f\left(s_{T}\right)=$ mg. Hint: Use Theorem 2.3.1.
15. A $100-\mathrm{g}$ mass with initial velocity $v_{0} \leq 0$ falls in a medium that exerts a resistive force proportional to the fourth power of the speed. The resistance is .1 N if the speed is $3 \mathrm{~m} / \mathrm{s}$.
(a) Set up the initial value problem for the velocity $v$ of the mass for $t>0$.
(b) Use Exercise 14(c) to determine the terminal velocity of the object.
(c) Co confirm your answer to (b), use one of the numerical methods studied in Chapter 3 to compute approximate solutions on $[0,1]$ (seconds) of the initial value problem of (a), with initial values $v_{0}=0,-2,-4, \ldots,-12$. Present your results in graphical form similar to Figure 4.3.3.
16. A 64-lb object with initial velocity $v_{0} \leq 0$ falls through a dense fluid that exerts a resistive force proportional to the square root of the speed. The resistance is 64 lb if the speed is $16 \mathrm{ft} / \mathrm{s}$.
(a) Set up the initial value problem for the velocity $v$ of the mass for $t>0$.
(b) Use Exercise 14(c) to determine the terminal velocity of the object.
(c) C To confirm your answer to (b), use one of the numerical methods studied in Chapter 3 to compute approximate solutions on $[0,4]$ (seconds) of the initial value problem of (a), with initial values $v_{0}=0,-5,-10, \ldots,-30$. Present your results in graphical form similar to Figure 4.3.3.

In Exercises 17-20, assume that the force due to gravity is given by Newton's law of gravitation. Take the upward direction to be positive.
17. A space probe is to be launched from a space station 200 miles above Earth. Determine its escape velocity in miles/s. Take Earth's radius to be 3960 miles.
18. A space vehicle is to be launched from the moon, which has a radius of about 1080 miles. The acceleration due to gravity at the surface of the moon is about $5.31 \mathrm{ft} / \mathrm{s}^{2}$. Find the escape velocity in miles/s.
19. (a) Show that Eqn. (4.3.23) can be rewritten as

$$
v^{2}=\frac{h-y}{y+R} v_{e}^{2}+v_{0}^{2} .
$$

(b) Show that if $v_{0}=\rho v_{e}$ with $0 \leq \rho<1$, then the maximum altitude $y_{m}$ attained by the space vehicle is

$$
y_{m}=\frac{h+R \rho^{2}}{1-\rho^{2}}
$$

(c) By requiring that $v\left(y_{m}\right)=0$, use Eqn. (4.3.22) to deduce that if $v_{0}<v_{e}$ then

$$
|v|=v_{e}\left[\frac{\left(1-\rho^{2}\right)\left(y_{m}-y\right)}{y+R}\right]^{1 / 2},
$$

where $y_{m}$ and $\rho$ are as defined in (b) and $y \geq h$.
(d) Deduce from (c) that if $v<v_{e}$, the vehicle takes equal times to climb from $y=h$ to $y=y_{m}$ and to fall back from $y=y_{m}$ to $y=h$.
20. In the situation considered in the discussion of escape velocity, show that $\lim _{t \rightarrow \infty} y(t)=\infty$ if $v(t)>0$ for all $t>0$.
Hint: Use a proof by contradiction. Assume that there's a number $y_{m}$ such that $y(t) \leq y_{m}$ for all $t>0$. Deduce from this that there's positive number $\alpha$ such that $y^{\prime \prime}(t) \leq-\alpha$ for all $t \geq 0$. Show that this contradicts the assumption that $v(t)>0$ for all $t>0$.

### 4.4 AUTONOMOUS SECOND ORDER EQUATIONS

A second order differential equation that can be written as

$$
\begin{equation*}
y^{\prime \prime}=F\left(y, y^{\prime}\right) \tag{4.4.1}
\end{equation*}
$$

where $F$ is independent of $t$, is said to be autonomous. An autonomous second order equation can be converted into a first order equation relating $v=y^{\prime}$ and $y$. If we let $v=y^{\prime}$, (4.4.1) becomes

$$
\begin{equation*}
v^{\prime}=F(y, v) \tag{4.4.2}
\end{equation*}
$$

Since

$$
\begin{equation*}
v^{\prime}=\frac{d v}{d t}=\frac{d v}{d y} \frac{d y}{d t}=v \frac{d v}{d y} \tag{4.4.3}
\end{equation*}
$$

(4.4.2) can be rewritten as

$$
\begin{equation*}
v \frac{d v}{d y}=F(y, v) \tag{4.4.4}
\end{equation*}
$$

The integral curves of (4.4.4) can be plotted in the $(y, v)$ plane, which is called the Poincaré phase plane of (4.4.1). If $y$ is a solution of (4.4.1) then $y=y(t), v=y^{\prime}(t)$ is a parametric equation for an integral curve of (4.4.4). We'll call these integral curves trajectories of (4.4.1), and we'll call (4.4.4) the phase plane equivalent of (4.4.1).

In this section we'll consider autonomous equations that can be written as

$$
\begin{equation*}
y^{\prime \prime}+q\left(y, y^{\prime}\right) y^{\prime}+p(y)=0 . \tag{4.4.5}
\end{equation*}
$$

Equations of this form often arise in applications of Newton's second law of motion. For example, suppose $y$ is the displacement of a moving object with mass $m$. It's reasonable to think of two types of time-independent forces acting on the object. One type - such as gravity - depends only on position. We could write such a force as $-m p(y)$. The second type - such as atmospheric resistance or friction may depend on position and velocity. (Forces that depend on velocity are called damping forces.) We write this force as $-m q\left(y, y^{\prime}\right) y^{\prime}$, where $q\left(y, y^{\prime}\right)$ is usually a positive function and we've put the factor $y^{\prime}$ outside to make it explicit that the force is in the direction opposing the motion. In this case Newton's, second law of motion leads to (4.4.5).

The phase plane equivalent of (4.4.5) is

$$
\begin{equation*}
v \frac{d v}{d y}+q(y, v) v+p(y)=0 \tag{4.4.6}
\end{equation*}
$$

Some statements that we'll be making about the properties of (4.4.5) and (4.4.6) are intuitively reasonable, but difficult to prove. Therefore our presentation in this section will be informal: we'll just say things without proof, all of which are true if we assume that $p=p(y)$ is continuously differentiable for all $y$ and $q=q(y, v)$ is continuously differentiable for all $(y, v)$. We begin with the following statements:

- Statement 1. If $y_{0}$ and $v_{0}$ are arbitrary real numbers then (4.4.5) has a unique solution on $(-\infty, \infty)$ such that $y(0)=y_{0}$ and $y^{\prime}(0)=v_{0}$.
- Statement 2.) If $y=y(t)$ is a solution of (4.4.5) and $\tau$ is any constant then $y_{1}=y(t-\tau)$ is also a solution of (4.4.5), and $y$ and $y_{1}$ have the same trajectory.
- Statement 3. If two solutions $y$ and $y_{1}$ of (4.4.5) have the same trajectory then $y_{1}(t)=y(t-\tau)$ for some constant $\tau$.
- Statement 4. Distinct trajectories of (4.4.5) can't intersect; that is, if two trajectories of (4.4.5) intersect, they are identical.
- Statement 5. If the trajectory of a solution of (4.4.5) is a closed curve then $(y(t), v(t))$ traverses the trajectory in a finite time $T$, and the solution is periodic with period $T$; that is, $y(t+T)=y(t)$ for all $t$ in $(-\infty, \infty)$.

If $\bar{y}$ is a constant such that $p(\bar{y})=0$ then $y \equiv \bar{y}$ is a constant solution of (4.4.5). We say that $\bar{y}$ is an equilibrium of (4.4.5) and $(\bar{y}, 0)$ is a critical point of the phase plane equivalent equation (4.4.6). We say that the equilibrium and the critical point are stable if, for any given $\epsilon>0$ no matter how small, there's a $\delta>0$, sufficiently small, such that if

$$
\sqrt{\left(y_{0}-\bar{y}\right)^{2}+v_{0}^{2}}<\delta
$$

then the solution of the initial value problem

$$
y^{\prime \prime}+q\left(y, y^{\prime}\right) y^{\prime}+p(y)=0, \quad y(0)=y_{0}, \quad y^{\prime}(0)=v_{0}
$$

satisfies the inequality

$$
\sqrt{(y(t)-\bar{y})^{2}+(v(t))^{2}}<\epsilon
$$

for all $t>0$. Figure 4.4.1 illustrates the geometrical interpretation of this definition in the Poincaré phase plane: if $\left(y_{0}, v_{0}\right)$ is in the smaller shaded circle (with radius $\delta$ ), then $(y(t), v(t))$ must be in in the larger circle (with radius $\epsilon$ ) for all $t>0$.


Figure 4.4.1 Stability: if $\left(y_{0}, v_{0}\right)$ is in the smaller circle then $(y(t), v(t))$ is in the larger circle for all

$$
t>0
$$

If an equilibrium and the associated critical point are not stable, we say they are unstable. To see if you really understand what stable means, try to give a direct definition of unstable (Exercise 22). We'll illustrate both definitions in the following examples.

The Undamped Case
We'll begin with the case where $q \equiv 0$, so (4.4.5) reduces to

$$
\begin{equation*}
y^{\prime \prime}+p(y)=0 \tag{4.4.7}
\end{equation*}
$$

We say that this equation - as well as any physical situation that it may model - is undamped. The phase plane equivalent of (4.4.7) is the separable equation

$$
v \frac{d v}{d y}+p(y)=0
$$

Integrating this yields

$$
\begin{equation*}
\frac{v^{2}}{2}+P(y)=c \tag{4.4.8}
\end{equation*}
$$

where $c$ is a constant of integration and $P(y)=\int p(y) d y$ is an antiderivative of $p$.
If (4.4.7) is the equation of motion of an object of mass $m$, then $m v^{2} / 2$ is the kinetic energy and $m P(y)$ is the potential energy of the object; thus, (4.4.8) says that the total energy of the object remains constant, or is conserved. In particular, if a trajectory passes through a given point $\left(y_{0}, v_{0}\right)$ then

$$
c=\frac{v_{0}^{2}}{2}+P\left(y_{0}\right)
$$

Example 4.4.1 [The Undamped Spring - Mass System] Consider an object with mass $m$ suspended from a spring and moving vertically. Let $y$ be the displacement of the object from the position it occupies when suspended at rest from the spring (Figure 4.4.2).


Figure 4.4.2 (a) $y>0$ (b) $y=0$ (c) $y<0$

Assume that if the length of the spring is changed by an amount $\Delta L$ (positive or negative), then the spring exerts an opposing force with magnitude $k|\Delta L|$, where k is a positive constant. In Section 6.1 it will be shown that if the mass of the spring is negligible compared to $m$ and no other forces act on the object then Newton's second law of motion implies that

$$
\begin{equation*}
m y^{\prime \prime}=-k y \tag{4.4.9}
\end{equation*}
$$

which can be written in the form (4.4.7) with $p(y)=k y / m$. This equation can be solved easily by a method that we'll study in Section 5.2, but that method isn't available here. Instead, we'll consider the phase plane equivalent of (4.4.9).

From (4.4.3), we can rewrite (4.4.9) as the separable equation

$$
m v \frac{d v}{d y}=-k y
$$

Integrating this yields

$$
\frac{m v^{2}}{2}=-\frac{k y^{2}}{2}+c,
$$

which implies that

$$
\begin{equation*}
m v^{2}+k y^{2}=\rho \tag{4.4.10}
\end{equation*}
$$

( $\rho=2 c$ ). This defines an ellipse in the Poincaré phase plane (Figure 4.4.3).
We can identify $\rho$ by setting $t=0$ in (4.4.10); thus, $\rho=m v_{0}^{2}+k y_{0}^{2}$, where $y_{0}=y(0)$ and $v_{0}=v(0)$. To determine the maximum and minimum values of $y$ we set $v=0$ in (4.4.10); thus,

$$
\begin{equation*}
y_{\max }=R \quad \text { and } \quad y_{\min }=-R, \quad \text { with } R=\sqrt{\frac{\rho}{k}} . \tag{4.4.11}
\end{equation*}
$$

Equation (4.4.9) has exactly one equilibrium, $\bar{y}=0$, and it's stable. You can see intuitively why this is so: if the object is displaced in either direction from equilibrium, the spring tries to bring it back.

In this case we can find $y$ explicitly as a function of $t$. (Don't expect this to happen in more complicated problems!) If $v>0$ on an interval $I$, (4.4.10) implies that

$$
\frac{d y}{d t}=v=\sqrt{\frac{\rho-k y^{2}}{m}}
$$

on $I$. This is equivalent to

$$
\begin{equation*}
\frac{\sqrt{k}}{\sqrt{\rho-k y^{2}}} \frac{d y}{d t}=\omega_{0}, \quad \text { where } \quad \omega_{0}=\sqrt{\frac{k}{m}} . \tag{4.4.12}
\end{equation*}
$$



Figure 4.4.3 Trajectories of $m y^{\prime \prime}+k y=0$

Since

$$
\int \frac{\sqrt{k} d y}{\sqrt{\rho-k y^{2}}}=\sin ^{-1}\left(\sqrt{\frac{k}{\rho}} y\right)+c=\sin ^{-1}\left(\frac{y}{R}\right)+c
$$

(see (4.4.11)), (4.4.12) implies that that there's a constant $\phi$ such that

$$
\sin ^{-1}\left(\frac{y}{R}\right)=\omega_{0} t+\phi
$$

or

$$
y=R \sin \left(\omega_{0} t+\phi\right)
$$

for all $t$ in $I$. Although we obtained this function by assuming that $v>0$, you can easily verify that $y$ satisfies (4.4.9) for all values of $t$. Thus, the displacement varies periodically between $-R$ and $R$, with period $T=2 \pi / \omega_{0}$ (Figure 4.4.4). (If you've taken a course in elementary mechanics you may recognize this as simple harmonic motion.)

Example 4.4.2 [The Undamped Pendulum] Now we consider the motion of a pendulum with mass $m$, attached to the end of a weightless rod with length $L$ that rotates on a frictionless axle (Figure 4.4.5). We assume that there's no air resistance.

Let $y$ be the angle measured from the rest position (vertically downward) of the pendulum, as shown in Figure 4.4.5. Newton's second law of motion says that the product of $m$ and the tangential acceleration equals the tangential component of the gravitational force; therefore, from Figure 4.4.5,

$$
m L y^{\prime \prime}=-m g \sin y
$$

or

$$
\begin{equation*}
y^{\prime \prime}=-\frac{g}{L} \sin y \tag{4.4.13}
\end{equation*}
$$

Since $\sin n \pi=0$ if $n$ is any integer, (4.4.13) has infinitely many equilibria $\bar{y}_{n}=n \pi$. If $n$ is even, the mass is directly below the axle (Figure 4.4.6 (a)) and gravity opposes any deviation from the equilibrium. However, if $n$ is odd, the mass is directly above the axle (Figure 4.4.6 (b)) and gravity increases any deviation from the equilibrium. Therefore we conclude on physical grounds that $\bar{y}_{2 m}=2 m \pi$ is stable and $\bar{y}_{2 m+1}=(2 m+1) \pi$ is unstable.

The phase plane equivalent of (4.4.13) is

$$
v \frac{d v}{d y}=-\frac{g}{L} \sin y
$$



Figure 4.4.4 $y=R \sin \left(\omega_{0} t+\phi\right)$


Figure 4.4.5 The undamped pendulum


Figure 4.4.6 (a) Stable equilibrium (b) Unstable equilibrium
where $v=y^{\prime}$ is the angular velocity of the pendulum. Integrating this yields

$$
\begin{equation*}
\frac{v^{2}}{2}=\frac{g}{L} \cos y+c \tag{4.4.14}
\end{equation*}
$$

If $v=v_{0}$ when $y=0$, then

$$
c=\frac{v_{0}^{2}}{2}-\frac{g}{L}
$$

so (4.4.14) becomes

$$
\frac{v^{2}}{2}=\frac{v_{0}^{2}}{2}-\frac{g}{L}(1-\cos y)=\frac{v_{0}^{2}}{2}-\frac{2 g}{L} \sin ^{2} \frac{y}{2}
$$

which is equivalent to

$$
\begin{equation*}
v^{2}=v_{0}^{2}-v_{c}^{2} \sin ^{2} \frac{y}{2}, \tag{4.4.15}
\end{equation*}
$$

where

$$
v_{c}=2 \sqrt{\frac{g}{L}}
$$

The curves defined by (4.4.15) are the trajectories of (4.4.13). They are periodic with period $2 \pi$ in $y$, which isn't surprising, since if $y=y(t)$ is a solution of (4.4.13) then so is $y_{n}=y(t)+2 n \pi$ for any integer $n$. Figure 4.4 .7 shows trajectories over the interval $[-\pi, \pi]$. From (4.4.15), you can see that if $\left|v_{0}\right|>v_{c}$ then $v$ is nonzero for all $t$, which means that the object whirls in the same direction forever, as in Figure 4.4.8. The trajectories associated with this whirling motion are above the upper dashed curve and below the lower dashed curve in Figure 4.4.7. You can also see from (4.4.15) that if $0<\left|v_{0}\right|<v_{c}$,then $v=0$ when $y= \pm y_{\text {max }}$, where

$$
y_{\max }=2 \sin ^{-1}\left(\left|v_{0}\right| / v_{c}\right)
$$

In this case the pendulum oscillates periodically between $-y_{\max }$ and $y_{\max }$, as shown in Figure 4.4.9. The trajectories associated with this kind of motion are the ovals between the dashed curves in Figure 4.4.7. It can be shown (see Exercise 21 for a partial proof) that the period of the oscillation is

$$
\begin{equation*}
T=8 \int_{0}^{\pi / 2} \frac{d \theta}{\sqrt{v_{c}^{2}-v_{0}^{2} \sin ^{2} \theta}} \tag{4.4.16}
\end{equation*}
$$

Although this integral can't be evaluated in terms of familiar elementary functions, you can see that it's finite if $\left|v_{0}\right|<v_{c}$.

The dashed curves in Figure 4.4 . 7 contain four trajectories. The critical points $(\pi, 0)$ and $(-\pi, 0)$ are the trajectories of the unstable equilibrium solutions $\bar{y}= \pm \pi$. The upper dashed curve connecting (but not including) them is obtained from initial conditions of the form $y\left(t_{0}\right)=0, v\left(t_{0}\right)=v_{c}$. If $y$ is any solution with this trajectory then

$$
\lim _{t \rightarrow \infty} y(t)=\pi \quad \text { and } \quad \lim _{t \rightarrow-\infty} y(t)=-\pi
$$

The lower dashed curve connecting (but not including) them is obtained from initial conditions of the form $y\left(t_{0}\right)=0, v\left(t_{0}\right)=-v_{c}$. If $y$ is any solution with this trajectory then

$$
\lim _{t \rightarrow \infty} y(t)=-\pi \quad \text { and } \quad \lim _{t \rightarrow-\infty} y(t)=\pi
$$

Consistent with this, the integral (4.4.16) diverges to $\infty$ if $v_{0}= \pm v_{c}$. (Exercise 21).
Since the dashed curves separate trajectories of whirling solutions from trajectories of oscillating solutions, each of these curves is called a separatrix.

In general, if (4.4.7) has both stable and unstable equilibria then the separatrices are the curves given by (4.4.8) that pass through unstable critical points. Thus, if $(\bar{y}, 0)$ is an unstable critical point, then

$$
\begin{equation*}
\frac{v^{2}}{2}+P(y)=P(\bar{y}) \tag{4.4.17}
\end{equation*}
$$

defines a separatrix passing through $(\bar{y}, 0)$.


Figure 4.4.7 Trajectories of the undamped pendulum



Figure 4.4.8 The whirling undamped pendulum

Stability and Instability Conditions for $y^{\prime \prime}+p(y)=0$
It can be shown (Exercise 23) that an equilibrium $\bar{y}$ of an undamped equation

$$
\begin{equation*}
y^{\prime \prime}+p(y)=0 \tag{4.4.18}
\end{equation*}
$$

is stable if there's an open interval $(a, b)$ containing $\bar{y}$ such that

$$
\begin{equation*}
p(y)<0 \text { if } a<y<\bar{y} \text { and } p(y)>0 \text { if } \bar{y}<y<b . \tag{4.4.19}
\end{equation*}
$$

If we regard $p(y)$ as a force acting on a unit mass, (4.4.19) means that the force resists all sufficiently small displacements from $\bar{y}$.

We've already seen examples illustrating this principle. The equation (4.4.9) for the undamped springmass system is of the form (4.4.18) with $p(y)=k y / m$, which has only the stable equilibrium $\bar{y}=0$. In this case (4.4.19) holds with $a=-\infty$ and $b=\infty$. The equation (4.4.13) for the undamped pendulum is of the form (4.4.18) with $p(y)=(g / L) \sin y$. We've seen that $\bar{y}=2 m \pi$ is a stable equilibrium if $m$ is an integer. In this case

$$
p(y)=\sin y<0 \text { if }(2 m-1) \pi<y<2 m \pi
$$

and

$$
p(y)>0 \text { if } 2 m \pi<y<(2 m+1) \pi .
$$

It can also be shown (Exercise 24) that $\bar{y}$ is unstable if there's a $b>\bar{y}$ such that

$$
\begin{equation*}
p(y)<0 \text { if } \bar{y}<y<b \tag{4.4.20}
\end{equation*}
$$

or an $a<\bar{y}$ such that

$$
\begin{equation*}
p(y)>0 \text { if } a<y<\bar{y} . \tag{4.4.21}
\end{equation*}
$$

If we regard $p(y)$ as a force acting on a unit mass, (4.4.20) means that the force tends to increase all sufficiently small positive displacements from $\bar{y}$, while (4.4.21) means that the force tends to increase the magnitude of all sufficiently small negative displacements from $\bar{y}$.

The undamped pendulum also illustrates this principle. We've seen that $\bar{y}=(2 m+1) \pi$ is an unstable equilibrium if $m$ is an integer. In this case

$$
\sin y<0 \text { if }(2 m+1) \pi<y<(2 m+2) \pi,
$$

so (4.4.20) holds with $b=(2 m+2) \pi$, and

$$
\sin y>0 \text { if } 2 m \pi<y<(2 m+1) \pi
$$

so (4.4.21) holds with $a=2 m \pi$.
Example 4.4.3 The equation

$$
\begin{equation*}
y^{\prime \prime}+y(y-1)=0 \tag{4.4.22}
\end{equation*}
$$

is of the form (4.4.18) with $p(y)=y(y-1)$. Therefore $\bar{y}=0$ and $\bar{y}=1$ are the equilibria of (4.4.22). Since

$$
\begin{aligned}
y(y-1)>0 & \text { if } y<0 \text { or } y>1, \\
<0 & \text { if } 0<y<1,
\end{aligned}
$$

$\bar{y}=0$ is unstable and $\bar{y}=1$ is stable.
The phase plane equivalent of (4.4.22) is the separable equation

$$
v \frac{d v}{d y}+y(y-1)=0
$$

Integrating yields

$$
\frac{v^{2}}{2}+\frac{y^{3}}{3}-\frac{y^{2}}{2}=C
$$

which we rewrite as

$$
\begin{equation*}
v^{2}+\frac{1}{3} y^{2}(2 y-3)=c \tag{4.4.23}
\end{equation*}
$$

after renaming the constant of integration. These are the trajectories of (4.4.22). If $y$ is any solution of (4.4.22), the point $(y(t), v(t))$ moves along the trajectory of $y$ in the direction of increasing $y$ in the upper half plane ( $v=y^{\prime}>0$ ), or in the direction of decreasing $y$ in the lower half plane ( $v=y^{\prime}<0$ ).

Figure 4.4.10 shows typical trajectories. The dashed curve through the critical point $(0,0)$, obtained by setting $c=0$ in (4.4.23), separates the $y-v$ plane into regions that contain different kinds of trajectories; again, we call this curve a separatrix. Trajectories in the region bounded by the closed loop (b) are closed curves, so solutions associated with them are periodic. Solutions associated with other trajectories are not periodic. If $y$ is any such solution with trajectory not on the separatrix, then

$$
\begin{aligned}
& \lim _{t \rightarrow \infty} y(t)=-\infty, \quad \lim _{t \rightarrow-\infty} y(t)=-\infty \\
& \lim _{t \rightarrow \infty} v(t)=-\infty, \quad \lim _{t \rightarrow-\infty} v(t)=\infty
\end{aligned}
$$



Figure 4.4.10 Trajectories of $y^{\prime \prime}+y(y-1)=0$
The separatrix contains four trajectories of (4.4.22). One is the point $(0,0)$, the trajectory of the equilibrium $\bar{y}=0$. Since distinct trajectories can't intersect, the segments of the separatrix marked (a), (b), and (c) - which don't include $(0,0)$ - are distinct trajectories, none of which can be traversed in finite time. Solutions with these trajectories have the following asymptotic behavior:

$$
\begin{aligned}
\lim _{t \rightarrow \infty} y(t) & =0, \quad \lim _{t \rightarrow-\infty} y(t) \\
\lim _{t \rightarrow \infty} v(t) & =-\infty \\
\lim _{t \rightarrow-\infty} v(t) & =\infty \quad(\text { on (a)) } \\
\lim _{t \rightarrow \infty} y(t) & =0, \lim _{t \rightarrow-\infty} y(t)
\end{aligned} \quad=0, \quad(\text { on (b)) })
$$

The Damped Case
The phase plane equivalent of the damped autonomous equation

$$
\begin{equation*}
y^{\prime \prime}+q\left(y, y^{\prime}\right) y^{\prime}+p(y)=0 \tag{4.4.24}
\end{equation*}
$$

is

$$
v \frac{d v}{d y}+q(y, v) v+p(y)=0 .
$$

This equation isn't separable, so we can't solve it for $v$ in terms of $y$, as we did in the undamped case, and conservation of energy doesn't hold. (For example, energy expended in overcoming friction is lost.) However, we can study the qualitative behavior of its solutions by rewriting it as

$$
\begin{equation*}
\frac{d v}{d y}=-q(y, v)-\frac{p(y)}{v} \tag{4.4.25}
\end{equation*}
$$

and considering the direction fields for this equation. In the following examples we'll also be showing computer generated trajectories of this equation, obtained by numerical methods. The exercises call for similar computations. The methods discussed in Chapter 3 are not suitable for this task, since $p(y) / v$ in (4.4.25) is undefined on the $y$ axis of the Poincaré phase plane. Therefore we're forced to apply numerical methods briefly discussed in Section 10.1 to the system

$$
\begin{aligned}
y^{\prime} & =v \\
v^{\prime} & =-q(y, v) v-p(y)
\end{aligned}
$$

which is equivalent to (4.4.24) in the sense defined in Section 10.1. Fortunately, most differential equation software packages enable you to do this painlessly.

In the text we'll confine ourselves to the case where $q$ is constant, so (4.4.24) and (4.4.25) reduce to

$$
\begin{equation*}
y^{\prime \prime}+c y^{\prime}+p(y)=0 \tag{4.4.26}
\end{equation*}
$$

and

$$
\frac{d v}{d y}=-c-\frac{p(y)}{v}
$$

(We'll consider more general equations in the exercises.) The constant $c$ is called the damping constant. In situations where (4.4.26) is the equation of motion of an object, $c$ is positive; however, there are situations where $c$ may be negative.
The Damped Spring-Mass System
Earlier we considered the spring - mass system under the assumption that the only forces acting on the object were gravity and the spring's resistance to changes in its length. Now we'll assume that some mechanism (for example, friction in the spring or atmospheric resistance) opposes the motion of the object with a force proportional to its velocity. In Section 6.1 it will be shown that in this case Newton's second law of motion implies that

$$
\begin{equation*}
m y^{\prime \prime}+c y^{\prime}+k y=0 \tag{4.4.27}
\end{equation*}
$$

where $c>0$ is the damping constant. Again, this equation can be solved easily by a method that we'll study in Section 5.2, but that method isn't available here. Instead, we'll consider its phase plane equivalent, which can be written in the form (4.4.25) as

$$
\begin{equation*}
\frac{d v}{d y}=-\frac{c}{m}-\frac{k y}{m v} \tag{4.4.28}
\end{equation*}
$$

(A minor note: the $c$ in (4.4.26) actually corresponds to $c / m$ in this equation.) Figure 4.4.11 shows a typical direction field for an equation of this form. Recalling that motion along a trajectory must be in the direction of increasing $y$ in the upper half plane $(v>0)$ and in the direction of decreasing $y$ in the lower half plane ( $v<0$ ), you can infer that all trajectories approach the origin in clockwise fashion. To confirm this, Figure 4.4.12 shows the same direction field with some trajectories filled in. All the trajectories shown there correspond to solutions of the initial value problem

$$
m y^{\prime \prime}+c y^{\prime}+k y=0, \quad y(0)=y_{0}, \quad y^{\prime}(0)=v_{0}
$$

where

$$
m v_{0}^{2}+k y_{0}^{2}=\rho \quad(\text { a positive constant })
$$

thus, if there were no damping $(c=0)$, all the solutions would have the same dashed elliptic trajectory, shown in Figure 4.4.14.


Figure 4.4.11 A typical direction field for $m y^{\prime \prime}+c y^{\prime}+k y=0$ with $0<c<c_{1}$


Figure 4.4.12 Figure 4.4.11 with some trajectories added

Solutions corresponding to the trajectories in Figure 4.4 .12 cross the $y$-axis infinitely many times. The corresponding solutions are said to be oscillatory (Figure 4.4.13) It is shown in Section 6.2 that there's
a number $c_{1}$ such that if $0 \leq c<c_{1}$ then all solutions of (4.4.27) are oscillatory, while if $c \geq c_{1}$, no solutions of (4.4.27) have this property. (In fact, no solution not identically zero can have more than two zeros in this case.) Figure 4.4 .14 shows a direction field and some integral curves for (4.4.28) in this case.


Figure 4.4.13 An oscillatory solution of $m y^{\prime \prime}+c y^{\prime}+k y=0$

Example 4.4.4 (The Damped Pendulum) Now we return to the pendulum. If we assume that some mechanism (for example, friction in the axle or atmospheric resistance) opposes the motion of the pendulum with a force proportional to its angular velocity, Newton's second law of motion implies that

$$
\begin{equation*}
m L y^{\prime \prime}=-c y^{\prime}-m g \sin y, \tag{4.4.29}
\end{equation*}
$$

where $c>0$ is the damping constant. (Again, a minor note: the $c$ in (4.4.26) actually corresponds to $c / m L$ in this equation.) To plot a direction field for (4.4.29) we write its phase plane equivalent as

$$
\frac{d v}{d y}=-\frac{c}{m L}-\frac{g}{L v} \sin y .
$$

Figure 4.4.15 shows trajectories of four solutions of (4.4.29), all satisfying $y(0)=0$. For each $m=0,1$, 2,3 , imparting the initial velocity $v(0)=v_{m}$ causes the pendulum to make $m$ complete revolutions and then settle into decaying oscillation about the stable equilibrium $\bar{y}=2 m \pi$.


Figure 4.4.14 A typical direction field for $m y^{\prime \prime}+c y^{\prime}+k y=0$ with $c>c_{1}$


Figure 4.4.15 Four trajectories of the damped pendulum

### 4.4 Exercises

In Exercises 1-4 find the equations of the trajectories of the given undamped equation. Identify the equilibrium solutions, determine whether they are stable or unstable, and plot some trajectories. Hint: Use Eqn. (4.4.8) to obtain the equations of the trajectories.

1. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y^{3}=0$
2. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y^{2}=0$
3. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y|y|=0$
4. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y e^{-y}=0$

In Exercises 5-8 find the equations of the trajectories of the given undamped equation. Identify the equilibrium solutions, determine whether they are stable or unstable, and find the equations of the separatrices (that is, the curves through the unstable equilibria). Plot the separatrices and some trajectories in each of the regions of Poincaré plane determined by them. Hint: Use Eqn. (4.4.17) to determine the separatrices.
5. C/G $y^{\prime \prime}-y^{3}+4 y=0$
6. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y^{3}-4 y=0$
7. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y\left(y^{2}-1\right)\left(y^{2}-4\right)=0$
8. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+y(y-2)(y-1)(y+2)=0$

In Exercises 9-12 plot some trajectories of the given equation for various values (positive, negative, zero) of the parameter a. Find the equilibria of the equation and classify them as stable or unstable. Explain why the phase plane plots corresponding to positive and negative values of a differ so markedly. Can you think of a reason why zero deserves to be called the critical value of a?
9. $\mathrm{L} y^{\prime \prime}+y^{2}-a=0$
10. $\mathrm{L} y^{\prime \prime}+y^{3}-a y=0$
11. L $y^{\prime \prime}-y^{3}+a y=0$
12. L $y^{\prime \prime}+y-a y^{3}=0$

In Exercises 13-18 plot trajectories of the given equation for $c=0$ and small nonzero (positive and negative) values of $c$ to observe the effects of damping.
13. L $y^{\prime \prime}+c y^{\prime}+y^{3}=0$
14. L $y^{\prime \prime}+c y^{\prime}-y=0$
15. L $y^{\prime \prime}+c y^{\prime}+y^{3}=0$
16. $\mathrm{L} y^{\prime \prime}+c y^{\prime}+y^{2}=0$
17. L $y^{\prime \prime}+c y^{\prime}+y|y|=0$
18. L $y^{\prime \prime}+y(y-1)+c y=0$
19. L The van der Pol equation

$$
\begin{equation*}
y^{\prime \prime}-\mu\left(1-y^{2}\right) y^{\prime}+y=0 \tag{A}
\end{equation*}
$$

where $\mu$ is a positive constant and $y$ is electrical current (Section 6.3), arises in the study of an electrical circuit whose resistive properties depend upon the current. The damping term $-\mu\left(1-y^{2}\right) y^{\prime}$ works to reduce $|y|$ if $|y|<1$ or to increase $|y|$ if $|y|>1$. It can be shown that van der Pol's equation has exactly one closed trajectory, which is called a limit cycle. Trajectories inside the limit cycle spiral outward to it, while trajectories outside the limit cycle spiral inward to it (Figure 4.4.16). Use your favorite differential equations software to verify this for $\mu=.5,1.1 .5,2$. Use a grid with $-4<y<4$ and $-4<v<4$.


Figure 4.4.16 Trajectories of van der Pol's equation
20. $L$ Rayleigh's equation,

$$
y^{\prime \prime}-\mu\left(1-\left(y^{\prime}\right)^{2} / 3\right) y^{\prime}+y=0
$$

also has a limit cycle. Follow the directions of Exercise 19 for this equation.
21. In connection with Eqn (4.4.15), suppose $y(0)=0$ and $y^{\prime}(0)=v_{0}$, where $0<v_{0}<v_{c}$.
(a) Let $T_{1}$ be the time required for $y$ to increase from zero to $y_{\max }=2 \sin ^{-1}\left(v_{0} / v_{c}\right)$. Show that

$$
\begin{equation*}
\frac{d y}{d t}=\sqrt{v_{0}^{2}-v_{c}^{2} \sin ^{2} y / 2}, \quad 0 \leq t<T_{1} \tag{A}
\end{equation*}
$$

(b) Separate variables in (A) and show that

$$
\begin{equation*}
T_{1}=\int_{0}^{y_{\max }} \frac{d u}{\sqrt{v_{0}^{2}-v_{c}^{2} \sin ^{2} u / 2}} \tag{B}
\end{equation*}
$$

(c) Substitute $\sin u / 2=\left(v_{0} / v_{c}\right) \sin \theta$ in (B) to obtain

$$
\begin{equation*}
T_{1}=2 \int_{0}^{\pi / 2} \frac{d \theta}{\sqrt{v_{c}^{2}-v_{0}^{2} \sin ^{2} \theta}} \tag{C}
\end{equation*}
$$

(d) Conclude from symmetry that the time required for $(y(t), v(t))$ to traverse the trajectory

$$
v^{2}=v_{0}^{2}-v_{c}^{2} \sin ^{2} y / 2
$$

is $T=4 T_{1}$, and that consequently $y(t+T)=y(t)$ and $v(t+T)=v(t)$; that is, the oscillation is periodic with period $T$.
(e) Show that if $v_{0}=v_{c}$, the integral in (C) is improper and diverges to $\infty$. Conclude from this that $y(t)<\pi$ for all $t$ and $\lim _{t \rightarrow \infty} y(t)=\pi$.
22. Give a direct definition of an unstable equilibrium of $y^{\prime \prime}+p(y)=0$.
23. Let $p$ be continuous for all $y$ and $p(0)=0$. Suppose there's a positive number $\rho$ such that $p(y)>0$ if $0<y \leq \rho$ and $p(y)<0$ if $-\rho \leq y<0$. For $0<r \leq \rho$ let $\alpha(r)=\min \left\{\int_{0}^{r} p(x) d x, \int_{-r}^{0}|p(x)| d x\right\} \quad$ and $\quad \beta(r)=\max \left\{\int_{0}^{r} p(x) d x, \int_{-r}^{0}|p(x)| d x\right\}$.

Let $y$ be the solution of the initial value problem

$$
y^{\prime \prime}+p(y)=0, \quad y(0)=v_{0}, \quad y^{\prime}(0)=v_{0}
$$

and define $c\left(y_{0}, v_{0}\right)=v_{0}^{2}+2 \int_{0}^{y_{0}} p(x) d x$.
(a) Show that

$$
0<c\left(y_{0}, v_{0}\right)<v_{0}^{2}+2 \beta\left(\left|y_{0}\right|\right) \quad \text { if } \quad 0<\left|y_{0}\right| \leq \rho .
$$

(b) Show that

$$
v^{2}+2 \int_{0}^{y} p(x) d x=c\left(y_{0}, v_{0}\right), \quad t>0
$$

(c) Conclude from (b) that if $c\left(y_{0}, v_{0}\right)<2 \alpha(r)$ then $|y|<r, t>0$.
(d) Given $\epsilon>0$, let $\delta>0$ be chosen so that

$$
\delta^{2}+2 \beta(\delta)<\max \left\{\epsilon^{2} / 2,2 \alpha(\epsilon / \sqrt{2})\right\}
$$

Show that if $\sqrt{y_{0}^{2}+v_{0}^{2}}<\delta$ then $\sqrt{y^{2}+v^{2}}<\epsilon$ for $t>0$, which implies that $\bar{y}=0$ is a stable equilibrium of $y^{\prime \prime}+p(y)=0$.
(e) Now let $p$ be continuous for all $y$ and $p(\bar{y})=0$, where $\bar{y}$ is not necessarily zero. Suppose there's a positive number $\rho$ such that $p(y)>0$ if $\bar{y}<y \leq \bar{y}+\rho$ and $p(y)<0$ if $\bar{y}-\rho \leq$ $y<\bar{y}$. Show that $\bar{y}$ is a stable equilibrium of $y^{\prime \prime}+p(y)=0$.
24. Let $p$ be continuous for all $y$.
(a) Suppose $p(0)=0$ and there's a positive number $\rho$ such that $p(y)<0$ if $0<y \leq \rho$. Let $\epsilon$ be any number such that $0<\epsilon<\rho$. Show that if $y$ is the solution of the initial value problem

$$
y^{\prime \prime}+p(y)=0, \quad y(0)=y_{0}, \quad y^{\prime}(0)=0
$$

with $0<y_{0}<\epsilon$, then $y(t) \geq \epsilon$ for some $t>0$. Conclude that $\bar{y}=0$ is an unstable equilibrium of $y^{\prime \prime}+p(y)=0$. Hint: Let $k=\min _{y_{0} \leq x \leq \epsilon}(-p(x))$, which is positive. Show that if $y(t)<\epsilon$ for $0 \leq t<T$ then $k T^{2}<2\left(\epsilon-y_{0}\right)$.
(b) Now let $p(\bar{y})=0$, where $\bar{y}$ isn't necessarily zero. Suppose there's a positive number $\rho$ such that $p(y)<0$ if $\bar{y}<y \leq \bar{y}+\rho$. Show that $\bar{y}$ is an unstable equilibrium of $y^{\prime \prime}+p(y)=0$.
(c) Modify your proofs of (a) and (b) to show that if there's a positive number $\rho$ such that $p(y)>0$ if $\bar{y}-\rho \leq y<\bar{y}$, then $\bar{y}$ is an unstable equilibrium of $y^{\prime \prime}+p(y)=0$.

## CHAPTER 5 Linear Second Order Equations

IN THIS CHAPTER we study a particularly important class of second order equations. Because of their many applications in science and engineering, second order differential equation have historically been the most thoroughly studied class of differential equations. Research on the theory of second order differential equations continues to the present day. This chapter is devoted to second order equations that can be written in the form

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F(x)
$$

Such equations are said to be linear. As in the case of first order linear equations, (A) is said to be homogeneous if $F \equiv 0$, or nonhomogeneous if $F \not \equiv 0$.

SECTION 5.1 is devoted to the theory of homogeneous linear equations.
SECTION 5.2 deals with homogeneous equations of the special form

$$
a y^{\prime \prime}+b y^{\prime}+c y=0
$$

where $a, b$, and $c$ are constant $(a \neq 0)$. When you've completed this section you'll know everything there is to know about solving such equations.

SECTION 5.3 presents the theory of nonhomogeneous linear equations.
SECTIONS 5.4 AND 5.5 present the method of undetermined coefficients, which can be used to solve nonhomogeneous equations of the form

$$
a y^{\prime \prime}+b y^{\prime}+c y=F(x),
$$

where $a, b$, and $c$ are constants and $F$ has a special form that is still sufficiently general to occur in many applications. In this section we make extensive use of the idea of variation of parameters introduced in Chapter 2.

SECTION 5.6 deals with reduction of order, a technique based on the idea of variation of parameters, which enables us to find the general solution of a nonhomogeneous linear second order equation provided that we know one nontrivial (not identically zero) solution of the associated homogeneous equation.

SECTION 5.7 deals with the method traditionally called variation of parameters, which enables us to find the general solution of a nonhomogeneous linear second order equation provided that we know two nontrivial solutions (with nonconstant ratio) of the associated homogeneous equation.

### 5.1 HOMOGENEOUS LINEAR EQUATIONS

A second order differential equation is said to be linear if it can be written as

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f(x) . \tag{5.1.1}
\end{equation*}
$$

We call the function $f$ on the right a forcing function, since in physical applications it's often related to a force acting on some system modeled by the differential equation. We say that (5.1.1) is homogeneous if $f \equiv 0$ or nonhomogeneous if $f \not \equiv 0$. Since these definitions are like the corresponding definitions in Section 2.1 for the linear first order equation

$$
\begin{equation*}
y^{\prime}+p(x) y=f(x) \tag{5.1.2}
\end{equation*}
$$

it's natural to expect similarities between methods of solving (5.1.1) and (5.1.2). However, solving (5.1.1) is more difficult than solving (5.1.2). For example, while Theorem 2.1.1 gives a formula for the general solution of (5.1.2) in the case where $f \equiv 0$ and Theorem 2.1.2 gives a formula for the case where $f \not \equiv 0$, there are no formulas for the general solution of (5.1.1) in either case. Therefore we must be content to solve linear second order equations of special forms.

In Section 2.1 we considered the homogeneous equation $y^{\prime}+p(x) y=0$ first, and then used a nontrivial solution of this equation to find the general solution of the nonhomogeneous equation $y^{\prime}+p(x) y=f(x)$. Although the progression from the homogeneous to the nonhomogeneous case isn't that simple for the linear second order equation, it's still necessary to solve the homogeneous equation

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.1.3}
\end{equation*}
$$

in order to solve the nonhomogeneous equation (5.1.1). This section is devoted to (5.1.3).
The next theorem gives sufficient conditions for existence and uniqueness of solutions of initial value problems for (5.1.3). We omit the proof.

Theorem 5.1.1 Suppose $p$ and $q$ are continuous on an open interval $(a, b)$, let $x_{0}$ be any point in $(a, b)$, and let $k_{0}$ and $k_{1}$ be arbitrary real numbers. Then the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0, y\left(x_{0}\right)=k_{0}, y^{\prime}\left(x_{0}\right)=k_{1}
$$

has a unique solution on $(a, b)$.
Since $y \equiv 0$ is obviously a solution of (5.1.3) we call it the trivial solution. Any other solution is nontrivial. Under the assumptions of Theorem 5.1.1, the only solution of the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0, y\left(x_{0}\right)=0, y^{\prime}\left(x_{0}\right)=0
$$

on $(a, b)$ is the trivial solution (Exercise 24).
The next three examples illustrate concepts that we'll develop later in this section. You shouldn't be concerned with how to find the given solutions of the equations in these examples. This will be explained in later sections.

Example 5.1.1 The coefficients of $y^{\prime}$ and $y$ in

$$
\begin{equation*}
y^{\prime \prime}-y=0 \tag{5.1.4}
\end{equation*}
$$

are the constant functions $p \equiv 0$ and $q \equiv-1$, which are continuous on $(-\infty, \infty)$. Therefore Theorem 5.1.1 implies that every initial value problem for (5.1.4) has a unique solution on $(-\infty, \infty)$.
(a) Verify that $y_{1}=e^{x}$ and $y_{2}=e^{-x}$ are solutions of (5.1.4) on $(-\infty, \infty)$.
(b) Verify that if $c_{1}$ and $c_{2}$ are arbitrary constants, $y=c_{1} e^{x}+c_{2} e^{-x}$ is a solution of (5.1.4) on $(-\infty, \infty)$.
(c) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}-y=0, \quad y(0)=1, \quad y^{\prime}(0)=3 . \tag{5.1.5}
\end{equation*}
$$

$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ If $y_{1}=e^{x}$ then $y_{1}^{\prime}=e^{x}$ and $y_{1}^{\prime \prime}=e^{x}=y_{1}$, so $y_{1}^{\prime \prime}-y_{1}=0$. If $y_{2}=e^{-x}$, then $y_{2}^{\prime}=-e^{-x}$ and $y_{2}^{\prime \prime}=e^{-x}=y_{2}$, so $y_{2}^{\prime \prime}-y_{2}=0$.
$\underline{\text { SOLUTION(b) If }}$

$$
\begin{equation*}
y=c_{1} e^{x}+c_{2} e^{-x} \tag{5.1.6}
\end{equation*}
$$

then

$$
\begin{equation*}
y^{\prime}=c_{1} e^{x}-c_{2} e^{-x} \tag{5.1.7}
\end{equation*}
$$

and

$$
y^{\prime \prime}=c_{1} e^{x}+c_{2} e^{-x}
$$

so

$$
\begin{aligned}
y^{\prime \prime}-y & =\left(c_{1} e^{x}+c_{2} e^{-x}\right)-\left(c_{1} e^{x}+c_{2} e^{-x}\right) \\
& =c_{1}\left(e^{x}-e^{x}\right)+c_{2}\left(e^{-x}-e^{-x}\right)=0
\end{aligned}
$$

for all $x$. Therefore $y=c_{1} e^{x}+c_{2} e^{-x}$ is a solution of (5.1.4) on $(-\infty, \infty)$.
 Setting $x=0$ in (5.1.6) and (5.1.7) shows that this is equivalent to

$$
\begin{aligned}
& c_{1}+c_{2}=1 \\
& c_{1}-c_{2}=3
\end{aligned}
$$

Solving these equations yields $c_{1}=2$ and $c_{2}=-1$. Therefore $y=2 e^{x}-e^{-x}$ is the unique solution of (5.1.5) on $(-\infty, \infty)$.

Example 5.1.2 Let $\omega$ be a positive constant. The coefficients of $y^{\prime}$ and $y$ in

$$
\begin{equation*}
y^{\prime \prime}+\omega^{2} y=0 \tag{5.1.8}
\end{equation*}
$$

are the constant functions $p \equiv 0$ and $q \equiv \omega^{2}$, which are continuous on $(-\infty, \infty)$. Therefore Theorem 5.1.1 implies that every initial value problem for (5.1.8) has a unique solution on $(-\infty, \infty)$.
(a) Verify that $y_{1}=\cos \omega x$ and $y_{2}=\sin \omega x$ are solutions of (5.1.8) on $(-\infty, \infty)$.
(b) Verify that if $c_{1}$ and $c_{2}$ are arbitrary constants then $y=c_{1} \cos \omega x+c_{2} \sin \omega x$ is a solution of (5.1.8) on $(-\infty, \infty)$.
(c) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+\omega^{2} y=0, \quad y(0)=1, \quad y^{\prime}(0)=3 \tag{5.1.9}
\end{equation*}
$$

SOLUTION(a) If $y_{1}=\cos \omega x$ then $y_{1}^{\prime}=-\omega \sin \omega x$ and $y_{1}^{\prime \prime}=-\omega^{2} \cos \omega x=-\omega^{2} y_{1}$, so $y_{1}^{\prime \prime}+\omega^{2} y_{1}=0$. If $y_{2}=\sin \omega x$ then, $y_{2}^{\prime}=\omega \cos \omega x$ and $y_{2}^{\prime \prime}=-\omega^{2} \sin \omega x=-\omega^{2} y_{2}$, so $y_{2}^{\prime \prime}+\omega^{2} y_{2}=0$.

## $\underline{\text { SOLUTION(b) If }}$

$$
\begin{equation*}
y=c_{1} \cos \omega x+c_{2} \sin \omega x \tag{5.1.10}
\end{equation*}
$$

then

$$
\begin{equation*}
y^{\prime}=\omega\left(-c_{1} \sin \omega x+c_{2} \cos \omega x\right) \tag{5.1.11}
\end{equation*}
$$

and

$$
y^{\prime \prime}=-\omega^{2}\left(c_{1} \cos \omega x+c_{2} \sin \omega x\right)
$$

so

$$
\begin{aligned}
y^{\prime \prime}+\omega^{2} y & =-\omega^{2}\left(c_{1} \cos \omega x+c_{2} \sin \omega x\right)+\omega^{2}\left(c_{1} \cos \omega x+c_{2} \sin \omega x\right) \\
& =c_{1} \omega^{2}(-\cos \omega x+\cos \omega x)+c_{2} \omega^{2}(-\sin \omega x+\sin \omega x)=0
\end{aligned}
$$

for all $x$. Therefore $y=c_{1} \cos \omega x+c_{2} \sin \omega x$ is a solution of (5.1.8) on $(-\infty, \infty)$.
SOLUTION(c) To solve (5.1.9), we must choosing $c_{1}$ and $c_{2}$ in (5.1.10) so that $y(0)=1$ and $y^{\prime}(0)=3$. Setting $x=0$ in (5.1.10) and (5.1.11) shows that $c_{1}=1$ and $c_{2}=3 / \omega$. Therefore

$$
y=\cos \omega x+\frac{3}{\omega} \sin \omega x
$$

is the unique solution of (5.1.9) on $(-\infty, \infty)$.
Theorem 5.1.1 implies that if $k_{0}$ and $k_{1}$ are arbitrary real numbers then the initial value problem

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1} \tag{5.1.12}
\end{equation*}
$$

has a unique solution on an interval $(a, b)$ that contains $x_{0}$, provided that $P_{0}, P_{1}$, and $P_{2}$ are continuous and $P_{0}$ has no zeros on $(a, b)$. To see this, we rewrite the differential equation in (5.1.12) as

$$
y^{\prime \prime}+\frac{P_{1}(x)}{P_{0}(x)} y^{\prime}+\frac{P_{2}(x)}{P_{0}(x)} y=0
$$

and apply Theorem 5.1.1 with $p=P_{1} / P_{0}$ and $q=P_{2} / P_{0}$.
Example 5.1.3 The equation

$$
\begin{equation*}
x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0 \tag{5.1.13}
\end{equation*}
$$

has the form of the differential equation in (5.1.12), with $P_{0}(x)=x^{2}, P_{1}(x)=x$, and $P_{2}(x)=-4$, which are are all continuous on $(-\infty, \infty)$. However, since $P(0)=0$ we must consider solutions of (5.1.13) on $(-\infty, 0)$ and $(0, \infty)$. Since $P_{0}$ has no zeros on these intervals, Theorem 5.1.1 implies that the initial value problem

$$
x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

has a unique solution on $(0, \infty)$ if $x_{0}>0$, or on $(-\infty, 0)$ if $x_{0}<0$.
(a) Verify that $y_{1}=x^{2}$ is a solution of (5.1.13) on $(-\infty, \infty)$ and $y_{2}=1 / x^{2}$ is a solution of (5.1.13) on $(-\infty, 0)$ and $(0, \infty)$.
(b) Verify that if $c_{1}$ and $c_{2}$ are any constants then $y=c_{1} x^{2}+c_{2} / x^{2}$ is a solution of (5.1.13) on $(-\infty, 0)$ and $(0, \infty)$.
(c) Solve the initial value problem

$$
\begin{equation*}
x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0, \quad y(1)=2, \quad y^{\prime}(1)=0 \tag{5.1.14}
\end{equation*}
$$

(d) Solve the initial value problem

$$
\begin{equation*}
x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0, \quad y(-1)=2, \quad y^{\prime}(-1)=0 \tag{5.1.15}
\end{equation*}
$$

$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ If $y_{1}=x^{2}$ then $y_{1}^{\prime}=2 x$ and $y_{1}^{\prime \prime}=2$, so

$$
x^{2} y_{1}^{\prime \prime}+x y_{1}^{\prime}-4 y_{1}=x^{2}(2)+x(2 x)-4 x^{2}=0
$$

for $x$ in $(-\infty, \infty)$. If $y_{2}=1 / x^{2}$, then $y_{2}^{\prime}=-2 / x^{3}$ and $y_{2}^{\prime \prime}=6 / x^{4}$, so

$$
x^{2} y_{2}^{\prime \prime}+x y_{2}^{\prime}-4 y_{2}=x^{2}\left(\frac{6}{x^{4}}\right)-x\left(\frac{2}{x^{3}}\right)-\frac{4}{x^{2}}=0
$$

for $x$ in $(-\infty, 0)$ or $(0, \infty)$.
$\underline{\text { SOLUTION(b) If }}$

$$
\begin{equation*}
y=c_{1} x^{2}+\frac{c_{2}}{x^{2}} \tag{5.1.16}
\end{equation*}
$$

then

$$
\begin{equation*}
y^{\prime}=2 c_{1} x-\frac{2 c_{2}}{x^{3}} \tag{5.1.17}
\end{equation*}
$$

and

$$
y^{\prime \prime}=2 c_{1}+\frac{6 c_{2}}{x^{4}},
$$

SO

$$
\begin{aligned}
x^{2} y^{\prime \prime}+x y^{\prime}-4 y & =x^{2}\left(2 c_{1}+\frac{6 c_{2}}{x^{4}}\right)+x\left(2 c_{1} x-\frac{2 c_{2}}{x^{3}}\right)-4\left(c_{1} x^{2}+\frac{c_{2}}{x^{2}}\right) \\
& =c_{1}\left(2 x^{2}+2 x^{2}-4 x^{2}\right)+c_{2}\left(\frac{6}{x^{2}}-\frac{2}{x^{2}}-\frac{4}{x^{2}}\right) \\
& =c_{1} \cdot 0+c_{2} \cdot 0=0
\end{aligned}
$$

for $x$ in $(-\infty, 0)$ or $(0, \infty)$.
 $x=1$ in (5.1.16) and (5.1.17) shows that this is equivalent to

$$
\begin{aligned}
c_{1}+c_{2} & =2 \\
2 c_{1}-2 c_{2} & =0 .
\end{aligned}
$$

Solving these equations yields $c_{1}=1$ and $c_{2}=1$. Therefore $y=x^{2}+1 / x^{2}$ is the unique solution of (5.1.14) on $(0, \infty)$.
$\underline{\text { SOLUTION(d) }}$ We can solve (5.1.15) by choosing $c_{1}$ and $c_{2}$ in (5.1.16) so that $y(-1)=2$ and $y^{\prime}(-1)=$ 0 . Setting $x=-1$ in (5.1.16) and (5.1.17) shows that this is equivalent to

$$
\begin{aligned}
c_{1}+c_{2} & =2 \\
-2 c_{1}+2 c_{2} & =0 .
\end{aligned}
$$

Solving these equations yields $c_{1}=1$ and $c_{2}=1$. Therefore $y=x^{2}+1 / x^{2}$ is the unique solution of (5.1.15) on $(-\infty, 0)$.

Although the formulas for the solutions of (5.1.14) and (5.1.15) are both $y=x^{2}+1 / x^{2}$, you should not conclude that these two initial value problems have the same solution. Remember that a solution of an initial value problem is defined on an interval that contains the initial point; therefore, the solution of (5.1.14) is $y=x^{2}+1 / x^{2}$ on the interval $(0, \infty)$, which contains the initial point $x_{0}=1$, while the solution of (5.1.15) is $y=x^{2}+1 / x^{2}$ on the interval $(-\infty, 0)$, which contains the initial point $x_{0}=-1$.

The General Solution of a Homogeneous Linear Second Order Equation
If $y_{1}$ and $y_{2}$ are defined on an interval $(a, b)$ and $c_{1}$ and $c_{2}$ are constants, then

$$
y=c_{1} y_{1}+c_{2} y_{2}
$$

is a linear combination of $y_{1}$ and $y_{2}$. For example, $y=2 \cos x+7 \sin x$ is a linear combination of $y_{1}=\cos x$ and $y_{2}=\sin x$, with $c_{1}=2$ and $c_{2}=7$.

The next theorem states a fact that we've already verified in Examples 5.1.1, 5.1.2, and 5.1.3.
Theorem 5.1.2 If $y_{1}$ and $y_{2}$ are solutions of the homogeneous equation

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.1.18}
\end{equation*}
$$

on $(a, b)$, then any linear combination

$$
\begin{equation*}
y=c_{1} y_{1}+c_{2} y_{2} \tag{5.1.19}
\end{equation*}
$$

of $y_{1}$ and $y_{2}$ is also a solution of (5.1.18) on $(a, b)$.
Proof If

$$
y=c_{1} y_{1}+c_{2} y_{2}
$$

then

$$
y^{\prime}=c_{1} y_{1}^{\prime}+c_{2} y_{2}^{\prime} \quad \text { and } \quad y^{\prime \prime}=c_{1} y_{1}^{\prime \prime}+c_{2} y_{2}^{\prime \prime}
$$

Therefore

$$
\begin{aligned}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y & =\left(c_{1} y_{1}^{\prime \prime}+c_{2} y_{2}^{\prime \prime}\right)+p(x)\left(c_{1} y_{1}^{\prime}+c_{2} y_{2}^{\prime}\right)+q(x)\left(c_{1} y_{1}+c_{2} y_{2}\right) \\
& =c_{1}\left(y_{1}^{\prime \prime}+p(x) y_{1}^{\prime}+q(x) y_{1}\right)+c_{2}\left(y_{2}^{\prime \prime}+p(x) y_{2}^{\prime}+q(x) y_{2}\right) \\
& =c_{1} \cdot 0+c_{2} \cdot 0=0
\end{aligned}
$$

since $y_{1}$ and $y_{2}$ are solutions of (5.1.18).
We say that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.1.18) on $(a, b)$ if every solution of (5.1.18) on $(a, b)$ can be written as a linear combination of $y_{1}$ and $y_{2}$ as in (5.1.19). In this case we say that (5.1.19) is general solution of (5.1.18) on $(a, b)$.

Linear Independence

We need a way to determine whether a given set $\left\{y_{1}, y_{2}\right\}$ of solutions of (5.1.18) is a fundamental set. The next definition will enable us to state necessary and sufficient conditions for this.

We say that two functions $y_{1}$ and $y_{2}$ defined on an interval $(a, b)$ are linearly independent on $(a, b)$ if neither is a constant multiple of the other on $(a, b)$. (In particular, this means that neither can be the trivial solution of (5.1.18), since, for example, if $y_{1} \equiv 0$ we could write $y_{1}=0 y_{2}$.) We'll also say that the set $\left\{y_{1}, y_{2}\right\}$ is linearly independent on $(a, b)$.

Theorem 5.1.3 Suppose $p$ and $q$ are continuous on $(a, b)$. Then a set $\left\{y_{1}, y_{2}\right\}$ of solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.1.20}
\end{equation*}
$$

on $(a, b)$ is a fundamental set if and only if $\left\{y_{1}, y_{2}\right\}$ is linearly independent on $(a, b)$.
We'll present the proof of Theorem 5.1.3 in steps worth regarding as theorems in their own right. However, let's first interpret Theorem 5.1.3 in terms of Examples 5.1.1, 5.1.2, and 5.1.3.

## Example 5.1.4

(a) Since $e^{x} / e^{-x}=e^{2 x}$ is nonconstant, Theorem 5.1.3 implies that $y=c_{1} e^{x}+c_{2} e^{-x}$ is the general solution of $y^{\prime \prime}-y=0$ on $(-\infty, \infty)$.
(b) Since $\cos \omega x / \sin \omega x=\cot \omega x$ is nonconstant, Theorem 5.1.3 implies that $y=c_{1} \cos \omega x+$ $c_{2} \sin \omega x$ is the general solution of $y^{\prime \prime}+\omega^{2} y=0$ on $(-\infty, \infty)$.
(c) Since $x^{2} / x^{-2}=x^{4}$ is nonconstant, Theorem 5.1.3 implies that $y=c_{1} x^{2}+c_{2} / x^{2}$ is the general solution of $x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0$ on $(-\infty, 0)$ and $(0, \infty)$.

The Wronskian and Abel's Formula
To motivate a result that we need in order to prove Theorem 5.1.3, let's see what is required to prove that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.1.20) on $(a, b)$. Let $x_{0}$ be an arbitrary point in $(a, b)$, and suppose $y$ is an arbitrary solution of $(5.1 .20)$ on $(a, b)$. Then $y$ is the unique solution of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1} \tag{5.1.21}
\end{equation*}
$$

that is, $k_{0}$ and $k_{1}$ are the numbers obtained by evaluating $y$ and $y^{\prime}$ at $x_{0}$. Moreover, $k_{0}$ and $k_{1}$ can be any real numbers, since Theorem 5.1.1 implies that (5.1.21) has a solution no matter how $k_{0}$ and $k_{1}$ are chosen. Therefore $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.1.20) on $(a, b)$ if and only if it's possible to write the solution of an arbitrary initial value problem (5.1.21) as $y=c_{1} y_{1}+c_{2} y_{2}$. This is equivalent to requiring that the system

$$
\begin{align*}
& c_{1} y_{1}\left(x_{0}\right)+c_{2} y_{2}\left(x_{0}\right)=k_{0} \\
& c_{1} y_{1}^{\prime}\left(x_{0}\right)+c_{2} y_{2}^{\prime}\left(x_{0}\right)=k_{1} \tag{5.1.22}
\end{align*}
$$

has a solution $\left(c_{1}, c_{2}\right)$ for every choice of $\left(k_{0}, k_{1}\right)$. Let's try to solve (5.1.22).
Multiplying the first equation in (5.1.22) by $y_{2}^{\prime}\left(x_{0}\right)$ and the second by $y_{2}\left(x_{0}\right)$ yields

$$
\begin{aligned}
& c_{1} y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)+c_{2} y_{2}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)=y_{2}^{\prime}\left(x_{0}\right) k_{0} \\
& c_{1} y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)+c_{2} y_{2}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)=y_{2}\left(x_{0}\right) k_{1}
\end{aligned}
$$

and subtracting the second equation here from the first yields

$$
\begin{equation*}
\left(y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)\right) c_{1}=y_{2}^{\prime}\left(x_{0}\right) k_{0}-y_{2}\left(x_{0}\right) k_{1} . \tag{5.1.23}
\end{equation*}
$$

Multiplying the first equation in (5.1.22) by $y_{1}^{\prime}\left(x_{0}\right)$ and the second by $y_{1}\left(x_{0}\right)$ yields

$$
\begin{aligned}
& c_{1} y_{1}\left(x_{0}\right) y_{1}^{\prime}\left(x_{0}\right)+c_{2} y_{2}\left(x_{0}\right) y_{1}^{\prime}\left(x_{0}\right)=y_{1}^{\prime}\left(x_{0}\right) k_{0} \\
& c_{1} y_{1}^{\prime}\left(x_{0}\right) y_{1}\left(x_{0}\right)+c_{2} y_{2}^{\prime}\left(x_{0}\right) y_{1}\left(x_{0}\right)=y_{1}\left(x_{0}\right) k_{1}
\end{aligned}
$$

and subtracting the first equation here from the second yields

$$
\begin{equation*}
\left(y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)\right) c_{2}=y_{1}\left(x_{0}\right) k_{1}-y_{1}^{\prime}\left(x_{0}\right) k_{0} . \tag{5.1.24}
\end{equation*}
$$

If

$$
y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)=0,
$$

it's impossible to satisfy (5.1.23) and (5.1.24) (and therefore (5.1.22)) unless $k_{0}$ and $k_{1}$ happen to satisfy

$$
\begin{aligned}
y_{1}\left(x_{0}\right) k_{1}-y_{1}^{\prime}\left(x_{0}\right) k_{0} & =0 \\
y_{2}^{\prime}\left(x_{0}\right) k_{0}-y_{2}\left(x_{0}\right) k_{1} & =0 .
\end{aligned}
$$

On the other hand, if

$$
\begin{equation*}
y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right) \neq 0 \tag{5.1.25}
\end{equation*}
$$

we can divide (5.1.23) and (5.1.24) through by the quantity on the left to obtain

$$
\begin{align*}
c_{1} & =\frac{y_{2}^{\prime}\left(x_{0}\right) k_{0}-y_{2}\left(x_{0}\right) k_{1}}{y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)}  \tag{5.1.26}\\
c_{2} & =\frac{y_{1}\left(x_{0}\right) k_{1}-y_{1}^{\prime}\left(x_{0}\right) k_{0}}{y_{1}\left(x_{0}\right) y_{2}^{\prime}\left(x_{0}\right)-y_{1}^{\prime}\left(x_{0}\right) y_{2}\left(x_{0}\right)},
\end{align*}
$$

no matter how $k_{0}$ and $k_{1}$ are chosen. This motivates us to consider conditions on $y_{1}$ and $y_{2}$ that imply (5.1.25).

Theorem 5.1.4 Suppose $p$ and $q$ are continuous on $(a, b)$, let $y_{1}$ and $y_{2}$ be solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.1.27}
\end{equation*}
$$

on $(a, b)$, and define

$$
\begin{equation*}
W=y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2} \tag{5.1.28}
\end{equation*}
$$

Let $x_{0}$ be any point in $(a, b)$. Then

$$
\begin{equation*}
W(x)=W\left(x_{0}\right) e^{-\int_{x_{0}}^{x} p(t) d t}, \quad a<x<b \tag{5.1.29}
\end{equation*}
$$

Therefore either $W$ has no zeros in $(a, b)$ or $W \equiv 0$ on $(a, b)$.
Proof Differentiating (5.1.28) yields

$$
\begin{equation*}
W^{\prime}=y_{1}^{\prime} y_{2}^{\prime}+y_{1} y_{2}^{\prime \prime}-y_{1}^{\prime} y_{2}^{\prime}-y_{1}^{\prime \prime} y_{2}=y_{1} y_{2}^{\prime \prime}-y_{1}^{\prime \prime} y_{2} . \tag{5.1.30}
\end{equation*}
$$

Since $y_{1}$ and $y_{2}$ both satisfy (5.1.27),

$$
y_{1}^{\prime \prime}=-p y_{1}^{\prime}-q y_{1} \quad \text { and } \quad y_{2}^{\prime \prime}=-p y_{2}^{\prime}-q y_{2} .
$$

Substituting these into (5.1.30) yields

$$
\begin{aligned}
W^{\prime} & =-y_{1}\left(p y_{2}^{\prime}+q y_{2}\right)+y_{2}\left(p y_{1}^{\prime}+q y_{1}\right) \\
& =-p\left(y_{1} y_{2}^{\prime}-y_{2} y_{1}^{\prime}\right)-q\left(y_{1} y_{2}-y_{2} y_{1}\right) \\
& =-p\left(y_{1} y_{2}^{\prime}-y_{2} y_{1}^{\prime}\right)=-p W .
\end{aligned}
$$

Therefore $W^{\prime}+p(x) W=0$; that is, $W$ is the solution of the initial value problem

$$
y^{\prime}+p(x) y=0, \quad y\left(x_{0}\right)=W\left(x_{0}\right)
$$

We leave it to you to verify by separation of variables that this implies (5.1.29). If $W\left(x_{0}\right) \neq 0$, (5.1.29) implies that $W$ has no zeros in $(a, b)$, since an exponential is never zero. On the other hand, if $W\left(x_{0}\right)=0$, (5.1.29) implies that $W(x)=0$ for all $x$ in $(a, b)$.

The function $W$ defined in (5.1.28) is the Wronskian of $\left\{y_{1}, y_{2}\right\}$. Formula (5.1.29) is Abel's formula.
The Wronskian of $\left\{y_{1}, y_{2}\right\}$ is usually written as the determinant

$$
W=\left|\begin{array}{ll}
y_{1} & y_{2} \\
y_{1}^{\prime} & y_{2}^{\prime}
\end{array}\right| .
$$

The expressions in (5.1.26) for $c_{1}$ and $c_{2}$ can be written in terms of determinants as

$$
c_{1}=\frac{1}{W\left(x_{0}\right)}\left|\begin{array}{ll}
k_{0} & y_{2}\left(x_{0}\right) \\
k_{1} & y_{2}^{\prime}\left(x_{0}\right)
\end{array}\right| \quad \text { and } \quad c_{2}=\frac{1}{W\left(x_{0}\right)}\left|\begin{array}{cc}
y_{1}\left(x_{0}\right) & k_{0} \\
y_{1}^{\prime}\left(x_{0}\right) & k_{1}
\end{array}\right| .
$$

If you've taken linear algebra you may recognize this as Cramer's rule.

Example 5.1.5 Verify Abel's formula for the following differential equations and the corresponding solutions, from Examples 5.1.1, 5.1.2, and 5.1.3:
(a) $y^{\prime \prime}-y=0 ; \quad y_{1}=e^{x}, y_{2}=e^{-x}$
(b) $y^{\prime \prime}+\omega^{2} y=0 ; \quad y_{1}=\cos \omega x, y_{2}=\sin \omega x$
(c) $x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0 ; \quad y_{1}=x^{2}, y_{2}=1 / x^{2}$
$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ Since $p \equiv 0$, we can verify Abel's formula by showing that $W$ is constant, which is true, since

$$
W(x)=\left|\begin{array}{rr}
e^{x} & e^{-x} \\
e^{x} & -e^{-x}
\end{array}\right|=e^{x}\left(-e^{-x}\right)-e^{x} e^{-x}=-2
$$

for all $x$.
$\underline{\text { SOLUTION(b) Again, since } p \equiv 0 \text {, we can verify Abel's formula by showing that } W \text { is constant, which }}$ is true, since

$$
\begin{aligned}
W(x) & =\left|\begin{array}{cc}
\cos \omega x & \sin \omega x \\
-\omega \sin \omega x & \omega \cos \omega x
\end{array}\right| \\
& =\cos \omega x(\omega \cos \omega x)-(-\omega \sin \omega x) \sin \omega x \\
& =\omega\left(\cos ^{2} \omega x+\sin ^{2} \omega x\right)=\omega
\end{aligned}
$$

for all $x$.
$\underline{\text { SOLUTION(c) Computing the Wronskian of } y_{1}=x^{2} \text { and } y_{2}=1 / x^{2} \text { directly yields }}$

$$
W=\left|\begin{array}{cc}
x^{2} & 1 / x^{2}  \tag{5.1.31}\\
2 x & -2 / x^{3}
\end{array}\right|=x^{2}\left(-\frac{2}{x^{3}}\right)-2 x\left(\frac{1}{x^{2}}\right)=-\frac{4}{x} .
$$

To verify Abel's formula we rewrite the differential equation as

$$
y^{\prime \prime}+\frac{1}{x} y^{\prime}-\frac{4}{x^{2}} y=0
$$

to see that $p(x)=1 / x$. If $x_{0}$ and $x$ are either both in $(-\infty, 0)$ or both in $(0, \infty)$ then

$$
\int_{x_{0}}^{x} p(t) d t=\int_{x_{0}}^{x} \frac{d t}{t}=\ln \left(\frac{x}{x_{0}}\right)
$$

so Abel's formula becomes

$$
\begin{aligned}
W(x) & =W\left(x_{0}\right) e^{-\ln \left(x / x_{0}\right)}=W\left(x_{0}\right) \frac{x_{0}}{x} \\
& =-\left(\frac{4}{x_{0}}\right)\left(\frac{x_{0}}{x}\right) \quad \text { from (5.1.31) } \\
& =-\frac{4}{x}
\end{aligned}
$$

which is consistent with (5.1.31).
The next theorem will enable us to complete the proof of Theorem 5.1.3.
Theorem 5.1.5 Suppose $p$ and $q$ are continuous on an open interval $(a, b)$, let $y_{1}$ and $y_{2}$ be solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.1.32}
\end{equation*}
$$

on $(a, b)$, and let $W=y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}$. Then $y_{1}$ and $y_{2}$ are linearly independent on $(a, b)$ if and only if $W$ has no zeros on $(a, b)$.

Proof We first show that if $W\left(x_{0}\right)=0$ for some $x_{0}$ in $(a, b)$, then $y_{1}$ and $y_{2}$ are linearly dependent on $(a, b)$. Let $I$ be a subinterval of $(a, b)$ on which $y_{1}$ has no zeros. (If there's no such subinterval, $y_{1} \equiv 0$ on
$(a, b)$, so $y_{1}$ and $y_{2}$ are linearly independent, and we're finished with this part of the proof.) Then $y_{2} / y_{1}$ is defined on $I$, and

$$
\begin{equation*}
\left(\frac{y_{2}}{y_{1}}\right)^{\prime}=\frac{y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}}{y_{1}^{2}}=\frac{W}{y_{1}^{2}} . \tag{5.1.33}
\end{equation*}
$$

However, if $W\left(x_{0}\right)=0$, Theorem 5.1.4 implies that $W \equiv 0$ on $(a, b)$. Therefore (5.1.33) implies that $\left(y_{2} / y_{1}\right)^{\prime} \equiv 0$, so $y_{2} / y_{1}=c$ (constant) on $I$. This shows that $y_{2}(x)=c y_{1}(x)$ for all $x$ in $I$. However, we want to show that $y_{2}=c y_{1}(x)$ for all $x$ in $(a, b)$. Let $Y=y_{2}-c y_{1}$. Then $Y$ is a solution of (5.1.32) on $(a, b)$ such that $Y \equiv 0$ on $I$, and therefore $Y^{\prime} \equiv 0$ on $I$. Consequently, if $x_{0}$ is chosen arbitrarily in $I$ then $Y$ is a solution of the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0, \quad y\left(x_{0}\right)=0, \quad y^{\prime}\left(x_{0}\right)=0
$$

which implies that $Y \equiv 0$ on $(a, b)$, by the paragraph following Theorem 5.1.1. (See also Exercise 24). Hence, $y_{2}-c y_{1} \equiv 0$ on $(a, b)$, which implies that $y_{1}$ and $y_{2}$ are not linearly independent on $(a, b)$.

Now suppose $W$ has no zeros on $(a, b)$. Then $y_{1}$ can't be identically zero on ( $a, b$ ) (why not?), and therefore there is a subinterval $I$ of $(a, b)$ on which $y_{1}$ has no zeros. Since (5.1.33) implies that $y_{2} / y_{1}$ is nonconstant on $I$, $y_{2}$ isn't a constant multiple of $y_{1}$ on $(a, b)$. A similar argument shows that $y_{1}$ isn't a constant multiple of $y_{2}$ on $(a, b)$, since

$$
\left(\frac{y_{1}}{y_{2}}\right)^{\prime}=\frac{y_{1}^{\prime} y_{2}-y_{1} y_{2}^{\prime}}{y_{2}^{2}}=-\frac{W}{y_{2}^{2}}
$$

on any subinterval of $(a, b)$ where $y_{2}$ has no zeros.
We can now complete the proof of Theorem 5.1.3. From Theorem 5.1.5, two solutions $y_{1}$ and $y_{2}$ of (5.1.32) are linearly independent on $(a, b)$ if and only if $W$ has no zeros on $(a, b)$. From Theorem 5.1.4 and the motivating comments preceding it, $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.1.32) if and only if $W$ has no zeros on $(a, b)$. Therefore $\left\{y_{1}, y_{2}\right\}$ is a fundamental set for (5.1.32) on $(a, b)$ if and only if $\left\{y_{1}, y_{2}\right\}$ is linearly independent on $(a, b)$.

The next theorem summarizes the relationships among the concepts discussed in this section.
Theorem 5.1.6 Suppose $p$ and $q$ are continuous on an open interval $(a, b)$ and let $y_{1}$ and $y_{2}$ be solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.1.34}
\end{equation*}
$$

on $(a, b)$. Then the following statements are equivalent; that is, they are either all true or all false.
(a) The general solution of (5.1.34) on $(a, b)$ is $y=c_{1} y_{1}+c_{2} y_{2}$.
(b) $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.1.34) on $(a, b)$.
(c) $\left\{y_{1}, y_{2}\right\}$ is linearly independent on $(a, b)$.
(d) The Wronskian of $\left\{y_{1}, y_{2}\right\}$ is nonzero at some point in $(a, b)$.
(e) The Wronskian of $\left\{y_{1}, y_{2}\right\}$ is nonzero at all points in $(a, b)$.

We can apply this theorem to an equation written as

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0
$$

on an interval $(a, b)$ where $P_{0}, P_{1}$, and $P_{2}$ are continuous and $P_{0}$ has no zeros.
Theorem 5.1.7 Suppose $c$ is in $(a, b)$ and $\alpha$ and $\beta$ are real numbers, not both zero. Under the assumptions of Theorem 5.1.7, suppose $y_{1}$ and $y_{2}$ are solutions of (5.1.34) such that

$$
\begin{equation*}
\alpha y_{1}(c)+\beta y_{1}^{\prime}(c)=0 \text { and } \quad \alpha y_{2}(c)+\beta y_{2}^{\prime}(c)=0 \tag{5.1.35}
\end{equation*}
$$

Then $\left\{y_{1}, y_{2}\right\}$ isn't linearly independent on $(a, b)$.
Proof Since $\alpha$ and $\beta$ are not both zero, (5.1.35) implies that

$$
\left|\begin{array}{ll}
y_{1}(c) & y_{1}^{\prime}(c) \\
y_{2}(c) & y_{2}^{\prime}(c)
\end{array}\right|=0, \text { so } \quad\left|\begin{array}{ll}
y_{1}(c) & y_{2}(c) \\
y_{1}^{\prime}(c) & y_{2}^{\prime}(c)
\end{array}\right|=0
$$

and Theorem 5.1.6 implies the stated conclusion.

### 5.1 Exercises

1. (a) Verify that $y_{1}=e^{2 x}$ and $y_{2}=e^{5 x}$ are solutions of

$$
\begin{equation*}
y^{\prime \prime}-7 y^{\prime}+10 y=0 \tag{A}
\end{equation*}
$$

on $(-\infty, \infty)$.
(b) Verify that if $c_{1}$ and $c_{2}$ are arbitrary constants then $y=c_{1} e^{2 x}+c_{2} e^{5 x}$ is a solution of (A) on $(-\infty, \infty)$
(c) Solve the initial value problem

$$
y^{\prime \prime}-7 y^{\prime}+10 y=0, \quad y(0)=-1, \quad y^{\prime}(0)=1
$$

(d) Solve the initial value problem

$$
y^{\prime \prime}-7 y^{\prime}+10 y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

2. (a) Verify that $y_{1}=e^{x} \cos x$ and $y_{2}=e^{x} \sin x$ are solutions of

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+2 y=0 \tag{A}
\end{equation*}
$$

on $(-\infty, \infty)$.
(b) Verify that if $c_{1}$ and $c_{2}$ are arbitrary constants then $y=c_{1} e^{x} \cos x+c_{2} e^{x} \sin x$ is a solution of (A) on $(-\infty, \infty)$.
(c) Solve the initial value problem

$$
y^{\prime \prime}-2 y^{\prime}+2 y=0, \quad y(0)=3, \quad y^{\prime}(0)=-2
$$

(d) Solve the initial value problem

$$
y^{\prime \prime}-2 y^{\prime}+2 y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} .
$$

3. (a) Verify that $y_{1}=e^{x}$ and $y_{2}=x e^{x}$ are solutions of

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+y=0 \tag{A}
\end{equation*}
$$

on $(-\infty, \infty)$.
(b) Verify that if $c_{1}$ and $c_{2}$ are arbitrary constants then $y=e^{x}\left(c_{1}+c_{2} x\right)$ is a solution of (A) on $(-\infty, \infty)$
(c) Solve the initial value problem

$$
y^{\prime \prime}-2 y^{\prime}+y=0, \quad y(0)=7, \quad y^{\prime}(0)=4
$$

(d) Solve the initial value problem

$$
y^{\prime \prime}-2 y^{\prime}+y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} .
$$

4. (a) Verify that $y_{1}=1 /(x-1)$ and $y_{2}=1 /(x+1)$ are solutions of

$$
\begin{equation*}
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0 \tag{A}
\end{equation*}
$$

on $(-\infty,-1),(-1,1)$, and $(1, \infty)$. What is the general solution of $(\mathrm{A})$ on each of these intervals?
(b) Solve the initial value problem

$$
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0, \quad y(0)=-5, \quad y^{\prime}(0)=1 .
$$

What is the interval of validity of the solution?
(c) C/G Graph the solution of the initial value problem.
(d) Verify Abel's formula for $y_{1}$ and $y_{2}$, with $x_{0}=0$.
5. Compute the Wronskians of the given sets of functions.
(a) $\left\{1, e^{x}\right\}$
(b) $\left\{e^{x}, e^{x} \sin x\right\}$
(c) $\left\{x+1, x^{2}+2\right\}$
(d) $\left\{x^{1 / 2}, x^{-1 / 3}\right\}$
(e) $\left\{\frac{\sin x}{x}, \frac{\cos x}{x}\right\}$
(f) $\left\{x \ln |x|, x^{2} \ln |x|\right\}$
(g) $\left\{e^{x} \cos \sqrt{x}, e^{x} \sin \sqrt{x}\right\}$
6. Find the Wronskian of a given set $\left\{y_{1}, y_{2}\right\}$ of solutions of

$$
y^{\prime \prime}+3\left(x^{2}+1\right) y^{\prime}-2 y=0
$$

given that $W(\pi)=0$.
7. Find the Wronskian of a given set $\left\{y_{1}, y_{2}\right\}$ of solutions of

$$
\left(1-x^{2}\right) y^{\prime \prime}-2 x y^{\prime}+\alpha(\alpha+1) y=0
$$

given that $W(0)=1$. (This is Legendre's equation.)
8. Find the Wronskian of a given set $\left\{y_{1}, y_{2}\right\}$ of solutions of

$$
x^{2} y^{\prime \prime}+x y^{\prime}+\left(x^{2}-\nu^{2}\right) y=0
$$

given that $W(1)=1$. (This is Bessel's equation.)
9. (This exercise shows that if you know one nontrivial solution of $y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0$, you can use Abel's formula to find another.)
Suppose $p$ and $q$ are continuous and $y_{1}$ is a solution of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

that has no zeros on $(a, b)$. Let $P(x)=\int p(x) d x$ be any antiderivative of $p$ on $(a, b)$.
(a) Show that if $K$ is an arbitrary nonzero constant and $y_{2}$ satisfies

$$
\begin{equation*}
y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}=K e^{-P(x)} \tag{B}
\end{equation*}
$$

on $(a, b)$, then $y_{2}$ also satisfies (A) on $(a, b)$, and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions on (A) on $(a, b)$.
(b) Conclude from (a) that if $y_{2}=u y_{1}$ where $u^{\prime}=K \frac{e^{-P(x)}}{y_{1}^{2}(x)}$, then $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (A) on $(a, b)$.

In Exercises 10-23 use the method suggested by Exercise 9 to find a second solution $y_{2}$ that isn't a constant multiple of the solution $y_{1}$. Choose $K$ conveniently to simplify $y_{2}$.
10. $y^{\prime \prime}-2 y^{\prime}-3 y=0 ; \quad y_{1}=e^{3 x}$
11. $y^{\prime \prime}-6 y^{\prime}+9 y=0 ; \quad y_{1}=e^{3 x}$
12. $y^{\prime \prime}-2 a y^{\prime}+a^{2} y=0(a=$ constant $) ; \quad y_{1}=e^{a x}$
13. $x^{2} y^{\prime \prime}+x y^{\prime}-y=0 ; \quad y_{1}=x$
14. $x^{2} y^{\prime \prime}-x y^{\prime}+y=0 ; \quad y_{1}=x$
15. $x^{2} y^{\prime \prime}-(2 a-1) x y^{\prime}+a^{2} y=0(a=$ nonzero constant $) ; x>0 ; \quad y_{1}=x^{a}$
16. $4 x^{2} y^{\prime \prime}-4 x y^{\prime}+\left(3-16 x^{2}\right) y=0 ; \quad y_{1}=x^{1 / 2} e^{2 x}$
17. $(x-1) y^{\prime \prime}-x y^{\prime}+y=0 ; \quad y_{1}=e^{x}$
18. $x^{2} y^{\prime \prime}-2 x y^{\prime}+\left(x^{2}+2\right) y=0 ; \quad y_{1}=x \cos x$
19. $4 x^{2}(\sin x) y^{\prime \prime}-4 x(x \cos x+\sin x) y^{\prime}+(2 x \cos x+3 \sin x) y=0 ; \quad y_{1}=x^{1 / 2}$
20. $(3 x-1) y^{\prime \prime}-(3 x+2) y^{\prime}-(6 x-8) y=0 ; \quad y_{1}=e^{2 x}$
21. $\left(x^{2}-4\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0 ; \quad y_{1}=\frac{1}{x-2}$
22. $(2 x+1) x y^{\prime \prime}-2\left(2 x^{2}-1\right) y^{\prime}-4(x+1) y=0 ; \quad y_{1}=\frac{1}{x}$
23. $\left(x^{2}-2 x\right) y^{\prime \prime}+\left(2-x^{2}\right) y^{\prime}+(2 x-2) y=0 ; \quad y_{1}=e^{x}$
24. Suppose $p$ and $q$ are continuous on an open interval $(a, b)$ and let $x_{0}$ be in $(a, b)$. Use Theorem 5.1.1 to show that the only solution of the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0, \quad y\left(x_{0}\right)=0, \quad y^{\prime}\left(x_{0}\right)=0
$$

on $(a, b)$ is the trivial solution $y \equiv 0$.
25. Suppose $P_{0}, P_{1}$, and $P_{2}$ are continuous on $(a, b)$ and let $x_{0}$ be in $(a, b)$. Show that if either of the following statements is true then $P_{0}(x)=0$ for some $x$ in $(a, b)$.
(a) The initial value problem

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

has more than one solution on $(a, b)$.
(b) The initial value problem

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0, \quad y\left(x_{0}\right)=0, \quad y^{\prime}\left(x_{0}\right)=0
$$

has a nontrivial solution on $(a, b)$.
26. Suppose $p$ and $q$ are continuous on $(a, b)$ and $y_{1}$ and $y_{2}$ are solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

on $(a, b)$. Let

$$
z_{1}=\alpha y_{1}+\beta y_{2} \quad \text { and } \quad z_{2}=\gamma y_{1}+\delta y_{2}
$$

where $\alpha, \beta, \gamma$, and $\delta$ are constants. Show that if $\left\{z_{1}, z_{2}\right\}$ is a fundamental set of solutions of (A) on $(a, b)$ then so is $\left\{y_{1}, y_{2}\right\}$.
27. Suppose $p$ and $q$ are continuous on $(a, b)$ and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

on $(a, b)$. Let

$$
z_{1}=\alpha y_{1}+\beta y_{2} \quad \text { and } \quad z_{2}=\gamma y_{1}+\delta y_{2}
$$

where $\alpha, \beta, \gamma$, and $\delta$ are constants. Show that $\left\{z_{1}, z_{2}\right\}$ is a fundamental set of solutions of (A) on $(a, b)$ if and only if $\alpha \gamma-\beta \delta \neq 0$.
28. Suppose $y_{1}$ is differentiable on an interval $(a, b)$ and $y_{2}=k y_{1}$, where $k$ is a constant. Show that the Wronskian of $\left\{y_{1}, y_{2}\right\}$ is identically zero on $(a, b)$.
29. Let

$$
y_{1}=x^{3} \quad \text { and } \quad y_{2}=\left\{\begin{aligned}
x^{3}, & x \geq 0, \\
-x^{3}, & x<0 .
\end{aligned}\right.
$$

(a) Show that the Wronskian of $\left\{y_{1}, y_{2}\right\}$ is defined and identically zero on $(-\infty, \infty)$.
(b) Suppose $a<0<b$. Show that $\left\{y_{1}, y_{2}\right\}$ is linearly independent on $(a, b)$.
(c) Use Exercise 25(b) to show that these results don't contradict Theorem 5.1.5, because neither $y_{1}$ nor $y_{2}$ can be a solution of an equation

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0
$$

on $(a, b)$ if $p$ and $q$ are continuous on $(a, b)$.
30. Suppose $p$ and $q$ are continuous on $(a, b)$ and $\left\{y_{1}, y_{2}\right\}$ is a set of solutions of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0
$$

on $(a, b)$ such that either $y_{1}\left(x_{0}\right)=y_{2}\left(x_{0}\right)=0$ or $y_{1}^{\prime}\left(x_{0}\right)=y_{2}^{\prime}\left(x_{0}\right)=0$ for some $x_{0}$ in $(a, b)$. Show that $\left\{y_{1}, y_{2}\right\}$ is linearly dependent on $(a, b)$.
31. Suppose $p$ and $q$ are continuous on $(a, b)$ and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0
$$

on $(a, b)$. Show that if $y_{1}\left(x_{1}\right)=y_{1}\left(x_{2}\right)=0$, where $a<x_{1}<x_{2}<b$, then $y_{2}(x)=0$ for some $x$ in $\left(x_{1}, x_{2}\right)$. HinT: Show that if $y_{2}$ has no zeros in $\left(x_{1}, x_{2}\right)$, then $y_{1} / y_{2}$ is either strictly increasing or strictly decreasing on $\left(x_{1}, x_{2}\right)$, and deduce a contradiction.
32. Suppose $p$ and $q$ are continuous on $(a, b)$ and every solution of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

on $(a, b)$ can be written as a linear combination of the twice differentiable functions $\left\{y_{1}, y_{2}\right\}$. Use Theorem 5.1.1 to show that $y_{1}$ and $y_{2}$ are themselves solutions of (A) on $(a, b)$.
33. Suppose $p_{1}, p_{2}, q_{1}$, and $q_{2}$ are continuous on $(a, b)$ and the equations

$$
y^{\prime \prime}+p_{1}(x) y^{\prime}+q_{1}(x) y=0 \quad \text { and } \quad y^{\prime \prime}+p_{2}(x) y^{\prime}+q_{2}(x) y=0
$$

have the same solutions on $(a, b)$. Show that $p_{1}=p_{2}$ and $q_{1}=q_{2}$ on $(a, b)$. Hint: Use Abel's formula.
34. (For this exercise you have to know about $3 \times 3$ determinants.) Show that if $y_{1}$ and $y_{2}$ are twice continuously differentiable on $(a, b)$ and the Wronskian $W$ of $\left\{y_{1}, y_{2}\right\}$ has no zeros in $(a, b)$ then the equation

$$
\frac{1}{W}\left|\begin{array}{ccc}
y & y_{1} & y_{2} \\
y^{\prime} & y_{1}^{\prime} & y_{2}^{\prime} \\
y^{\prime \prime} & y_{1}^{\prime \prime} & y_{2}^{\prime \prime}
\end{array}\right|=0
$$

can be written as

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

where $p$ and $q$ are continuous on $(a, b)$ and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (A) on $(a, b)$. Hint: Expand the determinant by cofactors of its first column.
35. Use the method suggested by Exercise 34 to find a linear homogeneous equation for which the given functions form a fundamental set of solutions on some interval.
(a) $e^{x} \cos 2 x, \quad e^{x} \sin 2 x$
(b) $x, \quad e^{2 x}$
(c) $x, \quad x \ln x$
(d) $\cos (\ln x), \quad \sin (\ln x)$
(e) $\cosh x, \quad \sinh x$
(f) $x^{2}-1, \quad x^{2}+1$
36. Suppose $p$ and $q$ are continuous on $(a, b)$ and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

on $(a, b)$. Show that if $y$ is a solution of (A) on $(a, b)$, there's exactly one way to choose $c_{1}$ and $c_{2}$ so that $y=c_{1} y_{1}+c_{2} y_{2}$ on $(a, b)$.
37. Suppose $p$ and $q$ are continuous on $(a, b)$ and $x_{0}$ is in $(a, b)$. Let $y_{1}$ and $y_{2}$ be the solutions of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{A}
\end{equation*}
$$

such that

$$
y_{1}\left(x_{0}\right)=1, \quad y_{1}^{\prime}\left(x_{0}\right)=0 \quad \text { and } \quad y_{2}\left(x_{0}\right)=0, y_{2}^{\prime}\left(x_{0}\right)=1 .
$$

(Theorem 5.1.1 implies that each of these initial value problems has a unique solution on $(a, b)$.)
(a) Show that $\left\{y_{1}, y_{2}\right\}$ is linearly independent on $(a, b)$.
(b) Show that an arbitrary solution $y$ of (A) on $(a, b)$ can be written as $y=y\left(x_{0}\right) y_{1}+y^{\prime}\left(x_{0}\right) y_{2}$.
(c) Express the solution of the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

as a linear combination of $y_{1}$ and $y_{2}$.
38. Find solutions $y_{1}$ and $y_{2}$ of the equation $y^{\prime \prime}=0$ that satisfy the initial conditions

$$
y_{1}\left(x_{0}\right)=1, \quad y_{1}^{\prime}\left(x_{0}\right)=0 \quad \text { and } \quad y_{2}\left(x_{0}\right)=0, \quad y_{2}^{\prime}\left(x_{0}\right)=1 .
$$

Then use Exercise 37 (c) to write the solution of the initial value problem

$$
y^{\prime \prime}=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

as a linear combination of $y_{1}$ and $y_{2}$.
39. Let $x_{0}$ be an arbitrary real number. Given (Example 5.1.1) that $e^{x}$ and $e^{-x}$ are solutions of $y^{\prime \prime}-y=$ 0 , find solutions $y_{1}$ and $y_{2}$ of $y^{\prime \prime}-y=0$ such that

$$
y_{1}\left(x_{0}\right)=1, \quad y_{1}^{\prime}\left(x_{0}\right)=0 \quad \text { and } \quad y_{2}\left(x_{0}\right)=0, y_{2}^{\prime}\left(x_{0}\right)=1 .
$$

Then use Exercise 37 (c) to write the solution of the initial value problem

$$
y^{\prime \prime}-y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

as a linear combination of $y_{1}$ and $y_{2}$.
40. Let $x_{0}$ be an arbitrary real number. Given (Example 5.1.2) that $\cos \omega x$ and $\sin \omega x$ are solutions of $y^{\prime \prime}+\omega^{2} y=0$, find solutions of $y^{\prime \prime}+\omega^{2} y=0$ such that

$$
y_{1}\left(x_{0}\right)=1, \quad y_{1}^{\prime}\left(x_{0}\right)=0 \quad \text { and } \quad y_{2}\left(x_{0}\right)=0, y_{2}^{\prime}\left(x_{0}\right)=1 .
$$

Then use Exercise 37 (c) to write the solution of the initial value problem

$$
y^{\prime \prime}+\omega^{2} y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

as a linear combination of $y_{1}$ and $y_{2}$. Use the identities

$$
\begin{aligned}
\cos (A+B) & =\cos A \cos B-\sin A \sin B \\
\sin (A+B) & =\sin A \cos B+\cos A \sin B
\end{aligned}
$$

to simplify your expressions for $y_{1}, y_{2}$, and $y$.
41. Recall from Exercise 4 that $1 /(x-1)$ and $1 /(x+1)$ are solutions of

$$
\begin{equation*}
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0 \tag{A}
\end{equation*}
$$

on $(-1,1)$. Find solutions of (A) such that

$$
y_{1}(0)=1, \quad y_{1}^{\prime}(0)=0 \quad \text { and } \quad y_{2}(0)=0, \quad y_{2}^{\prime}(0)=1 .
$$

Then use Exercise 37 (c) to write the solution of initial value problem

$$
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

as a linear combination of $y_{1}$ and $y_{2}$.
42. (a) Verify that $y_{1}=x^{2}$ and $y_{2}=x^{3}$ satisfy

$$
\begin{equation*}
x^{2} y^{\prime \prime}-4 x y^{\prime}+6 y=0 \tag{A}
\end{equation*}
$$

on $(-\infty, \infty)$ and that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of $(\mathrm{A})$ on $(-\infty, 0)$ and $(0, \infty)$.
(b) Let $a_{1}, a_{2}, b_{1}$, and $b_{2}$ be constants. Show that

$$
y=\left\{\begin{array}{cc}
a_{1} x^{2}+a_{2} x^{3}, & x \geq 0 \\
b_{1} x^{2}+b_{2} x^{3}, & x<0
\end{array}\right.
$$

is a solution of $(\mathrm{A})$ on $(-\infty, \infty)$ if and only if $a_{1}=b_{1}$. From this, justify the statement that $y$ is a solution of $(\mathrm{A})$ on $(-\infty, \infty)$ if and only if

$$
y= \begin{cases}c_{1} x^{2}+c_{2} x^{3}, & x \geq 0, \\ c_{1} x^{2}+c_{3} x^{3}, & x<0,\end{cases}
$$

where $c_{1}, c_{2}$, and $c_{3}$ are arbitrary constants.
(c) For what values of $k_{0}$ and $k_{1}$ does the initial value problem

$$
x^{2} y^{\prime \prime}-4 x y^{\prime}+6 y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

have a solution? What are the solutions?
(d) Show that if $x_{0} \neq 0$ and $k_{0}, k_{1}$ are arbitrary constants, the initial value problem

$$
\begin{equation*}
x^{2} y^{\prime \prime}-4 x y^{\prime}+6 y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1} \tag{B}
\end{equation*}
$$

has infinitely many solutions on $(-\infty, \infty)$. On what interval does (B) have a unique solution?
43. (a) Verify that $y_{1}=x$ and $y_{2}=x^{2}$ satisfy

$$
\begin{equation*}
x^{2} y^{\prime \prime}-2 x y^{\prime}+2 y=0 \tag{A}
\end{equation*}
$$

on $(-\infty, \infty)$ and that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (A) on $(-\infty, 0)$ and $(0, \infty)$.
(b) Let $a_{1}, a_{2}, b_{1}$, and $b_{2}$ be constants. Show that

$$
y=\left\{\begin{array}{cc}
a_{1} x+a_{2} x^{2}, & x \geq 0 \\
b_{1} x+b_{2} x^{2}, & x<0
\end{array}\right.
$$

is a solution of $(\mathrm{A})$ on $(-\infty, \infty)$ if and only if $a_{1}=b_{1}$ and $a_{2}=b_{2}$. From this, justify the statement that the general solution of $(\mathrm{A})$ on $(-\infty, \infty)$ is $y=c_{1} x+c_{2} x^{2}$, where $c_{1}$ and $c_{2}$ are arbitrary constants.
(c) For what values of $k_{0}$ and $k_{1}$ does the initial value problem

$$
x^{2} y^{\prime \prime}-2 x y^{\prime}+2 y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

have a solution? What are the solutions?
(d) Show that if $x_{0} \neq 0$ and $k_{0}, k_{1}$ are arbitrary constants then the initial value problem

$$
x^{2} y^{\prime \prime}-2 x y^{\prime}+2 y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

has a unique solution on $(-\infty, \infty)$.
44. (a) Verify that $y_{1}=x^{3}$ and $y_{2}=x^{4}$ satisfy

$$
\begin{equation*}
x^{2} y^{\prime \prime}-6 x y^{\prime}+12 y=0 \tag{A}
\end{equation*}
$$

on $(-\infty, \infty)$, and that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of $(\mathrm{A})$ on $(-\infty, 0)$ and $(0, \infty)$.
(b) Show that $y$ is a solution of (A) on $(-\infty, \infty)$ if and only if

$$
y=\left\{\begin{array}{cc}
a_{1} x^{3}+a_{2} x^{4}, & x \geq 0 \\
b_{1} x^{3}+b_{2} x^{4}, & x<0
\end{array}\right.
$$

where $a_{1}, a_{2}, b_{1}$, and $b_{2}$ are arbitrary constants.
(c) For what values of $k_{0}$ and $k_{1}$ does the initial value problem

$$
x^{2} y^{\prime \prime}-6 x y^{\prime}+12 y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

have a solution? What are the solutions?
(d) Show that if $x_{0} \neq 0$ and $k_{0}, k_{1}$ are arbitrary constants then the initial value problem

$$
\begin{equation*}
x^{2} y^{\prime \prime}-6 x y^{\prime}+12 y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1} \tag{B}
\end{equation*}
$$

has infinitely many solutions on $(-\infty, \infty)$. On what interval does (B) have a unique solution?

### 5.2 CONSTANT COEFFICIENT HOMOGENEOUS EQUATIONS

If $a, b$, and $c$ are real constants and $a \neq 0$, then

$$
a y^{\prime \prime}+b y^{\prime}+c y=F(x)
$$

is said to be a constant coefficient equation. In this section we consider the homogeneous constant coefficient equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{5.2.1}
\end{equation*}
$$

As we'll see, all solutions of (5.2.1) are defined on $(-\infty, \infty)$. This being the case, we'll omit references to the interval on which solutions are defined, or on which a given set of solutions is a fundamental set, etc., since the interval will always be $(-\infty, \infty)$.
The key to solving (5.2.1) is that if $y=e^{r x}$ where $r$ is a constant then the left side of (5.2.1) is a multiple of $e^{r x}$; thus, if $y=e^{r x}$ then $y^{\prime}=r e^{r x}$ and $y^{\prime \prime}=r^{2} e^{r x}$, so

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=a r^{2} e^{r x}+b r e^{r x}+c e^{r x}=\left(a r^{2}+b r+c\right) e^{r x} \tag{5.2.2}
\end{equation*}
$$

The quadratic polynomial

$$
p(r)=a r^{2}+b r+c
$$

is the characteristic polynomial of (5.2.1), and $p(r)=0$ is the characteristic equation. From (5.2.2) we can see that $y=e^{r x}$ is a solution of (5.2.1) if and only if $p(r)=0$.
The roots of the characteristic equation are given by the quadratic formula

$$
\begin{equation*}
r=\frac{-b \pm \sqrt{b^{2}-4 a c}}{2 a} \tag{5.2.3}
\end{equation*}
$$

We consider three cases:
CASE 1. $b^{2}-4 a c>0$, so the characteristic equation has two distinct real roots.
CASE 2. $b^{2}-4 a c=0$, so the characteristic equation has a repeated real root.
CASE 3. $b^{2}-4 a c<0$, so the characteristic equation has complex roots.
In each case we'll start with an example.
Case 1: Distinct Real Roots

## Example 5.2.1

(a) Find the general solution of

$$
\begin{equation*}
y^{\prime \prime}+6 y^{\prime}+5 y=0 . \tag{5.2.4}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+6 y^{\prime}+5 y=0, \quad y(0)=3, y^{\prime}(0)=-1 . \tag{5.2.5}
\end{equation*}
$$

$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ The characteristic polynomial of (5.2.4) is

$$
p(r)=r^{2}+6 r+5=(r+1)(r+5) .
$$

Since $p(-1)=p(-5)=0, y_{1}=e^{-x}$ and $y_{2}=e^{-5 x}$ are solutions of (5.2.4). Since $y_{2} / y_{1}=e^{-4 x}$ is nonconstant, 5.1.6 implies that the general solution of (5.2.4) is

$$
\begin{equation*}
y=c_{1} e^{-x}+c_{2} e^{-5 x} \tag{5.2.6}
\end{equation*}
$$

SOLUTION(b) We must determine $c_{1}$ and $c_{2}$ in (5.2.6) so that $y$ satisfies the initial conditions in (5.2.5). Differentiating (5.2.6) yields

$$
\begin{equation*}
y^{\prime}=-c_{1} e^{-x}-5 c_{2} e^{-5 x} \tag{5.2.7}
\end{equation*}
$$



Figure 5.2.1 $y=\frac{7}{2} e^{-x}-\frac{1}{2} e^{-5 x}$

Imposing the initial conditions $y(0)=3, y^{\prime}(0)=-1$ in (5.2.6) and (5.2.7) yields

$$
\begin{aligned}
c_{1}+c_{2} & = & 3 \\
-c_{1}-5 c_{2} & = & -1 .
\end{aligned}
$$

The solution of this system is $c_{1}=7 / 2, c_{2}=-1 / 2$. Therefore the solution of (5.2.5) is

$$
y=\frac{7}{2} e^{-x}-\frac{1}{2} e^{-5 x}
$$

Figure 5.2.1 is a graph of this solution.
If the characteristic equation has arbitrary distinct real roots $r_{1}$ and $r_{2}$, then $y_{1}=e^{r_{1} x}$ and $y_{2}=e^{r_{2} x}$ are solutions of $a y^{\prime \prime}+b y^{\prime}+c y=0$. Since $y_{2} / y_{1}=e^{\left(r_{2}-r_{1}\right) x}$ is nonconstant, Theorem 5.1.6 implies that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of $a y^{\prime \prime}+b y^{\prime}+c y=0$.
Case 2: A Repeated Real Root

## Example 5.2.2

(a) Find the general solution of

$$
\begin{equation*}
y^{\prime \prime}+6 y^{\prime}+9 y=0 . \tag{5.2.8}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+6 y^{\prime}+9 y=0, \quad y(0)=3, y^{\prime}(0)=-1 \tag{5.2.9}
\end{equation*}
$$

$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ The characteristic polynomial of (5.2.8) is

$$
p(r)=r^{2}+6 r+9=(r+3)^{2},
$$

so the characteristic equation has the repeated real root $r_{1}=-3$. Therefore $y_{1}=e^{-3 x}$ is a solution of (5.2.8). Since the characteristic equation has no other roots, (5.2.8) has no other solutions of the form $e^{r x}$. We look for solutions of the form $y=u y_{1}=u e^{-3 x}$, where $u$ is a function that we'll now determine. (This should remind you of the method of variation of parameters used in Section 2.1 to
solve the nonhomogeneous equation $y^{\prime}+p(x) y=f(x)$, given a solution $y_{1}$ of the complementary equation $y^{\prime}+p(x) y=0$. It's also a special case of a method called reduction of order that we'll study in Section 5.6. For other ways to obtain a second solution of (5.2.8) that's not a multiple of $e^{-3 x}$, see Exercises 5.1.9, 5.1.12, and 33.
If $y=u e^{-3 x}$, then

$$
y^{\prime}=u^{\prime} e^{-3 x}-3 u e^{-3 x} \quad \text { and } \quad y^{\prime \prime}=u^{\prime \prime} e^{-3 x}-6 u^{\prime} e^{-3 x}+9 u e^{-3 x}
$$

so

$$
\begin{aligned}
y^{\prime \prime}+6 y^{\prime}+9 y & =e^{-3 x}\left[\left(u^{\prime \prime}-6 u^{\prime}+9 u\right)+6\left(u^{\prime}-3 u\right)+9 u\right] \\
& =e^{-3 x}\left[u^{\prime \prime}-(6-6) u^{\prime}+(9-18+9) u\right]=u^{\prime \prime} e^{-3 x}
\end{aligned}
$$

Therefore $y=u e^{-3 x}$ is a solution of (5.2.8) if and only if $u^{\prime \prime}=0$, which is equivalent to $u=c_{1}+c_{2} x$, where $c_{1}$ and $c_{2}$ are constants. Therefore any function of the form

$$
\begin{equation*}
y=e^{-3 x}\left(c_{1}+c_{2} x\right) \tag{5.2.10}
\end{equation*}
$$

is a solution of (5.2.8). Letting $c_{1}=1$ and $c_{2}=0$ yields the solution $y_{1}=e^{-3 x}$ that we already knew. Letting $c_{1}=0$ and $c_{2}=1$ yields the second solution $y_{2}=x e^{-3 x}$. Since $y_{2} / y_{1}=x$ is nonconstant, 5.1.6 implies that $\left\{y_{1}, y_{2}\right\}$ is fundamental set of solutions of (5.2.8), and (5.2.10) is the general solution.
$\underline{\text { SOLUTION(b) }}$ Differentiating (5.2.10) yields

$$
\begin{equation*}
y^{\prime}=-3 e^{-3 x}\left(c_{1}+c_{2} x\right)+c_{2} e^{-3 x} \tag{5.2.11}
\end{equation*}
$$

Imposing the initial conditions $y(0)=3, y^{\prime}(0)=-1$ in (5.2.10) and (5.2.11) yields $c_{1}=3$ and $-3 c_{1}+$ $c_{2}=-1$, so $c_{2}=8$. Therefore the solution of (5.2.9) is

$$
y=e^{-3 x}(3+8 x)
$$

Figure 5.2.2 is a graph of this solution.


Figure 5.2.2 $y=e^{-3 x}(3+8 x)$

If the characteristic equation of $a y^{\prime \prime}+b y^{\prime}+c y=0$ has an arbitrary repeated root $r_{1}$, the characteristic polynomial must be

$$
p(r)=a\left(r-r_{1}\right)^{2}=a\left(r^{2}-2 r_{1} r+r_{1}^{2}\right) .
$$

Therefore

$$
a r^{2}+b r+c=a r^{2}-\left(2 a r_{1}\right) r+a r_{1}^{2}
$$

which implies that $b=-2 a r_{1}$ and $c=a r_{1}^{2}$. Therefore $a y^{\prime \prime}+b y^{\prime}+c y=0$ can be written as $a\left(y^{\prime \prime}-\right.$ $\left.2 r_{1} y^{\prime}+r_{1}^{2} y\right)=0$. Since $a \neq 0$ this equation has the same solutions as

$$
\begin{equation*}
y^{\prime \prime}-2 r_{1} y^{\prime}+r_{1}^{2} y=0 \tag{5.2.12}
\end{equation*}
$$

Since $p\left(r_{1}\right)=0, \mathrm{t} y_{1}=e^{r_{1} x}$ is a solution of $a y^{\prime \prime}+b y^{\prime}+c y=0$, and therefore of (5.2.12). Proceeding as in Example 5.2.2, we look for other solutions of (5.2.12) of the form $y=u e^{r_{1} x}$; then

$$
y^{\prime}=u^{\prime} e^{r_{1} x}+r u e^{r_{1} x} \quad \text { and } \quad y^{\prime \prime}=u^{\prime \prime} e^{r_{1} x}+2 r_{1} u^{\prime} e^{r_{1} x}+r_{1}^{2} u e^{r_{1} x}
$$

so

$$
\begin{aligned}
y^{\prime \prime}-2 r_{1} y^{\prime}+r_{1}^{2} y & =e^{r x}\left[\left(u^{\prime \prime}+2 r_{1} u^{\prime}+r_{1}^{2} u\right)-2 r_{1}\left(u^{\prime}+r_{1} u\right)+r_{1}^{2} u\right] \\
& =e^{r_{1} x}\left[u^{\prime \prime}+\left(2 r_{1}-2 r_{1}\right) u^{\prime}+\left(r_{1}^{2}-2 r_{1}^{2}+r_{1}^{2}\right) u\right]=u^{\prime \prime} e^{r_{1} x}
\end{aligned}
$$

Therefore $y=u e^{r_{1} x}$ is a solution of (5.2.12) if and only if $u^{\prime \prime}=0$, which is equivalent to $u=c_{1}+c_{2} x$, where $c_{1}$ and $c_{2}$ are constants. Hence, any function of the form

$$
\begin{equation*}
y=e^{r_{1} x}\left(c_{1}+c_{2} x\right) \tag{5.2.13}
\end{equation*}
$$

is a solution of (5.2.12). Letting $c_{1}=1$ and $c_{2}=0$ here yields the solution $y_{1}=e^{r_{1} x}$ that we already knew. Letting $c_{1}=0$ and $c_{2}=1$ yields the second solution $y_{2}=x e^{r_{1} x}$. Since $y_{2} / y_{1}=x$ is nonconstant, 5.1.6 implies that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.2.12), and (5.2.13) is the general solution.
Case 3: Complex Conjugate Roots

## Example 5.2.3

(a) Find the general solution of

$$
\begin{equation*}
y^{\prime \prime}+4 y^{\prime}+13 y=0 \tag{5.2.14}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+4 y^{\prime}+13 y=0, \quad y(0)=2, \quad y^{\prime}(0)=-3 . \tag{5.2.15}
\end{equation*}
$$

$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ The characteristic polynomial of (5.2.14) is

$$
p(r)=r^{2}+4 r+13=r^{2}+4 r+4+9=(r+2)^{2}+9
$$

The roots of the characteristic equation are $r_{1}=-2+3 i$ and $r_{2}=-2-3 i$. By analogy with Case 1 , it's reasonable to expect that $e^{(-2+3 i) x}$ and $e^{(-2-3 i) x}$ are solutions of (5.2.14). This is true (see Exercise 34); however, there are difficulties here, since you are probably not familiar with exponential functions with complex arguments, and even if you are, it's inconvenient to work with them, since they are complexvalued. We'll take a simpler approach, which we motivate as follows: the exponential notation suggests that

$$
e^{(-2+3 i) x}=e^{-2 x} e^{3 i x} \quad \text { and } \quad e^{(-2-3 i) x}=e^{-2 x} e^{-3 i x},
$$

so even though we haven't defined $e^{3 i x}$ and $e^{-3 i x}$, it's reasonable to expect that every linear combination of $e^{(-2+3 i) x}$ and $e^{(-2-3 i) x}$ can be written as $y=u e^{-2 x}$, where $u$ depends upon $x$. To determine $u$, we note that if $y=u e^{-2 x}$ then

$$
y^{\prime}=u^{\prime} e^{-2 x}-2 u e^{-2 x} \quad \text { and } \quad y^{\prime \prime}=u^{\prime \prime} e^{-2 x}-4 u^{\prime} e^{-2 x}+4 u e^{-2 x}
$$

so

$$
\begin{aligned}
y^{\prime \prime}+4 y^{\prime}+13 y & =e^{-2 x}\left[\left(u^{\prime \prime}-4 u^{\prime}+4 u\right)+4\left(u^{\prime}-2 u\right)+13 u\right] \\
& =e^{-2 x}\left[u^{\prime \prime}-(4-4) u^{\prime}+(4-8+13) u\right]=e^{-2 x}\left(u^{\prime \prime}+9 u\right)
\end{aligned}
$$



Figure 5.2.3 $y=e^{-2 x}\left(2 \cos 3 x+\frac{1}{3} \sin 3 x\right)$

Therefore $y=u e^{-2 x}$ is a solution of (5.2.14) if and only if

$$
u^{\prime \prime}+9 u=0 .
$$

From Example 5.1.2, the general solution of this equation is

$$
u=c_{1} \cos 3 x+c_{2} \sin 3 x .
$$

Therefore any function of the form

$$
\begin{equation*}
y=e^{-2 x}\left(c_{1} \cos 3 x+c_{2} \sin 3 x\right) \tag{5.2.16}
\end{equation*}
$$

is a solution of (5.2.14). Letting $c_{1}=1$ and $c_{2}=0$ yields the solution $y_{1}=e^{-2 x} \cos 3 x$. Letting $c_{1}=0$ and $c_{2}=1$ yields the second solution $y_{2}=e^{-2 x} \sin 3 x$. Since $y_{2} / y_{1}=\tan 3 x$ is nonconstant, 5.1.6 implies that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.2.14), and (5.2.16) is the general solution.
$\underline{\text { SOLUTION(b) }}$ Imposing the condition $y(0)=2$ in (5.2.16) shows that $c_{1}=2$. Differentiating (5.2.16) yields

$$
y^{\prime}=-2 e^{-2 x}\left(c_{1} \cos 3 x+c_{2} \sin 3 x\right)+3 e^{-2 x}\left(-c_{1} \sin 3 x+c_{2} \cos 3 x\right),
$$

and imposing the initial condition $y^{\prime}(0)=-3$ here yields $-3=-2 c_{1}+3 c_{2}=-4+3 c_{2}$, so $c_{2}=1 / 3$. Therefore the solution of (5.2.15) is

$$
y=e^{-2 x}\left(2 \cos 3 x+\frac{1}{3} \sin 3 x\right) .
$$

Figure 5.2.3 is a graph of this function.
Now suppose the characteristic equation of $a y^{\prime \prime}+b y^{\prime}+c y=0$ has arbitrary complex roots; thus, $b^{2}-4 a c<0$ and, from (5.2.3), the roots are

$$
r_{1}=\frac{-b+i \sqrt{4 a c-b^{2}}}{2 a}, \quad r_{2}=\frac{-b-i \sqrt{4 a c-b^{2}}}{2 a}
$$

which we rewrite as

$$
\begin{equation*}
r_{1}=\lambda+i \omega, \quad r_{2}=\lambda-i \omega, \tag{5.2.17}
\end{equation*}
$$

with

$$
\lambda=-\frac{b}{2 a}, \quad \omega=\frac{\sqrt{4 a c-b^{2}}}{2 a} .
$$

Don't memorize these formulas. Just remember that $r_{1}$ and $r_{2}$ are of the form (5.2.17), where $\lambda$ is an arbitrary real number and $\omega$ is positive; $\lambda$ and $\omega$ are the real and imaginary parts, respectively, of $r_{1}$. Similarly, $\lambda$ and $-\omega$ are the real and imaginary parts of $r_{2}$. We say that $r_{1}$ and $r_{2}$ are complex conjugates, which means that they have the same real part and their imaginary parts have the same absolute values, but opposite signs.

As in Example 5.2.3, it's reasonable to to expect that the solutions of $a y^{\prime \prime}+b y^{\prime}+c y=0$ are linear combinations of $e^{(\lambda+i \omega) x}$ and $e^{(\lambda-i \omega) x}$. Again, the exponential notation suggests that

$$
e^{(\lambda+i \omega) x}=e^{\lambda x} e^{i \omega x} \quad \text { and } \quad e^{(\lambda-i \omega) x}=e^{\lambda x} e^{-i \omega x}
$$

so even though we haven't defined $e^{i \omega x}$ and $e^{-i \omega x}$, it's reasonable to expect that every linear combination of $e^{(\lambda+i \omega) x}$ and $e^{(\lambda-i \omega) x}$ can be written as $y=u e^{\lambda x}$, where $u$ depends upon $x$. To determine $u$ we first observe that since $r_{1}=\lambda+i \omega$ and $r_{2}=\lambda-i \omega$ are the roots of the characteristic equation, $p$ must be of the form

$$
\begin{aligned}
p(r) & =a\left(r-r_{1}\right)\left(r-r_{2}\right) \\
& =a(r-\lambda-i \omega)(r-\lambda+i \omega) \\
& =a\left[(r-\lambda)^{2}+\omega^{2}\right] \\
& =a\left(r^{2}-2 \lambda r+\lambda^{2}+\omega^{2}\right) .
\end{aligned}
$$

Therefore $a y^{\prime \prime}+b y^{\prime}+c y=0$ can be written as

$$
a\left[y^{\prime \prime}-2 \lambda y^{\prime}+\left(\lambda^{2}+\omega^{2}\right) y\right]=0 .
$$

Since $a \neq 0$ this equation has the same solutions as

$$
\begin{equation*}
y^{\prime \prime}-2 \lambda y^{\prime}+\left(\lambda^{2}+\omega^{2}\right) y=0 \tag{5.2.18}
\end{equation*}
$$

To determine $u$ we note that if $y=u e^{\lambda x}$ then

$$
y^{\prime}=u^{\prime} e^{\lambda x}+\lambda u e^{\lambda x} \quad \text { and } \quad y^{\prime \prime}=u^{\prime \prime} e^{\lambda x}+2 \lambda u^{\prime} e^{\lambda x}+\lambda^{2} u e^{\lambda x} .
$$

Substituting these expressions into (5.2.18) and dropping the common factor $e^{\lambda x}$ yields

$$
\left(u^{\prime \prime}+2 \lambda u^{\prime}+\lambda^{2} u\right)-2 \lambda\left(u^{\prime}+\lambda u\right)+\left(\lambda^{2}+\omega^{2}\right) u=0
$$

which simplifies to

$$
u^{\prime \prime}+\omega^{2} u=0
$$

From Example 5.1.2, the general solution of this equation is

$$
u=c_{1} \cos \omega x+c_{2} \sin \omega x
$$

Therefore any function of the form

$$
\begin{equation*}
y=e^{\lambda x}\left(c_{1} \cos \omega x+c_{2} \sin \omega x\right) \tag{5.2.19}
\end{equation*}
$$

is a solution of (5.2.18). Letting $c_{1}=1$ and $c_{2}=0$ here yields the solution $y_{1}=e^{\lambda x} \cos \omega x$. Letting $c_{1}=0$ and $c_{2}=1$ yields a second solution $y_{2}=e^{\lambda x} \sin \omega x$. Since $y_{2} / y_{1}=\tan \omega x$ is nonconstant, so Theorem 5.1.6 implies that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.2.18), and (5.2.19) is the general solution.

Summary
The next theorem summarizes the results of this section.
Theorem 5.2.1 Let $p(r)=a r^{2}+b r+c$ be the characteristic polynomial of

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 . \tag{5.2.20}
\end{equation*}
$$

Then:
(a) If $p(r)=0$ has distinct real roots $r_{1}$ and $r_{2}$, then the general solution of (5.2.20) is

$$
y=c_{1} e^{r_{1} x}+c_{2} e^{r_{2} x} .
$$

(b) If $p(r)=0$ has a repeated root $r_{1}$, then the general solution of (5.2.20) is

$$
y=e^{r_{1} x}\left(c_{1}+c_{2} x\right)
$$

(c) If $p(r)=0$ has complex conjugate roots $r_{1}=\lambda+i \omega$ and $r_{2}=\lambda-i \omega(w h e r e \omega>0)$, then the general solution of (5.2.20) is

$$
y=e^{\lambda x}\left(c_{1} \cos \omega x+c_{2} \sin \omega x\right)
$$

### 5.2 Exercises

In Exercises 1-12 find the general solution.

1. $y^{\prime \prime}+5 y^{\prime}-6 y=0$
2. $y^{\prime \prime}-4 y^{\prime}+5 y=0$
3. $y^{\prime \prime}+8 y^{\prime}+7 y=0$
4. $y^{\prime \prime}-4 y^{\prime}+4 y=0$
5. $y^{\prime \prime}+2 y^{\prime}+10 y=0$
6. $y^{\prime \prime}+6 y^{\prime}+10 y=0$
7. $y^{\prime \prime}-8 y^{\prime}+16 y=0$
8. $y^{\prime \prime}+y^{\prime}=0$
9. $y^{\prime \prime}-2 y^{\prime}+3 y=0$
10. $y^{\prime \prime}+6 y^{\prime}+13 y=0$
11. $4 y^{\prime \prime}+4 y^{\prime}+10 y=0$
12. $10 y^{\prime \prime}-3 y^{\prime}-y=0$

In Exercises 13-17 solve the initial value problem.
13. $y^{\prime \prime}+14 y^{\prime}+50 y=0, \quad y(0)=2, \quad y^{\prime}(0)=-17$
14. $6 y^{\prime \prime}-y^{\prime}-y=0, \quad y(0)=10, \quad y^{\prime}(0)=0$
15. $6 y^{\prime \prime}+y^{\prime}-y=0, \quad y(0)=-1, \quad y^{\prime}(0)=3$
16. $4 y^{\prime \prime}-4 y^{\prime}-3 y=0, \quad y(0)=\frac{13}{12}, \quad y^{\prime}(0)=\frac{23}{24}$
17. $4 y^{\prime \prime}-12 y^{\prime}+9 y=0, \quad y(0)=3, \quad y^{\prime}(0)=\frac{5}{2}$

In Exercises 18-21 solve the initial value problem and graph the solution.
18. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+7 y^{\prime}+12 y=0, \quad y(0)=-1, \quad y^{\prime}(0)=0$
19. $\mathrm{C} / \mathrm{G} y^{\prime \prime}-6 y^{\prime}+9 y=0, \quad y(0)=0, \quad y^{\prime}(0)=2$
20. $\mathrm{C} / \mathrm{G} 36 y^{\prime \prime}-12 y^{\prime}+y=0, \quad y(0)=3, \quad y^{\prime}(0)=\frac{5}{2}$
21. $\mathrm{C} / \mathrm{G} \quad y^{\prime \prime}+4 y^{\prime}+10 y=0, \quad y(0)=3, \quad y^{\prime}(0)=-2$
22. (a) Suppose $y$ is a solution of the constant coefficient homogeneous equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{A}
\end{equation*}
$$

Let $z(x)=y\left(x-x_{0}\right)$, where $x_{0}$ is an arbitrary real number. Show that

$$
a z^{\prime \prime}+b z^{\prime}+c z=0
$$

(b) Let $z_{1}(x)=y_{1}\left(x-x_{0}\right)$ and $z_{2}(x)=y_{2}\left(x-x_{0}\right)$, where $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (A). Show that $\left\{z_{1}, z_{2}\right\}$ is also a fundamental set of solutions of (A).
(c) The statement of Theorem 5.2.1 is convenient for solving an initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

where the initial conditions are imposed at $x_{0}=0$. However, if the initial value problem is

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1} \tag{B}
\end{equation*}
$$

where $x_{0} \neq 0$, then determining the constants in

$$
y=c_{1} e^{r_{1} x}+c_{2} e^{r_{2} x}, \quad y=e^{r_{1} x}\left(c_{1}+c_{2} x\right), \text { or } y=e^{\lambda x}\left(c_{1} \cos \omega x+c_{2} \sin \omega x\right)
$$

(whichever is applicable) is more complicated. Use (b) to restate Theorem 5.2.1 in a form more convenient for solving (B).

In Exercises 23-28 use a method suggested by Exercise 22 to solve the initial value problem.
23. $y^{\prime \prime}+3 y^{\prime}+2 y=0, \quad y(1)=-1, \quad y^{\prime}(1)=4$
24. $y^{\prime \prime}-6 y^{\prime}-7 y=0, \quad y(2)=-\frac{1}{3}, \quad y^{\prime}(2)=-5$
25. $\quad y^{\prime \prime}-14 y^{\prime}+49 y=0, \quad y(1)=2, \quad y^{\prime}(1)=11$
26. $9 y^{\prime \prime}+6 y^{\prime}+y=0, \quad y(2)=2, \quad y^{\prime}(2)=-\frac{14}{3}$
27. $9 y^{\prime \prime}+4 y=0, \quad y(\pi / 4)=2, \quad y^{\prime}(\pi / 4)=-2$
28. $\quad y^{\prime \prime}+3 y=0, \quad y(\pi / 3)=2, \quad y^{\prime}(\pi / 3)=-1$
29. Prove: If the characteristic equation of

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{A}
\end{equation*}
$$

has a repeated negative root or two roots with negative real parts, then every solution of (A) approaches zero as $x \rightarrow \infty$.
30. Suppose the characteristic polynomial of $a y^{\prime \prime}+b y^{\prime}+c y=0$ has distinct real roots $r_{1}$ and $r_{2}$. Use a method suggested by Exercise 22 to find a formula for the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1} .
$$

31. Suppose the characteristic polynomial of $a y^{\prime \prime}+b y^{\prime}+c y=0$ has a repeated real root $r_{1}$. Use a method suggested by Exercise 22 to find a formula for the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

32. Suppose the characteristic polynomial of $a y^{\prime \prime}+b y^{\prime}+c y=0$ has complex conjugate roots $\lambda \pm i \omega$. Use a method suggested by Exercise 22 to find a formula for the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

33. Suppose the characteristic equation of

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{A}
\end{equation*}
$$

has a repeated real root $r_{1}$. Temporarily, think of $e^{r x}$ as a function of two real variables $x$ and $r$.
(a) Show that

$$
\begin{equation*}
a \frac{\partial^{2}}{\partial^{2} x}\left(e^{r x}\right)+b \frac{\partial}{\partial x}\left(e^{r x}\right)+c e^{r x}=a\left(r-r_{1}\right)^{2} e^{r x} \tag{B}
\end{equation*}
$$

(b) Differentiate (B) with respect to $r$ to obtain

$$
\begin{equation*}
a \frac{\partial}{\partial r}\left(\frac{\partial^{2}}{\partial^{2} x}\left(e^{r x}\right)\right)+b \frac{\partial}{\partial r}\left(\frac{\partial}{\partial x}\left(e^{r x}\right)\right)+c\left(x e^{r x}\right)=\left[2+\left(r-r_{1}\right) x\right] a\left(r-r_{1}\right) e^{r x} \tag{C}
\end{equation*}
$$

(c) Reverse the orders of the partial differentiations in the first two terms on the left side of (C) to obtain

$$
\begin{equation*}
a \frac{\partial^{2}}{\partial x^{2}}\left(x e^{r x}\right)+b \frac{\partial}{\partial x}\left(x e^{r x}\right)+c\left(x e^{r x}\right)=\left[2+\left(r-r_{1}\right) x\right] a\left(r-r_{1}\right) e^{r x} . \tag{D}
\end{equation*}
$$

(d) Set $r=r_{1}$ in (B) and (D) to see that $y_{1}=e^{r_{1} x}$ and $y_{2}=x e^{r_{1} x}$ are solutions of (A)
34. In calculus you learned that $e^{u}, \cos u$, and $\sin u$ can be represented by the infinite series

$$
\begin{gather*}
e^{u}=\sum_{n=0}^{\infty} \frac{u^{n}}{n!}=1+\frac{u}{1!}+\frac{u^{2}}{2!}+\frac{u^{3}}{3!}+\cdots+\frac{u^{n}}{n!}+\cdots  \tag{A}\\
\cos u=\sum_{n=0}^{\infty}(-1)^{n} \frac{u^{2 n}}{(2 n)!}=1-\frac{u^{2}}{2!}+\frac{u^{4}}{4!}+\cdots+(-1)^{n} \frac{u^{2 n}}{(2 n)!}+\cdots, \tag{B}
\end{gather*}
$$

and

$$
\begin{equation*}
\sin u=\sum_{n=0}^{\infty}(-1)^{n} \frac{u^{2 n+1}}{(2 n+1)!}=u-\frac{u^{3}}{3!}+\frac{u^{5}}{5!}+\cdots+(-1)^{n} \frac{u^{2 n+1}}{(2 n+1)!}+\cdots \tag{C}
\end{equation*}
$$

for all real values of $u$. Even though you have previously considered (A) only for real values of $u$, we can set $u=i \theta$, where $\theta$ is real, to obtain

$$
\begin{equation*}
e^{i \theta}=\sum_{n=0}^{\infty} \frac{(i \theta)^{n}}{n!} . \tag{D}
\end{equation*}
$$

Given the proper background in the theory of infinite series with complex terms, it can be shown that the series in (D) converges for all real $\theta$.
(a) Recalling that $i^{2}=-1$, write enough terms of the sequence $\left\{i^{n}\right\}$ to convince yourself that the sequence is repetitive:

$$
1, i,-1,-i, 1, i,-1,-i, 1, i,-1,-i, 1, i,-1,-i, \cdots
$$

Use this to group the terms in (D) as

$$
\begin{aligned}
e^{i \theta} & =\left(1-\frac{\theta^{2}}{2}+\frac{\theta^{4}}{4}+\cdots\right)+i\left(\theta-\frac{\theta^{3}}{3!}+\frac{\theta^{5}}{5!}+\cdots\right) \\
& =\sum_{n=0}^{\infty}(-1)^{n} \frac{\theta^{2 n}}{(2 n)!}+i \sum_{n=0}^{\infty}(-1)^{n} \frac{\theta^{2 n+1}}{(2 n+1)!}
\end{aligned}
$$

By comparing this result with (B) and (C), conclude that

$$
\begin{equation*}
e^{i \theta}=\cos \theta+i \sin \theta \tag{E}
\end{equation*}
$$

This is Euler's identity.
(b) Starting from

$$
e^{i \theta_{1}} e^{i \theta_{2}}=\left(\cos \theta_{1}+i \sin \theta_{1}\right)\left(\cos \theta_{2}+i \sin \theta_{2}\right)
$$

collect the real part (the terms not multiplied by $i$ ) and the imaginary part (the terms multiplied by $i$ ) on the right, and use the trigonometric identities

$$
\begin{aligned}
\cos \left(\theta_{1}+\theta_{2}\right) & =\cos \theta_{1} \cos \theta_{2}-\sin \theta_{1} \sin \theta_{2} \\
\sin \left(\theta_{1}+\theta_{2}\right) & =\sin \theta_{1} \cos \theta_{2}+\cos \theta_{1} \sin \theta_{2}
\end{aligned}
$$

to verify that

$$
e^{i\left(\theta_{1}+\theta_{2}\right)}=e^{i \theta_{1}} e^{i \theta_{2}}
$$

as you would expect from the use of the exponential notation $e^{i \theta}$.
(c) If $\alpha$ and $\beta$ are real numbers, define

$$
\begin{equation*}
e^{\alpha+i \beta}=e^{\alpha} e^{i \beta}=e^{\alpha}(\cos \beta+i \sin \beta) . \tag{F}
\end{equation*}
$$

Show that if $z_{1}=\alpha_{1}+i \beta_{1}$ and $z_{2}=\alpha_{2}+i \beta_{2}$ then

$$
e^{z_{1}+z_{2}}=e^{z_{1}} e^{z_{2}}
$$

(d) Let $a, b$, and $c$ be real numbers, with $a \neq 0$. Let $z=u+i v$ where $u$ and $v$ are real-valued functions of $x$. Then we say that $z$ is a solution of

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{G}
\end{equation*}
$$

if $u$ and $v$ are both solutions of (G). Use Theorem 5.2.1(c) to verify that if the characteristic equation of $(\mathbf{G})$ has complex conjugate roots $\lambda \pm i \omega$ then $z_{1}=e^{(\lambda+i \omega) x}$ and $z_{2}=e^{(\lambda-i \omega) x}$ are both solutions of (G).

### 5.3 NONHOMOGENEOUS LINEAR EQUATIONS

We'll now consider the nonhomogeneous linear second order equation

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f(x) \tag{5.3.1}
\end{equation*}
$$

where the forcing function $f$ isn't identically zero. The next theorem, an extension of Theorem 5.1.1, gives sufficient conditions for existence and uniqueness of solutions of initial value problems for (5.3.1). We omit the proof, which is beyond the scope of this book.

Theorem 5.3.1 Suppose $p, q$ and $f$ are continuous on an open interval $(a, b)$, let $x_{0}$ be any point in $(a, b)$, and let $k_{0}$ and $k_{1}$ be arbitrary real numbers. Then the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f(x), \quad y\left(x_{0}\right)=k_{0}, \quad y^{\prime}\left(x_{0}\right)=k_{1}
$$

has a unique solution on $(a, b)$.
To find the general solution of (5.3.1) on an interval $(a, b)$ where $p, q$, and $f$ are continuous, it's necessary to find the general solution of the associated homogeneous equation

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.3.2}
\end{equation*}
$$

on $(a, b)$. We call (5.3.2) the complementary equation for (5.3.1).
The next theorem shows how to find the general solution of (5.3.1) if we know one solution $y_{p}$ of (5.3.1) and a fundamental set of solutions of (5.3.2). We call $y_{p}$ a particular solution of (5.3.1); it can be any solution that we can find, one way or another.

Theorem 5.3.2 Suppose $p, q$, and $f$ are continuous on $(a, b)$. Let $y_{p}$ be a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f(x) \tag{5.3.3}
\end{equation*}
$$

on $(a, b)$, and let $\left\{y_{1}, y_{2}\right\}$ be a fundamental set of solutions of the complementary equation

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0 \tag{5.3.4}
\end{equation*}
$$

on $(a, b)$. Then $y$ is a solution of (5.3.3) on $(a, b)$ if and only if

$$
\begin{equation*}
y=y_{p}+c_{1} y_{1}+c_{2} y_{2} \tag{5.3.5}
\end{equation*}
$$

where $c_{1}$ and $c_{2}$ are constants.
Proof We first show that $y$ in (5.3.5) is a solution of (5.3.3) for any choice of the constants $c_{1}$ and $c_{2}$. Differentiating (5.3.5) twice yields

$$
y^{\prime}=y_{p}^{\prime}+c_{1} y_{1}^{\prime}+c_{2} y_{2}^{\prime} \quad \text { and } \quad y^{\prime \prime}=y_{p}^{\prime \prime}+c_{1} y_{1}^{\prime \prime}+c_{2} y_{2}^{\prime \prime}
$$

so

$$
\begin{aligned}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y= & \left(y_{p}^{\prime \prime}+c_{1} y_{1}^{\prime \prime}+c_{2} y_{2}^{\prime \prime}\right)+p(x)\left(y_{p}^{\prime}+c_{1} y_{1}^{\prime}+c_{2} y_{2}^{\prime}\right) \\
& +q(x)\left(y_{p}+c_{1} y_{1}+c_{2} y_{2}\right) \\
= & \left(y_{p}^{\prime \prime}+p(x) y_{p}^{\prime}+q(x) y_{p}\right)+c_{1}\left(y_{1}^{\prime \prime}+p(x) y_{1}^{\prime}+q(x) y_{1}\right) \\
& +c_{2}\left(y_{2}^{\prime \prime}+p(x) y_{2}^{\prime}+q(x) y_{2}\right) \\
= & f+c_{1} \cdot 0+c_{2} \cdot 0=f
\end{aligned}
$$

since $y_{p}$ satisfies (5.3.3) and $y_{1}$ and $y_{2}$ satisfy (5.3.4).
Now we'll show that every solution of (5.3.3) has the form (5.3.5) for some choice of the constants $c_{1}$ and $c_{2}$. Suppose $y$ is a solution of (5.3.3). We'll show that $y-y_{p}$ is a solution of (5.3.4), and therefore of the form $y-y_{p}=c_{1} y_{1}+c_{2} y_{2}$, which implies (5.3.5). To see this, we compute

$$
\begin{aligned}
\left(y-y_{p}\right)^{\prime \prime}+p(x)\left(y-y_{p}\right)^{\prime}+q(x)\left(y-y_{p}\right)= & \left(y^{\prime \prime}-y_{p}^{\prime \prime}\right)+p(x)\left(y^{\prime}-y_{p}^{\prime}\right) \\
& +q(x)\left(y-y_{p}\right) \\
= & \left(y^{\prime \prime}+p(x) y^{\prime}+q(x) y\right) \\
& -\left(y_{p}^{\prime \prime}+p(x) y_{p}^{\prime}+q(x) y_{p}\right) \\
= & f(x)-f(x)=0,
\end{aligned}
$$

since $y$ and $y_{p}$ both satisfy (5.3.3).
We say that (5.3.5) is the general solution of (5.3.3) on $(a, b)$.
If $P_{0}, P_{1}$, and $F$ are continuous and $P_{0}$ has no zeros on $(a, b)$, then Theorem 5.3.2 implies that the general solution of

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F(x) \tag{5.3.6}
\end{equation*}
$$

on $(a, b)$ is $y=y_{p}+c_{1} y_{1}+c_{2} y_{2}$, where $y_{p}$ is a particular solution of (5.3.6) on $(a, b)$ and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0
$$

on $(a, b)$. To see this, we rewrite (5.3.6) as

$$
y^{\prime \prime}+\frac{P_{1}(x)}{P_{0}(x)} y^{\prime}+\frac{P_{2}(x)}{P_{0}(x)} y=\frac{F(x)}{P_{0}(x)}
$$

and apply Theorem 5.3.2 with $p=P_{1} / P_{0}, q=P_{2} / P_{0}$, and $f=F / P_{0}$.
To avoid awkward wording in examples and exercises, we won't specify the interval $(a, b)$ when we ask for the general solution of a specific linear second order equation, or for a fundamental set of solutions of a homogeneous linear second order equation. Let's agree that this always means that we want the general solution (or a fundamental set of solutions, as the case may be) on every open interval on which $p, q$, and $f$ are continuous if the equation is of the form (5.3.3), or on which $P_{0}, P_{1}, P_{2}$, and $F$ are continuous and $P_{0}$ has no zeros, if the equation is of the form (5.3.6). We leave it to you to identify these intervals in specific examples and exercises.

For completeness, we point out that if $P_{0}, P_{1}, P_{2}$, and $F$ are all continuous on an open interval $(a, b)$, but $P_{0}$ does have a zero in $(a, b)$, then (5.3.6) may fail to have a general solution on $(a, b)$ in the sense just defined. Exercises 42-44 illustrate this point for a homogeneous equation.

In this section we to limit ourselves to applications of Theorem 5.3.2 where we can guess at the form of the particular solution.

## Example 5.3.1

(a) Find the general solution of

$$
\begin{equation*}
y^{\prime \prime}+y=1 \tag{5.3.7}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+y=1, \quad y(0)=2, \quad y^{\prime}(0)=7 . \tag{5.3.8}
\end{equation*}
$$

 and $f \equiv 1$ in (5.3.7) are continuous on $(-\infty, \infty)$. By inspection we see that $y_{p} \equiv 1$ is a particular solution of (5.3.7). Since $y_{1}=\cos x$ and $y_{2}=\sin x$ form a fundamental set of solutions of the complementary equation $y^{\prime \prime}+y=0$, the general solution of (5.3.7) is

$$
\begin{equation*}
y=1+c_{1} \cos x+c_{2} \sin x \tag{5.3.9}
\end{equation*}
$$

$\underline{\text { SOLUTION }(\mathbf{b})}$ Imposing the initial condition $y(0)=2$ in (5.3.9) yields $2=1+c_{1}$, so $c_{1}=1$. Differentiating (5.3.9) yields

$$
y^{\prime}=-c_{1} \sin x+c_{2} \cos x
$$



Figure 5.3.1 $y=1+\cos x+7 \sin x$

Imposing the initial condition $y^{\prime}(0)=7$ here yields $c_{2}=7$, so the solution of (5.3.8) is

$$
y=1+\cos x+7 \sin x
$$

Figure 5.3.1 is a graph of this function.

## Example 5.3.2

(a) Find the general solution of

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+y=-3-x+x^{2} . \tag{5.3.10}
\end{equation*}
$$

(b) Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+y=-3-x+x^{2}, \quad y(0)=-2, \quad y^{\prime}(0)=1 . \tag{5.3.11}
\end{equation*}
$$

SOLUTION(a) The characteristic polynomial of the complementary equation

$$
y^{\prime \prime}-2 y^{\prime}+y=0
$$

is $r^{2}-2 r+1=(r-1)^{2}$, so $y_{1}=e^{x}$ and $y_{2}=x e^{x}$ form a fundamental set of solutions of the complementary equation. To guess a form for a particular solution of (5.3.10), we note that substituting a second degree polynomial $y_{p}=A+B x+C x^{2}$ into the left side of (5.3.10) will produce another second degree polynomial with coefficients that depend upon $A, B$, and $C$. The trick is to choose $A, B$, and $C$ so the polynomials on the two sides of (5.3.10) have the same coefficients; thus, if

$$
y_{p}=A+B x+C x^{2} \quad \text { then } \quad y_{p}^{\prime}=B+2 C x \quad \text { and } \quad y_{p}^{\prime \prime}=2 C
$$

so

$$
\begin{aligned}
y_{p}^{\prime \prime}-2 y_{p}^{\prime}+y_{p} & =2 C-2(B+2 C x)+\left(A+B x+C x^{2}\right) \\
& =(2 C-2 B+A)+(-4 C+B) x+C x^{2}=-3-x+x^{2}
\end{aligned}
$$

Equating coefficients of like powers of $x$ on the two sides of the last equality yields

$$
\begin{array}{rlr}
C & =1 \\
B-4 C & = & -1 \\
A-2 B+2 C & = & -3,
\end{array}
$$



Figure 5.3.2 $y=1+3 x+x^{2}-e^{x}(3-x)$
so $C=1, B=-1+4 C=3$, and $A=-3-2 C+2 B=1$. Therefore $y_{p}=1+3 x+x^{2}$ is a particular solution of (5.3.10) and Theorem 5.3.2 implies that

$$
\begin{equation*}
y=1+3 x+x^{2}+e^{x}\left(c_{1}+c_{2} x\right) \tag{5.3.12}
\end{equation*}
$$

is the general solution of (5.3.10).
$\underline{\text { SOLUTION }(\mathbf{b})}$ Imposing the initial condition $y(0)=-2$ in (5.3.12) yields $-2=1+c_{1}$, so $c_{1}=-3$. Differentiating (5.3.12) yields

$$
y^{\prime}=3+2 x+e^{x}\left(c_{1}+c_{2} x\right)+c_{2} e^{x},
$$

and imposing the initial condition $y^{\prime}(0)=1$ here yields $1=3+c_{1}+c_{2}$, so $c_{2}=1$. Therefore the solution of (5.3.11) is

$$
y=1+3 x+x^{2}-e^{x}(3-x) .
$$

Figure 5.3.2 is a graph of this solution.
Example 5.3.3 Find the general solution of

$$
\begin{equation*}
x^{2} y^{\prime \prime}+x y^{\prime}-4 y=2 x^{4} \tag{5.3.13}
\end{equation*}
$$

on $(-\infty, 0)$ and $(0, \infty)$.

Solution In Example 5.1.3, we verified that $y_{1}=x^{2}$ and $y_{2}=1 / x^{2}$ form a fundamental set of solutions of the complementary equation

$$
x^{2} y^{\prime \prime}+x y^{\prime}-4 y=0
$$

on $(-\infty, 0)$ and $(0, \infty)$. To find a particular solution of (5.3.13), we note that if $y_{p}=A x^{4}$, where $A$ is a constant then both sides of (5.3.13) will be constant multiples of $x^{4}$ and we may be able to choose $A$ so the two sides are equal. This is true in this example, since if $y_{p}=A x^{4}$ then

$$
x^{2} y_{p}^{\prime \prime}+x y_{p}^{\prime}-4 y_{p}=x^{2}\left(12 A x^{2}\right)+x\left(4 A x^{3}\right)-4 A x^{4}=12 A x^{4}=2 x^{4}
$$

if $A=1 / 6$; therefore, $y_{p}=x^{4} / 6$ is a particular solution of (5.3.13) on $(-\infty, \infty)$. Theorem 5.3.2 implies that the general solution of $(5.3 .13)$ on $(-\infty, 0)$ and $(0, \infty)$ is

$$
y=\frac{x^{4}}{6}+c_{1} x^{2}+\frac{c_{2}}{x^{2}}
$$

The Principle of Superposition
The next theorem enables us to break a nonhomogeous equation into simpler parts, find a particular solution for each part, and then combine their solutions to obtain a particular solution of the original problem.

Theorem 5.3.3 [The Principle of Superposition] Suppose $y_{p_{1}}$ is a particular solution of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f_{1}(x)
$$

on $(a, b)$ and $y_{p_{2}}$ is a particular solution of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f_{2}(x)
$$

on $(a, b)$. Then

$$
y_{p}=y_{p_{1}}+y_{p_{2}}
$$

is a particular solution of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f_{1}(x)+f_{2}(x)
$$

on $(a, b)$.
Proof If $y_{p}=y_{p_{1}}+y_{p_{2}}$ then

$$
\begin{aligned}
y_{p}^{\prime \prime}+p(x) y_{p}^{\prime}+q(x) y_{p} & =\left(y_{p_{1}}+y_{p_{2}}\right)^{\prime \prime}+p(x)\left(y_{p_{1}}+y_{p_{2}}\right)^{\prime}+q(x)\left(y_{p_{1}}+y_{p_{2}}\right) \\
& =\left(y_{p_{1}}^{\prime \prime}+p(x) y_{p_{1}}^{\prime}+q(x) y_{p_{1}}\right)+\left(y_{p_{2}}^{\prime \prime}+p(x) y_{p_{2}}^{\prime}+q(x) y_{p_{2}}\right) \\
& =f_{1}(x)+f_{2}(x)
\end{aligned}
$$

It's easy to generalize Theorem 5.3.3 to the equation

$$
\begin{equation*}
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f(x) \tag{5.3.14}
\end{equation*}
$$

where

$$
f=f_{1}+f_{2}+\cdots+f_{k} ;
$$

thus, if $y_{p_{i}}$ is a particular solution of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f_{i}(x)
$$

on $(a, b)$ for $i=1,2, \ldots, k$, then $y_{p_{1}}+y_{p_{2}}+\cdots+y_{p_{k}}$ is a particular solution of (5.3.14) on $(a, b)$. Moreover, by a proof similar to the proof of Theorem 5.3.3 we can formulate the principle of superposition in terms of a linear equation written in the form

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F(x)
$$

(Exercise 39); that is, if $y_{p_{1}}$ is a particular solution of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F_{1}(x)
$$

on $(a, b)$ and $y_{p_{2}}$ is a particular solution of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F_{2}(x)
$$

on $(a, b)$, then $y_{p_{1}}+y_{p_{2}}$ is a solution of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F_{1}(x)+F_{2}(x)
$$

on ( $a, b$ ).
Example 5.3.4 The function $y_{p_{1}}=x^{4} / 15$ is a particular solution of

$$
\begin{equation*}
x^{2} y^{\prime \prime}+4 x y^{\prime}+2 y=2 x^{4} \tag{5.3.15}
\end{equation*}
$$

on $(-\infty, \infty)$ and $y_{p_{2}}=x^{2} / 3$ is a particular solution of

$$
\begin{equation*}
x^{2} y^{\prime \prime}+4 x y^{\prime}+2 y=4 x^{2} \tag{5.3.16}
\end{equation*}
$$

on $(-\infty, \infty)$. Use the principle of superposition to find a particular solution of

$$
\begin{equation*}
x^{2} y^{\prime \prime}+4 x y^{\prime}+2 y=2 x^{4}+4 x^{2} \tag{5.3.17}
\end{equation*}
$$

on $(-\infty, \infty)$.

Solution The right side $F(x)=2 x^{4}+4 x^{2}$ in (5.3.17) is the sum of the right sides

$$
F_{1}(x)=2 x^{4} \quad \text { and } \quad F_{2}(x)=4 x^{2} .
$$

in (5.3.15) and (5.3.16). Therefore the principle of superposition implies that

$$
y_{p}=y_{p_{1}}+y_{p_{2}}=\frac{x^{4}}{15}+\frac{x^{2}}{3}
$$

is a particular solution of (5.3.17).

### 5.3 Exercises

In Exercises 1-6 find a particular solution by the method used in Example 5.3.2. Then find the general solution and, where indicated, solve the initial value problem and graph the solution.

1. $y^{\prime \prime}+5 y^{\prime}-6 y=22+18 x-18 x^{2}$
2. $y^{\prime \prime}-4 y^{\prime}+5 y=1+5 x$
3. $y^{\prime \prime}+8 y^{\prime}+7 y=-8-x+24 x^{2}+7 x^{3}$
4. $y^{\prime \prime}-4 y^{\prime}+4 y=2+8 x-4 x^{2}$
5. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+2 y^{\prime}+10 y=4+26 x+6 x^{2}+10 x^{3}, \quad y(0)=2, \quad y^{\prime}(0)=9$
6. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+6 y^{\prime}+10 y=22+20 x, \quad y(0)=2, y^{\prime}(0)=-2$
7. Show that the method used in Example 5.3.2 won't yield a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+y^{\prime}=1+2 x+x^{2} \tag{A}
\end{equation*}
$$

that is, (A) does'nt have a particular solution of the form $y_{p}=A+B x+C x^{2}$, where $A, B$, and $C$ are constants.

## In Exercises 8-13 find a particular solution by the method used in Example 5.3.3.

8. $x^{2} y^{\prime \prime}+7 x y^{\prime}+8 y=\frac{6}{x}$
9. $x^{2} y^{\prime \prime}-7 x y^{\prime}+7 y=13 x^{1 / 2}$
10. $x^{2} y^{\prime \prime}-x y^{\prime}+y=2 x^{3}$
11. $x^{2} y^{\prime \prime}+5 x y^{\prime}+4 y=\frac{1}{x^{3}}$
12. $x^{2} y^{\prime \prime}+x y^{\prime}+y=10 x^{1 / 3}$
13. $x^{2} y^{\prime \prime}-3 x y^{\prime}+13 y=2 x^{4}$
14. Show that the method suggested for finding a particular solution in Exercises 8-13 won't yield a particular solution of

$$
\begin{equation*}
x^{2} y^{\prime \prime}+3 x y^{\prime}-3 y=\frac{1}{x^{3}} \tag{A}
\end{equation*}
$$

that is, (A) doesn't have a particular solution of the form $y_{p}=A / x^{3}$.
15. Prove: If $a, b, c, \alpha$, and $M$ are constants and $M \neq 0$ then

$$
a x^{2} y^{\prime \prime}+b x y^{\prime}+c y=M x^{\alpha}
$$

has a particular solution $y_{p}=A x^{\alpha}(A=$ constant $)$ if and only if $a \alpha(\alpha-1)+b \alpha+c \neq 0$.
If $a, b, c$, and $\alpha$ are constants, then

$$
a\left(e^{\alpha x}\right)^{\prime \prime}+b\left(e^{\alpha x}\right)^{\prime}+c e^{\alpha x}=\left(a \alpha^{2}+b \alpha+c\right) e^{\alpha x} .
$$

Use this in Exercises 16-21 to find a particular solution. Then find the general solution and, where indicated, solve the initial value problem and graph the solution.
16. $y^{\prime \prime}+5 y^{\prime}-6 y=6 e^{3 x}$
17. $y^{\prime \prime}-4 y^{\prime}+5 y=e^{2 x}$
18. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+8 y^{\prime}+7 y=10 e^{-2 x}, \quad y(0)=-2, y^{\prime}(0)=10$
19. $\mathrm{C} / \mathrm{G} y^{\prime \prime}-4 y^{\prime}+4 y=e^{x}, \quad y(0)=2, \quad y^{\prime}(0)=0$
20. $y^{\prime \prime}+2 y^{\prime}+10 y=e^{x / 2}$
21. $y^{\prime \prime}+6 y^{\prime}+10 y=e^{-3 x}$
22. Show that the method suggested for finding a particular solution in Exercises 16-21 won't yield a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-7 y^{\prime}+12 y=5 e^{4 x} \tag{A}
\end{equation*}
$$

that is, (A) doesn't have a particular solution of the form $y_{p}=A e^{4 x}$.
23. Prove: If $\alpha$ and $M$ are constants and $M \neq 0$ then constant coefficient equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=M e^{\alpha x}
$$

has a particular solution $y_{p}=A e^{\alpha x}(A=$ constant $)$ if and only if $e^{\alpha x}$ isn't a solution of the complementary equation.

If $\omega$ is a constant, differentiating a linear combination of $\cos \omega x$ and $\sin \omega x$ with respect to $x$ yields another linear combination of $\cos \omega x$ and $\sin \omega x$. In Exercises 24-29 use this to find a particular solution of the equation. Then find the general solution and, where indicated, solve the initial value problem and graph the solution.
24. $y^{\prime \prime}-8 y^{\prime}+16 y=23 \cos x-7 \sin x$
25. $y^{\prime \prime}+y^{\prime}=-8 \cos 2 x+6 \sin 2 x$
26. $y^{\prime \prime}-2 y^{\prime}+3 y=-6 \cos 3 x+6 \sin 3 x$
27. $y^{\prime \prime}+6 y^{\prime}+13 y=18 \cos x+6 \sin x$
28. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+7 y^{\prime}+12 y=-2 \cos 2 x+36 \sin 2 x, \quad y(0)=-3, \quad y^{\prime}(0)=3$
29. $\mathrm{C} / \mathrm{G} y^{\prime \prime}-6 y^{\prime}+9 y=18 \cos 3 x+18 \sin 3 x, \quad y(0)=2, \quad y^{\prime}(0)=2$
30. Find the general solution of

$$
y^{\prime \prime}+\omega_{0}^{2} y=M \cos \omega x+N \sin \omega x
$$

where $M$ and $N$ are constants and $\omega$ and $\omega_{0}$ are distinct positive numbers.
31. Show that the method suggested for finding a particular solution in Exercises 24-29 won't yield a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+y=\cos x+\sin x \tag{A}
\end{equation*}
$$

that is, (A) does not have a particular solution of the form $y_{p}=A \cos x+B \sin x$.
32. Prove: If $M, N$ are constants (not both zero) and $\omega>0$, the constant coefficient equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=M \cos \omega x+N \sin \omega x \tag{A}
\end{equation*}
$$

has a particular solution that's a linear combination of $\cos \omega x$ and $\sin \omega x$ if and only if the left side of (A) is not of the form $a\left(y^{\prime \prime}+\omega^{2} y\right)$, so that $\cos \omega x$ and $\sin \omega x$ are solutions of the complementary equation.

In Exercises 33-38 refer to the cited exercises and use the principal of superposition to find a particular solution. Then find the general solution.
33. $y^{\prime \prime}+5 y^{\prime}-6 y=22+18 x-18 x^{2}+6 e^{3 x}$ (See Exercises 1 and 16.)
34. $y^{\prime \prime}-4 y^{\prime}+5 y=1+5 x+e^{2 x}$ (See Exercises 2 and 17.)
35. $y^{\prime \prime}+8 y^{\prime}+7 y=-8-x+24 x^{2}+7 x^{3}+10 e^{-2 x}$ (See Exercises 3 and 18.)
36. $y^{\prime \prime}-4 y^{\prime}+4 y=2+8 x-4 x^{2}+e^{x}$ (See Exercises 4 and 19.)
37. $y^{\prime \prime}+2 y^{\prime}+10 y=4+26 x+6 x^{2}+10 x^{3}+e^{x / 2}$ (See Exercises 5 and 20.)
38. $y^{\prime \prime}+6 y^{\prime}+10 y=22+20 x+e^{-3 x}$ (See Exercises 6 and 21.)
39. Prove: If $y_{p_{1}}$ is a particular solution of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F_{1}(x)
$$

on $(a, b)$ and $y_{p_{2}}$ is a particular solution of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F_{2}(x)
$$

on $(a, b)$, then $y_{p}=y_{p_{1}}+y_{p_{2}}$ is a solution of

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F_{1}(x)+F_{2}(x)
$$

on $(a, b)$.
40. Suppose $p, q$, and $f$ are continuous on $(a, b)$. Let $y_{1}, y_{2}$, and $y_{p}$ be twice differentiable on $(a, b)$, such that $y=c_{1} y_{1}+c_{2} y_{2}+y_{p}$ is a solution of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f
$$

on $(a, b)$ for every choice of the constants $c_{1}, c_{2}$. Show that $y_{1}$ and $y_{2}$ are solutions of the complementary equation on $(a, b)$.

### 5.4 THE METHOD OF UNDETERMINED COEFFICIENTS I

In this section we consider the constant coefficient equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=e^{\alpha x} G(x) \tag{5.4.1}
\end{equation*}
$$

where $\alpha$ is a constant and $G$ is a polynomial.
From Theorem 5.3.2, the general solution of (5.4.1) is $y=y_{p}+c_{1} y_{1}+c_{2} y_{2}$, where $y_{p}$ is a particular solution of (5.4.1) and $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of the complementary equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=0 .
$$

In Section 5.2 we showed how to find $\left\{y_{1}, y_{2}\right\}$. In this section we'll show how to find $y_{p}$. The procedure that we'll use is called the method of undetermined coefficients.

Our first example is similar to Exercises 16-21.
Example 5.4.1 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-7 y^{\prime}+12 y=4 e^{2 x} . \tag{5.4.2}
\end{equation*}
$$

Then find the general solution.

Solution Substituting $y_{p}=A e^{2 x}$ for $y$ in (5.4.2) will produce a constant multiple of $A e^{2 x}$ on the left side of (5.4.2), so it may be possible to choose $A$ so that $y_{p}$ is a solution of (5.4.2). Let's try it; if $y_{p}=A e^{2 x}$ then

$$
y_{p}^{\prime \prime}-7 y_{p}^{\prime}+12 y_{p}=4 A e^{2 x}-14 A e^{2 x}+12 A e^{2 x}=2 A e^{2 x}=4 e^{2 x}
$$

if $A=2$. Therefore $y_{p}=2 e^{2 x}$ is a particular solution of (5.4.2). To find the general solution, we note that the characteristic polynomial of the complementary equation

$$
\begin{equation*}
y^{\prime \prime}-7 y^{\prime}+12 y=0 \tag{5.4.3}
\end{equation*}
$$

is $p(r)=r^{2}-7 r+12=(r-3)(r-4)$, so $\left\{e^{3 x}, e^{4 x}\right\}$ is a fundamental set of solutions of (5.4.3). Therefore the general solution of (5.4.2) is

$$
y=2 e^{2 x}+c_{1} e^{3 x}+c_{2} e^{4 x}
$$

Example 5.4.2 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-7 y^{\prime}+12 y=5 e^{4 x} \tag{5.4.4}
\end{equation*}
$$

Then find the general solution.

Solution Fresh from our success in finding a particular solution of (5.4.2) - where we chose $y_{p}=A e^{2 x}$ because the right side of (5.4.2) is a constant multiple of $e^{2 x}$ - it may seem reasonable to try $y_{p}=A e^{4 x}$ as a particular solution of (5.4.4). However, this won't work, since we saw in Example 5.4.1 that $e^{4 x}$ is a solution of the complementary equation (5.4.3), so substituting $y_{p}=A e^{4 x}$ into the left side of (5.4.4) produces zero on the left, no matter how we choose $A$. To discover a suitable form for $y_{p}$, we use the same approach that we used in Section 5.2 to find a second solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=0
$$

in the case where the characteristic equation has a repeated real root: we look for solutions of (5.4.4) in the form $y=u e^{4 x}$, where $u$ is a function to be determined. Substituting

$$
\begin{equation*}
y=u e^{4 x}, \quad y^{\prime}=u^{\prime} e^{4 x}+4 u e^{4 x}, \quad \text { and } \quad y^{\prime \prime}=u^{\prime \prime} e^{4 x}+8 u^{\prime} e^{4 x}+16 u e^{4 x} \tag{5.4.5}
\end{equation*}
$$

into (5.4.4) and canceling the common factor $e^{4 x}$ yields

$$
\left(u^{\prime \prime}+8 u^{\prime}+16 u\right)-7\left(u^{\prime}+4 u\right)+12 u=5
$$

or

$$
u^{\prime \prime}+u^{\prime}=5 .
$$

By inspection we see that $u_{p}=5 x$ is a particular solution of this equation, so $y_{p}=5 x e^{4 x}$ is a particular solution of (5.4.4). Therefore

$$
y=5 x e^{4 x}+c_{1} e^{3 x}+c_{2} e^{4 x}
$$

is the general solution.
Example 5.4.3 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-8 y^{\prime}+16 y=2 e^{4 x} \tag{5.4.6}
\end{equation*}
$$

Solution Since the characteristic polynomial of the complementary equation

$$
\begin{equation*}
y^{\prime \prime}-8 y^{\prime}+16 y=0 \tag{5.4.7}
\end{equation*}
$$

is $p(r)=r^{2}-8 r+16=(r-4)^{2}$, both $y_{1}=e^{4 x}$ and $y_{2}=x e^{4 x}$ are solutions of (5.4.7). Therefore (5.4.6) does not have a solution of the form $y_{p}=A e^{4 x}$ or $y_{p}=A x e^{4 x}$. As in Example 5.4.2, we look for solutions of (5.4.6) in the form $y=u e^{4 x}$, where $u$ is a function to be determined. Substituting from (5.4.5) into (5.4.6) and canceling the common factor $e^{4 x}$ yields

$$
\left(u^{\prime \prime}+8 u^{\prime}+16 u\right)-8\left(u^{\prime}+4 u\right)+16 u=2
$$

or

$$
u^{\prime \prime}=2
$$

Integrating twice and taking the constants of integration to be zero shows that $u_{p}=x^{2}$ is a particular solution of this equation, so $y_{p}=x^{2} e^{4 x}$ is a particular solution of (5.4.4). Therefore

$$
y=e^{4 x}\left(x^{2}+c_{1}+c_{2} x\right)
$$

is the general solution.
The preceding examples illustrate the following facts concerning the form of a particular solution $y_{p}$ of a constant coefficent equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=k e^{\alpha x}
$$

where $k$ is a nonzero constant:
(a) If $e^{\alpha x}$ isn't a solution of the complementary equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{5.4.8}
\end{equation*}
$$

then $y_{p}=A e^{\alpha x}$, where $A$ is a constant. (See Example 5.4.1).
(b) If $e^{\alpha x}$ is a solution of (5.4.8) but $x e^{\alpha x}$ is not, then $y_{p}=A x e^{\alpha x}$, where $A$ is a constant. (See Example 5.4.2.)
(c) If both $e^{\alpha x}$ and $x e^{\alpha x}$ are solutions of (5.4.8), then $y_{p}=A x^{2} e^{\alpha x}$, where $A$ is a constant. (See Example 5.4.3.)
See Exercise 30 for the proofs of these facts.
In all three cases you can just substitute the appropriate form for $y_{p}$ and its derivatives directly into

$$
a y_{p}^{\prime \prime}+b y_{p}^{\prime}+c y_{p}=k e^{\alpha x}
$$

and solve for the constant $A$, as we did in Example 5.4.1. (See Exercises 31-33.) However, if the equation is

$$
a y^{\prime \prime}+b y^{\prime}+c y=k e^{\alpha x} G(x)
$$

where $G$ is a polynomial of degree greater than zero, we recommend that you use the substitution $y=$ $u e^{\alpha x}$ as we did in Examples 5.4.2 and 5.4.3. The equation for $u$ will turn out to be

$$
\begin{equation*}
a u^{\prime \prime}+p^{\prime}(\alpha) u^{\prime}+p(\alpha) u=G(x) \tag{5.4.9}
\end{equation*}
$$

where $p(r)=a r^{2}+b r+c$ is the characteristic polynomial of the complementary equation and $p^{\prime}(r)=$ $2 a r+b$ (Exercise 30); however, you shouldn't memorize this since it's easy to derive the equation for $u$ in any particular case. Note, however, that if $e^{\alpha x}$ is a solution of the complementary equation then $p(\alpha)=0$, so (5.4.9) reduces to

$$
a u^{\prime \prime}+p^{\prime}(\alpha) u^{\prime}=G(x)
$$

while if both $e^{\alpha x}$ and $x e^{\alpha x}$ are solutions of the complementary equation then $p(r)=a(r-\alpha)^{2}$ and $p^{\prime}(r)=2 a(r-\alpha)$, so $p(\alpha)=p^{\prime}(\alpha)=0$ and (5.4.9) reduces to

$$
a u^{\prime \prime}=G(x) .
$$

Example 5.4.4 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-3 y^{\prime}+2 y=e^{3 x}\left(-1+2 x+x^{2}\right) \tag{5.4.10}
\end{equation*}
$$

Solution Substituting

$$
y=u e^{3 x}, \quad y^{\prime}=u^{\prime} e^{3 x}+3 u e^{3 x}, \quad \text { and } y^{\prime \prime}=u^{\prime \prime} e^{3 x}+6 u^{\prime} e^{3 x}+9 u e^{3 x}
$$

into (5.4.10) and canceling $e^{3 x}$ yields

$$
\left(u^{\prime \prime}+6 u^{\prime}+9 u\right)-3\left(u^{\prime}+3 u\right)+2 u=-1+2 x+x^{2}
$$

or

$$
\begin{equation*}
u^{\prime \prime}+3 u^{\prime}+2 u=-1+2 x+x^{2} \tag{5.4.11}
\end{equation*}
$$

As in Example 2, in order to guess a form for a particular solution of (5.4.11), we note that substituting a second degree polynomial $u_{p}=A+B x+C x^{2}$ for $u$ in the left side of (5.4.11) produces another second degree polynomial with coefficients that depend upon $A, B$, and $C$; thus,

$$
\text { if } \quad u_{p}=A+B x+C x^{2} \quad \text { then } \quad u_{p}^{\prime}=B+2 C x \quad \text { and } \quad u_{p}^{\prime \prime}=2 C .
$$

If $u_{p}$ is to satisfy (5.4.11), we must have

$$
\begin{aligned}
u_{p}^{\prime \prime}+3 u_{p}^{\prime}+2 u_{p} & =2 C+3(B+2 C x)+2\left(A+B x+C x^{2}\right) \\
& =(2 C+3 B+2 A)+(6 C+2 B) x+2 C x^{2}=-1+2 x+x^{2}
\end{aligned}
$$

Equating coefficients of like powers of $x$ on the two sides of the last equality yields

$$
\begin{aligned}
2 C & =1 \\
2 B+6 C & =2 \\
2 A+3 B+2 C & =-1 .
\end{aligned}
$$

Solving these equations for $C, B$, and $A$ (in that order) yields $C=1 / 2, B=-1 / 2, A=-1 / 4$. Therefore

$$
u_{p}=-\frac{1}{4}\left(1+2 x-2 x^{2}\right)
$$

is a particular solution of (5.4.11), and

$$
y_{p}=u_{p} e^{3 x}=-\frac{e^{3 x}}{4}\left(1+2 x-2 x^{2}\right)
$$

is a particular solution of (5.4.10).

Example 5.4.5 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-4 y^{\prime}+3 y=e^{3 x}\left(6+8 x+12 x^{2}\right) \tag{5.4.12}
\end{equation*}
$$

Solution Substituting

$$
y=u e^{3 x}, \quad y^{\prime}=u^{\prime} e^{3 x}+3 u e^{3 x}, \quad \text { and } y^{\prime \prime}=u^{\prime \prime} e^{3 x}+6 u^{\prime} e^{3 x}+9 u e^{3 x}
$$

into (5.4.12) and canceling $e^{3 x}$ yields

$$
\left(u^{\prime \prime}+6 u^{\prime}+9 u\right)-4\left(u^{\prime}+3 u\right)+3 u=6+8 x+12 x^{2}
$$

or

$$
\begin{equation*}
u^{\prime \prime}+2 u^{\prime}=6+8 x+12 x^{2} \tag{5.4.13}
\end{equation*}
$$

There's no $u$ term in this equation, since $e^{3 x}$ is a solution of the complementary equation for (5.4.12). (See Exercise 30.) Therefore (5.4.13) does not have a particular solution of the form $u_{p}=A+B x+C x^{2}$ that we used successfully in Example 5.4.4, since with this choice of $u_{p}$,

$$
u_{p}^{\prime \prime}+2 u_{p}^{\prime}=2 C+(B+2 C x)
$$

can't contain the last term $\left(12 x^{2}\right)$ on the right side of (5.4.13). Instead, let's try $u_{p}=A x+B x^{2}+C x^{3}$ on the grounds that

$$
u_{p}^{\prime}=A+2 B x+3 C x^{2} \quad \text { and } \quad u_{p}^{\prime \prime}=2 B+6 C x
$$

together contain all the powers of $x$ that appear on the right side of (5.4.13).
Substituting these expressions in place of $u^{\prime}$ and $u^{\prime \prime}$ in (5.4.13) yields

$$
(2 B+6 C x)+2\left(A+2 B x+3 C x^{2}\right)=(2 B+2 A)+(6 C+4 B) x+6 C x^{2}=6+8 x+12 x^{2} .
$$

Comparing coefficients of like powers of $x$ on the two sides of the last equality shows that $u_{p}$ satisfies (5.4.13) if

$$
\begin{aligned}
6 C & =12 \\
4 B+6 C & =8 \\
2 A+2 B & =6 .
\end{aligned}
$$

Solving these equations successively yields $C=2, B=-1$, and $A=4$. Therefore

$$
u_{p}=x\left(4-x+2 x^{2}\right)
$$

is a particular solution of (5.4.13), and

$$
y_{p}=u_{p} e^{3 x}=x e^{3 x}\left(4-x+2 x^{2}\right)
$$

is a particular solution of (5.4.12).
Example 5.4.6 Find a particular solution of

$$
\begin{equation*}
4 y^{\prime \prime}+4 y^{\prime}+y=e^{-x / 2}\left(-8+48 x+144 x^{2}\right) \tag{5.4.14}
\end{equation*}
$$

## Solution Substituting

$$
y=u e^{-x / 2}, \quad y^{\prime}=u^{\prime} e^{-x / 2}-\frac{1}{2} u e^{-x / 2}, \quad \text { and } \quad y^{\prime \prime}=u^{\prime \prime} e^{-x / 2}-u^{\prime} e^{-x / 2}+\frac{1}{4} u e^{-x / 2}
$$

into (5.4.14) and canceling $e^{-x / 2}$ yields

$$
4\left(u^{\prime \prime}-u^{\prime}+\frac{u}{4}\right)+4\left(u^{\prime}-\frac{u}{2}\right)+u=4 u^{\prime \prime}=-8+48 x+144 x^{2}
$$

or

$$
\begin{equation*}
u^{\prime \prime}=-2+12 x+36 x^{2}, \tag{5.4.15}
\end{equation*}
$$

which does not contain $u$ or $u^{\prime}$ because $e^{-x / 2}$ and $x e^{-x / 2}$ are both solutions of the complementary equation. (See Exercise 30.) To obtain a particular solution of (5.4.15) we integrate twice, taking the constants of integration to be zero; thus,

$$
u_{p}^{\prime}=-2 x+6 x^{2}+12 x^{3} \quad \text { and } \quad u_{p}=-x^{2}+2 x^{3}+3 x^{4}=x^{2}\left(-1+2 x+3 x^{2}\right)
$$

Therefore

$$
y_{p}=u_{p} e^{-x / 2}=x^{2} e^{-x / 2}\left(-1+2 x+3 x^{2}\right)
$$

is a particular solution of (5.4.14).
Summary
The preceding examples illustrate the following facts concerning particular solutions of a constant coefficent equation of the form

$$
a y^{\prime \prime}+b y^{\prime}+c y=e^{\alpha x} G(x),
$$

where $G$ is a polynomial (see Exercise 30):
(a) If $e^{\alpha x}$ isn't a solution of the complementary equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=0 \tag{5.4.16}
\end{equation*}
$$

then $y_{p}=e^{\alpha x} Q(x)$, where $Q$ is a polynomial of the same degree as $G$. (See Example 5.4.4).
(b) If $e^{\alpha x}$ is a solution of (5.4.16) but $x e^{\alpha x}$ is not, then $y_{p}=x e^{\alpha x} Q(x)$, where $Q$ is a polynomial of the same degree as $G$. (See Example 5.4.5.)
(c) If both $e^{\alpha x}$ and $x e^{\alpha x}$ are solutions of (5.4.16), then $y_{p}=x^{2} e^{\alpha x} Q(x)$, where $Q$ is a polynomial of the same degree as $G$. (See Example 5.4.6.)
In all three cases, you can just substitute the appropriate form for $y_{p}$ and its derivatives directly into

$$
a y_{p}^{\prime \prime}+b y_{p}^{\prime}+c y_{p}=e^{\alpha x} G(x)
$$

and solve for the coefficients of the polynomial $Q$. However, if you try this you will see that the computations are more tedious than those that you encounter by making the substitution $y=u e^{\alpha x}$ and finding a particular solution of the resulting equation for $u$. (See Exercises 34-36.) In Case (a) the equation for $u$ will be of the form

$$
a u^{\prime \prime}+p^{\prime}(\alpha) u^{\prime}+p(\alpha) u=G(x),
$$

with a particular solution of the form $u_{p}=Q(x)$, a polynomial of the same degree as $G$, whose coefficients can be found by the method used in Example 5.4.4. In Case (b) the equation for $u$ will be of the form

$$
a u^{\prime \prime}+p^{\prime}(\alpha) u^{\prime}=G(x)
$$

(no $u$ term on the left), with a particular solution of the form $u_{p}=x Q(x)$, where $Q$ is a polynomial of the same degree as $G$ whose coefficents can be found by the method used in Example 5.4.5. In Case (c) the equation for $u$ will be of the form

$$
a u^{\prime \prime}=G(x)
$$

with a particular solution of the form $u_{p}=x^{2} Q(x)$ that can be obtained by integrating $G(x) / a$ twice and taking the constants of integration to be zero, as in Example 5.4.6.
Using the Principle of Superposition
The next example shows how to combine the method of undetermined coefficients and Theorem 5.3.3, the principle of superposition.

Example 5.4.7 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-7 y^{\prime}+12 y=4 e^{2 x}+5 e^{4 x} \tag{5.4.17}
\end{equation*}
$$

Solution In Example 5.4.1 we found that $y_{p_{1}}=2 e^{2 x}$ is a particular solution of

$$
y^{\prime \prime}-7 y^{\prime}+12 y=4 e^{2 x}
$$

and in Example 5.4.2 we found that $y_{p_{2}}=5 x e^{4 x}$ is a particular solution of

$$
y^{\prime \prime}-7 y^{\prime}+12 y=5 e^{4 x}
$$

Therefore the principle of superposition implies that $y_{p}=2 e^{2 x}+5 x e^{4 x}$ is a particular solution of (5.4.17).

### 5.4 Exercises

In Exercises 1-14 find a particular solution.

1. $y^{\prime \prime}-3 y^{\prime}+2 y=e^{3 x}(1+x)$
2. $y^{\prime \prime}-6 y^{\prime}+5 y=e^{-3 x}(35-8 x)$
3. $y^{\prime \prime}-2 y^{\prime}-3 y=e^{x}(-8+3 x)$
4. $y^{\prime \prime}+2 y^{\prime}+y=e^{2 x}\left(-7-15 x+9 x^{2}\right)$
5. $y^{\prime \prime}+4 y=e^{-x}\left(7-4 x+5 x^{2}\right)$
6. $y^{\prime \prime}-y^{\prime}-2 y=e^{x}\left(9+2 x-4 x^{2}\right)$
7. $y^{\prime \prime}-4 y^{\prime}-5 y=-6 x e^{-x}$
8. $y^{\prime \prime}-3 y^{\prime}+2 y=e^{x}(3-4 x)$
9. $y^{\prime \prime}+y^{\prime}-12 y=e^{3 x}(-6+7 x)$
10. $2 y^{\prime \prime}-3 y^{\prime}-2 y=e^{2 x}(-6+10 x)$
11. $y^{\prime \prime}+2 y^{\prime}+y=e^{-x}(2+3 x)$
12. $y^{\prime \prime}-2 y^{\prime}+y=e^{x}(1-6 x)$
13. $y^{\prime \prime}-4 y^{\prime}+4 y=e^{2 x}\left(1-3 x+6 x^{2}\right)$
14. $9 y^{\prime \prime}+6 y^{\prime}+y=e^{-x / 3}\left(2-4 x+4 x^{2}\right)$

In Exercises 15-19 find the general solution.
15. $y^{\prime \prime}-3 y^{\prime}+2 y=e^{3 x}(1+x)$
16. $y^{\prime \prime}-6 y^{\prime}+8 y=e^{x}(11-6 x)$
17. $y^{\prime \prime}+6 y^{\prime}+9 y=e^{2 x}(3-5 x)$
18. $y^{\prime \prime}+2 y^{\prime}-3 y=-16 x e^{x}$
19. $y^{\prime \prime}-2 y^{\prime}+y=e^{x}(2-12 x)$

In Exercises 20-23 solve the initial value problem and plot the solution.
20. $\mathrm{C} / \mathrm{G} y^{\prime \prime}-4 y^{\prime}-5 y=9 e^{2 x}(1+x), \quad y(0)=0, \quad y^{\prime}(0)=-10$
21. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+3 y^{\prime}-4 y=e^{2 x}(7+6 x), \quad y(0)=2, \quad y^{\prime}(0)=8$
22. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y^{\prime}+3 y=-e^{-x}(2+8 x), \quad y(0)=1, \quad y^{\prime}(0)=2$
23. $\mathrm{C} / \mathrm{G} \quad y^{\prime \prime}-3 y^{\prime}-10 y=7 e^{-2 x}, \quad y(0)=1, \quad y^{\prime}(0)=-17$

In Exercises 24-29 use the principle of superposition to find a particular solution.
24. $y^{\prime \prime}+y^{\prime}+y=x e^{x}+e^{-x}(1+2 x)$
25. $y^{\prime \prime}-7 y^{\prime}+12 y=-e^{x}(17-42 x)-e^{3 x}$
26. $y^{\prime \prime}-8 y^{\prime}+16 y=6 x e^{4 x}+2+16 x+16 x^{2}$
27. $y^{\prime \prime}-3 y^{\prime}+2 y=-e^{2 x}(3+4 x)-e^{x}$
28. $y^{\prime \prime}-2 y^{\prime}+2 y=e^{x}(1+x)+e^{-x}\left(2-8 x+5 x^{2}\right)$
29. $y^{\prime \prime}+y=e^{-x}\left(2-4 x+2 x^{2}\right)+e^{3 x}\left(8-12 x-10 x^{2}\right)$
30. (a) Prove that $y$ is a solution of the constant coefficient equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=e^{\alpha x} G(x) \tag{A}
\end{equation*}
$$

if and only if $y=u e^{\alpha x}$, where $u$ satisfies

$$
\begin{equation*}
a u^{\prime \prime}+p^{\prime}(\alpha) u^{\prime}+p(\alpha) u=G(x) \tag{B}
\end{equation*}
$$

and $p(r)=a r^{2}+b r+c$ is the characteristic polynomial of the complementary equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=0
$$

For the rest of this exercise, let $G$ be a polynomial. Give the requested proofs for the case where

$$
G(x)=g_{0}+g_{1} x+g_{2} x^{2}+g_{3} x^{3}
$$

(b) Prove that if $e^{\alpha x}$ isn't a solution of the complementary equation then (B) has a particular solution of the form $u_{p}=A(x)$, where $A$ is a polynomial of the same degree as $G$, as in Example 5.4.4. Conclude that (A) has a particular solution of the form $y_{p}=e^{\alpha x} A(x)$.
(c) Show that if $e^{\alpha x}$ is a solution of the complementary equation and $x e^{\alpha x}$ isn't, then (B) has a particular solution of the form $u_{p}=x A(x)$, where $A$ is a polynomial of the same degree as $G$, as in Example 5.4.5. Conclude that (A) has a particular solution of the form $y_{p}=x e^{\alpha x} A(x)$.
(d) Show that if $e^{\alpha x}$ and $x e^{\alpha x}$ are both solutions of the complementary equation then (B) has a particular solution of the form $u_{p}=x^{2} A(x)$, where $A$ is a polynomial of the same degree as $G$, and $x^{2} A(x)$ can be obtained by integrating $G / a$ twice, taking the constants of integration to be zero, as in Example 5.4.6. Conclude that (A) has a particular solution of the form $y_{p}=x^{2} e^{\alpha x} A(x)$.

Exercises 31-36 treat the equations considered in Examples 5.4.1-5.4.6. Substitute the suggested form of $y_{p}$ into the equation and equate the resulting coefficients of like functions on the two sides of the resulting equation to derive a set of simultaneous equations for the coefficients in $y_{p}$. Then solve for the coefficients to obtain $y_{p}$. Compare the work you've done with the work required to obtain the same results in Examples 5.4.1-5.4.6.
31. Compare with Example 5.4.1:

$$
y^{\prime \prime}-7 y^{\prime}+12 y=4 e^{2 x} ; \quad y_{p}=A e^{2 x}
$$

32. Compare with Example 5.4.2:

$$
y^{\prime \prime}-7 y^{\prime}+12 y=5 e^{4 x} ; \quad y_{p}=A x e^{4 x}
$$

33. Compare with Example 5.4.3.

$$
y^{\prime \prime}-8 y^{\prime}+16 y=2 e^{4 x} ; \quad y_{p}=A x^{2} e^{4 x}
$$

34. Compare with Example 5.4.4:

$$
y^{\prime \prime}-3 y^{\prime}+2 y=e^{3 x}\left(-1+2 x+x^{2}\right), \quad y_{p}=e^{3 x}\left(A+B x+C x^{2}\right)
$$

35. Compare with Example 5.4.5:

$$
y^{\prime \prime}-4 y^{\prime}+3 y=e^{3 x}\left(6+8 x+12 x^{2}\right), \quad y_{p}=e^{3 x}\left(A x+B x^{2}+C x^{3}\right)
$$

36. Compare with Example 5.4.6:

$$
4 y^{\prime \prime}+4 y^{\prime}+y=e^{-x / 2}\left(-8+48 x+144 x^{2}\right), \quad y_{p}=e^{-x / 2}\left(A x^{2}+B x^{3}+C x^{4}\right)
$$

37. Write $y=u e^{\alpha x}$ to find the general solution.
(a) $y^{\prime \prime}+2 y^{\prime}+y=\frac{e^{-x}}{\sqrt{x}}$
(b) $y^{\prime \prime}+6 y^{\prime}+9 y=e^{-3 x} \ln x$
(c) $y^{\prime \prime}-4 y^{\prime}+4 y=\frac{e^{2 x}}{1+x}$
(d) $4 y^{\prime \prime}+4 y^{\prime}+y=4 e^{-x / 2}\left(\frac{1}{x}+x\right)$
38. Suppose $\alpha \neq 0$ and $k$ is a positive integer. In most calculus books integrals like $\int x^{k} e^{\alpha x} d x$ are evaluated by integrating by parts $k$ times. This exercise presents another method. Let

$$
y=\int e^{\alpha x} P(x) d x
$$

with

$$
P(x)=p_{0}+p_{1} x+\cdots+p_{k} x^{k}, \quad\left(\text { where } p_{k} \neq 0\right)
$$

(a) Show that $y=e^{\alpha x} u$, where

$$
\begin{equation*}
u^{\prime}+\alpha u=P(x) \tag{A}
\end{equation*}
$$

(b) Show that (A) has a particular solution of the form

$$
u_{p}=A_{0}+A_{1} x+\cdots+A_{k} x^{k}
$$

where $A_{k}, A_{k-1}, \ldots, A_{0}$ can be computed successively by equating coefficients of $x^{k}, x^{k-1}, \ldots, 1$ on both sides of the equation

$$
u_{p}^{\prime}+\alpha u_{p}=P(x)
$$

(c) Conclude that

$$
\int e^{\alpha x} P(x) d x=\left(A_{0}+A_{1} x+\cdots+A_{k} x^{k}\right) e^{\alpha x}+c
$$

where $c$ is a constant of integration.
39. Use the method of Exercise 38 to evaluate the integral.
(a) $\int e^{x}(4+x) d x$
(b) $\int e^{-x}\left(-1+x^{2}\right) d x$
(c) $\int x^{3} e^{-2 x} d x$
(d) $\int e^{x}(1+x)^{2} d x$
(e) $\int e^{3 x}\left(-14+30 x+27 x^{2}\right) d x$
(f) $\int e^{-x}\left(1+6 x^{2}-14 x^{3}+3 x^{4}\right) d x$
40. Use the method suggested in Exercise 38 to evaluate $\int x^{k} e^{\alpha x} d x$, where $k$ is an arbitrary positive integer and $\alpha \neq 0$.

### 5.5 THE METHOD OF UNDETERMINED COEFFICIENTS II

In this section we consider the constant coefficient equation

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=e^{\lambda x}(P(x) \cos \omega x+Q(x) \sin \omega x) \tag{5.5.1}
\end{equation*}
$$

where $\lambda$ and $\omega$ are real numbers, $\omega \neq 0$, and $P$ and $Q$ are polynomials. We want to find a particular solution of (5.5.1). As in Section 5.4, the procedure that we will use is called the method of undetermined coefficients.
Forcing Functions Without Exponential Factors
We begin with the case where $\lambda=0$ in (5.5.1); thus, we we want to find a particular solution of

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=P(x) \cos \omega x+Q(x) \sin \omega x \tag{5.5.2}
\end{equation*}
$$

where $P$ and $Q$ are polynomials.
Differentiating $x^{r} \cos \omega x$ and $x^{r} \sin \omega x$ yields
and $\quad \frac{d}{d x} x^{r} \sin \omega x=\omega x^{r} \cos \omega x+r x^{r-1} \sin \omega x$.
This implies that if

$$
y_{p}=A(x) \cos \omega x+B(x) \sin \omega x
$$

where $A$ and $B$ are polynomials, then

$$
a y_{p}^{\prime \prime}+b y_{p}^{\prime}+c y_{p}=F(x) \cos \omega x+G(x) \sin \omega x
$$

where $F$ and $G$ are polynomials with coefficients that can be expressed in terms of the coefficients of $A$ and $B$. This suggests that we try to choose $A$ and $B$ so that $F=P$ and $G=Q$, respectively. Then $y_{p}$ will be a particular solution of (5.5.2). The next theorem tells us how to choose the proper form for $y_{p}$. For the proof see Exercise 37.

Theorem 5.5.1 Suppose $\omega$ is a positive number and $P$ and $Q$ are polynomials. Let $k$ be the larger of the degrees of $P$ and $Q$. Then the equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=P(x) \cos \omega x+Q(x) \sin \omega x
$$

has a particular solution

$$
\begin{equation*}
y_{p}=A(x) \cos \omega x+B(x) \sin \omega x, \tag{5.5.3}
\end{equation*}
$$

where

$$
A(x)=A_{0}+A_{1} x+\cdots+A_{k} x^{k} \quad \text { and } \quad B(x)=B_{0}+B_{1} x+\cdots+B_{k} x^{k}
$$

provided that $\cos \omega x$ and $\sin \omega x$ are not solutions of the complementary equation. The solutions of

$$
a\left(y^{\prime \prime}+\omega^{2} y\right)=P(x) \cos \omega x+Q(x) \sin \omega x
$$

(for which $\cos \omega x$ and $\sin \omega x$ are solutions of the complementary equation) are of the form (5.5.3), where

$$
A(x)=A_{0} x+A_{1} x^{2}+\cdots+A_{k} x^{k+1} \quad \text { and } \quad B(x)=B_{0} x+B_{1} x^{2}+\cdots+B_{k} x^{k+1} .
$$

For an analog of this theorem that's applicable to (5.5.1), see Exercise 38.
Example 5.5.1 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+y=5 \cos 2 x+10 \sin 2 x . \tag{5.5.4}
\end{equation*}
$$

Solution In (5.5.4) the coefficients of $\cos 2 x$ and $\sin 2 x$ are both zero degree polynomials (constants). Therefore Theorem 5.5.1 implies that (5.5.4) has a particular solution

$$
y_{p}=A \cos 2 x+B \sin 2 x .
$$

Since

$$
y_{p}^{\prime}=-2 A \sin 2 x+2 B \cos 2 x \quad \text { and } \quad y_{p}^{\prime \prime}=-4(A \cos 2 x+B \sin 2 x),
$$

replacing $y$ by $y_{p}$ in (5.5.4) yields

$$
\begin{aligned}
y_{p}^{\prime \prime}-2 y_{p}^{\prime}+y_{p}= & -4(A \cos 2 x+B \sin 2 x)-4(-A \sin 2 x+B \cos 2 x) \\
& +(A \cos 2 x+B \sin 2 x) \\
= & (-3 A-4 B) \cos 2 x+(4 A-3 B) \sin 2 x
\end{aligned}
$$

Equating the coefficients of $\cos 2 x$ and $\sin 2 x$ here with the corresponding coefficients on the right side of (5.5.4) shows that $y_{p}$ is a solution of (5.5.4) if

$$
\begin{aligned}
-3 A-4 B & =5 \\
4 A-3 B & =10 .
\end{aligned}
$$

Solving these equations yields $A=1, B=-2$. Therefore

$$
y_{p}=\cos 2 x-2 \sin 2 x
$$

is a particular solution of (5.5.4).
Example 5.5.2 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+4 y=8 \cos 2 x+12 \sin 2 x . \tag{5.5.5}
\end{equation*}
$$

Solution The procedure used in Example 5.5.1 doesn't work here; substituting $y_{p}=A \cos 2 x+B \sin 2 x$ for $y$ in (5.5.5) yields

$$
y_{p}^{\prime \prime}+4 y_{p}=-4(A \cos 2 x+B \sin 2 x)+4(A \cos 2 x+B \sin 2 x)=0
$$

for any choice of $A$ and $B$, since $\cos 2 x$ and $\sin 2 x$ are both solutions of the complementary equation for (5.5.5). We're dealing with the second case mentioned in Theorem 5.5.1, and should therefore try a particular solution of the form

$$
\begin{equation*}
y_{p}=x(A \cos 2 x+B \sin 2 x) . \tag{5.5.6}
\end{equation*}
$$

Then

$$
\begin{aligned}
y_{p}^{\prime} & =A \cos 2 x+B \sin 2 x+2 x(-A \sin 2 x+B \cos 2 x) \\
y_{p}^{\prime \prime} & =-4 A \sin 2 x+4 B \cos 2 x-4 x(A \cos 2 x+B \sin 2 x) \\
& =-4 A \sin 2 x+4 B \cos 2 x-4 y_{p}(\operatorname{see}(5.5 .6))
\end{aligned}
$$

and
so

$$
y_{p}^{\prime \prime}+4 y_{p}=-4 A \sin 2 x+4 B \cos 2 x
$$

Therefore $y_{p}$ is a solution of (5.5.5) if

$$
-4 A \sin 2 x+4 B \cos 2 x=8 \cos 2 x+12 \sin 2 x
$$

which holds if $A=-3$ and $B=2$. Therefore

$$
y_{p}=-x(3 \cos 2 x-2 \sin 2 x)
$$

is a particular solution of (5.5.5).
Example 5.5.3 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+3 y^{\prime}+2 y=(16+20 x) \cos x+10 \sin x \tag{5.5.7}
\end{equation*}
$$

Solution The coefficients of $\cos x$ and $\sin x$ in (5.5.7) are polynomials of degree one and zero, respectively. Therefore Theorem 5.5.1 tells us to look for a particular solution of (5.5.7) of the form

$$
\begin{equation*}
y_{p}=\left(A_{0}+A_{1} x\right) \cos x+\left(B_{0}+B_{1} x\right) \sin x \tag{5.5.8}
\end{equation*}
$$

Then

$$
\begin{equation*}
y_{p}^{\prime}=\left(A_{1}+B_{0}+B_{1} x\right) \cos x+\left(B_{1}-A_{0}-A_{1} x\right) \sin x \tag{5.5.9}
\end{equation*}
$$

and

$$
\begin{equation*}
y_{p}^{\prime \prime}=\left(2 B_{1}-A_{0}-A_{1} x\right) \cos x-\left(2 A_{1}+B_{0}+B_{1} x\right) \sin x \tag{5.5.10}
\end{equation*}
$$

so

$$
\begin{align*}
y_{p}^{\prime \prime}+3 y_{p}^{\prime}+2 y_{p}= & {\left[A_{0}+3 A_{1}+3 B_{0}+2 B_{1}+\left(A_{1}+3 B_{1}\right) x\right] \cos x }  \tag{5.5.11}\\
& +\left[B_{0}+3 B_{1}-3 A_{0}-2 A_{1}+\left(B_{1}-3 A_{1}\right) x\right] \sin x
\end{align*}
$$

Comparing the coefficients of $x \cos x, x \sin x, \cos x$, and $\sin x$ here with the corresponding coefficients in (5.5.7) shows that $y_{p}$ is a solution of (5.5.7) if

$$
\begin{aligned}
A_{1}+3 B_{1} & =20 \\
-3 A_{1}+B_{1} & =0 \\
A_{0}+3 B_{0}+3 A_{1}+2 B_{1} & =16 \\
-3 A_{0}+B_{0}-2 A_{1}+3 B_{1} & =10 .
\end{aligned}
$$

Solving the first two equations yields $A_{1}=2, B_{1}=6$. Substituting these into the last two equations yields

$$
\begin{aligned}
A_{0}+3 B_{0} & =16-3 A_{1}-2 B_{1}=-2 \\
-3 A_{0}+B_{0} & =10+2 A_{1}-3 B_{1}=-4
\end{aligned}
$$

Solving these equations yields $A_{0}=1, B_{0}=-1$. Substituting $A_{0}=1, A_{1}=2, B_{0}=-1, B_{1}=6$ into (5.5.8) shows that

$$
y_{p}=(1+2 x) \cos x-(1-6 x) \sin x
$$

is a particular solution of (5.5.7).
A Useful Observation
In (5.5.9), (5.5.10), and (5.5.11) the polynomials multiplying $\sin x$ can be obtained by replacing $A_{0}, A_{1}, B_{0}$, and $B_{1}$ by $B_{0}, B_{1},-A_{0}$, and $-A_{1}$, respectively, in the polynomials mutiplying $\cos x$. An analogous result applies in general, as follows (Exercise 36).

## Theorem 5.5.2 If

$$
y_{p}=A(x) \cos \omega x+B(x) \sin \omega x
$$

where $A(x)$ and $B(x)$ are polynomials with coefficients $A_{0} \ldots, A_{k}$ and $B_{0}, \ldots, B_{k}$, then the polynomials multiplying $\sin \omega x$ in

$$
y_{p}^{\prime}, \quad y_{p}^{\prime \prime}, \quad a y_{p}^{\prime \prime}+b y_{p}^{\prime}+c y_{p} \quad \text { and } \quad y_{p}^{\prime \prime}+\omega^{2} y_{p}
$$

can be obtained by replacing $A_{0}, \ldots, A_{k}$ by $B_{0}, \ldots, B_{k}$ and $B_{0}, \ldots, B_{k}$ by $-A_{0}, \ldots,-A_{k}$ in the corresponding polynomials multiplying $\cos \omega x$.

We won't use this theorem in our examples, but we recommend that you use it to check your manipulations when you work the exercises.

Example 5.5.4 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+y=(8-4 x) \cos x-(8+8 x) \sin x \text {. } \tag{5.5.12}
\end{equation*}
$$

Solution According to Theorem 5.5.1, we should look for a particular solution of the form

$$
\begin{equation*}
y_{p}=\left(A_{0} x+A_{1} x^{2}\right) \cos x+\left(B_{0} x+B_{1} x^{2}\right) \sin x \tag{5.5.13}
\end{equation*}
$$

since $\cos x$ and $\sin x$ are solutions of the complementary equation. However, let's try

$$
\begin{equation*}
y_{p}=\left(A_{0}+A_{1} x\right) \cos x+\left(B_{0}+B_{1} x\right) \sin x \tag{5.5.14}
\end{equation*}
$$

first, so you can see why it doesn't work. From (5.5.10),

$$
y_{p}^{\prime \prime}=\left(2 B_{1}-A_{0}-A_{1} x\right) \cos x-\left(2 A_{1}+B_{0}+B_{1} x\right) \sin x
$$

which together with (5.5.14) implies that

$$
y_{p}^{\prime \prime}+y_{p}=2 B_{1} \cos x-2 A_{1} \sin x .
$$

Since the right side of this equation does not contain $x \cos x$ or $x \sin x$, (5.5.14) can't satisfy (5.5.12) no matter how we choose $A_{0}, A_{1}, B_{0}$, and $B_{1}$.

Now let $y_{p}$ be as in (5.5.13). Then

$$
\begin{aligned}
y_{p}^{\prime}= & {\left[A_{0}+\left(2 A_{1}+B_{0}\right) x+B_{1} x^{2}\right] \cos x } \\
& +\left[B_{0}+\left(2 B_{1}-A_{0}\right) x-A_{1} x^{2}\right] \sin x \\
\text { and } \quad y_{p}^{\prime \prime}= & {\left[2 A_{1}+2 B_{0}-\left(A_{0}-4 B_{1}\right) x-A_{1} x^{2}\right] \cos x } \\
& +\left[2 B_{1}-2 A_{0}-\left(B_{0}+4 A_{1}\right) x-B_{1} x^{2}\right] \sin x,
\end{aligned}
$$

so

$$
y_{p}^{\prime \prime}+y_{p}=\left(2 A_{1}+2 B_{0}+4 B_{1} x\right) \cos x+\left(2 B_{1}-2 A_{0}-4 A_{1} x\right) \sin x .
$$

Comparing the coefficients of $\cos x$ and $\sin x$ here with the corresponding coefficients in (5.5.12) shows that $y_{p}$ is a solution of (5.5.12) if

$$
\begin{array}{rlr}
4 B_{1} & =-4 \\
-4 A_{1} & =-8 \\
2 B_{0}+2 A_{1} & =8 \\
-2 A_{0}+2 B_{1} & = & -8 .
\end{array}
$$

The solution of this system is $A_{1}=2, B_{1}=-1, A_{0}=3, B_{0}=2$. Therefore

$$
y_{p}=x[(3+2 x) \cos x+(2-x) \sin x]
$$

is a particular solution of (5.5.12).
Forcing Functions with Exponential Factors
To find a particular solution of

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=e^{\lambda x}(P(x) \cos \omega x+Q(x) \sin \omega x) \tag{5.5.15}
\end{equation*}
$$

when $\lambda \neq 0$, we recall from Section 5.4 that substituting $y=u e^{\lambda x}$ into (5.5.15) will produce a constant coefficient equation for $u$ with the forcing function $P(x) \cos \omega x+Q(x) \sin \omega x$. We can find a particular solution $u_{p}$ of this equation by the procedure that we used in Examples 5.5.1-5.5.4. Then $y_{p}=u_{p} e^{\lambda x}$ is a particular solution of (5.5.15).

Example 5.5.5 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}-3 y^{\prime}+2 y=e^{-2 x}[2 \cos 3 x-(34-150 x) \sin 3 x] \tag{5.5.16}
\end{equation*}
$$

Solution Let $y=u e^{-2 x}$. Then

$$
\begin{aligned}
y^{\prime \prime}-3 y^{\prime}+2 y & =e^{-2 x}\left[\left(u^{\prime \prime}-4 u^{\prime}+4 u\right)-3\left(u^{\prime}-2 u\right)+2 u\right] \\
& =e^{-2 x}\left(u^{\prime \prime}-7 u^{\prime}+12 u\right) \\
& =e^{-2 x}[2 \cos 3 x-(34-150 x) \sin 3 x]
\end{aligned}
$$

if

$$
\begin{equation*}
u^{\prime \prime}-7 u^{\prime}+12 u=2 \cos 3 x-(34-150 x) \sin 3 x . \tag{5.5.17}
\end{equation*}
$$

Since $\cos 3 x$ and $\sin 3 x$ aren't solutions of the complementary equation

$$
u^{\prime \prime}-7 u^{\prime}+12 u=0
$$

Theorem 5.5.1 tells us to look for a particular solution of (5.5.17) of the form

$$
\begin{equation*}
u_{p}=\left(A_{0}+A_{1} x\right) \cos 3 x+\left(B_{0}+B_{1} x\right) \sin 3 x \tag{5.5.18}
\end{equation*}
$$

Then
and

$$
\begin{aligned}
u_{p}^{\prime} & =\left(A_{1}+3 B_{0}+3 B_{1} x\right) \cos 3 x+\left(B_{1}-3 A_{0}-3 A_{1} x\right) \sin 3 x \\
u_{p}^{\prime \prime} & =\left(-9 A_{0}+6 B_{1}-9 A_{1} x\right) \cos 3 x-\left(9 B_{0}+6 A_{1}+9 B_{1} x\right) \sin 3 x
\end{aligned}
$$

so

$$
\begin{aligned}
u_{p}^{\prime \prime}-7 u_{p}^{\prime}+12 u_{p}= & {\left[3 A_{0}-21 B_{0}-7 A_{1}+6 B_{1}+\left(3 A_{1}-21 B_{1}\right) x\right] \cos 3 x } \\
& +\left[21 A_{0}+3 B_{0}-6 A_{1}-7 B_{1}+\left(21 A_{1}+3 B_{1}\right) x\right] \sin 3 x
\end{aligned}
$$

Comparing the coefficients of $x \cos 3 x, x \sin 3 x, \cos 3 x$, and $\sin 3 x$ here with the corresponding coefficients on the right side of (5.5.17) shows that $u_{p}$ is a solution of (5.5.17) if

$$
\begin{array}{rlr}
3 A_{1}-21 B_{1} & = & 0 \\
21 A_{1}+3 B_{1} & = & 150 \\
3 A_{0}-21 B_{0}-7 A_{1}+6 B_{1} & = & 2  \tag{5.5.19}\\
21 A_{0}+3 B_{0}-6 A_{1}-7 B_{1} & = & -34
\end{array}
$$

Solving the first two equations yields $A_{1}=7, B_{1}=1$. Substituting these values into the last two equations of (5.5.19) yields

$$
\begin{array}{rrr}
3 A_{0}-21 B_{0} & = & 2+7 A_{1}-6 B_{1}=45 \\
21 A_{0}+3 B_{0} & = & -34+6 A_{1}+7 B_{1}=15
\end{array}
$$

Solving this system yields $A_{0}=1, B_{0}=-2$. Substituting $A_{0}=1, A_{1}=7, B_{0}=-2$, and $B_{1}=1$ into (5.5.18) shows that

$$
u_{p}=(1+7 x) \cos 3 x-(2-x) \sin 3 x
$$

is a particular solution of (5.5.17). Therefore

$$
y_{p}=e^{-2 x}[(1+7 x) \cos 3 x-(2-x) \sin 3 x]
$$

is a particular solution of (5.5.16).
Example 5.5.6 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+2 y^{\prime}+5 y=e^{-x}[(6-16 x) \cos 2 x-(8+8 x) \sin 2 x] \tag{5.5.20}
\end{equation*}
$$

Solution Let $y=u e^{-x}$. Then

$$
\begin{aligned}
y^{\prime \prime}+2 y^{\prime}+5 y & =e^{-x}\left[\left(u^{\prime \prime}-2 u^{\prime}+u\right)+2\left(u^{\prime}-u\right)+5 u\right] \\
& =e^{-x}\left(u^{\prime \prime}+4 u\right) \\
& =e^{-x}[(6-16 x) \cos 2 x-(8+8 x) \sin 2 x]
\end{aligned}
$$

if

$$
\begin{equation*}
u^{\prime \prime}+4 u=(6-16 x) \cos 2 x-(8+8 x) \sin 2 x . \tag{5.5.21}
\end{equation*}
$$

Since $\cos 2 x$ and $\sin 2 x$ are solutions of the complementary equation

$$
u^{\prime \prime}+4 u=0
$$

Theorem 5.5.1 tells us to look for a particular solution of (5.5.21) of the form

$$
u_{p}=\left(A_{0} x+A_{1} x^{2}\right) \cos 2 x+\left(B_{0} x+B_{1} x^{2}\right) \sin 2 x .
$$

Then

$$
\begin{aligned}
u_{p}^{\prime}= & {\left[A_{0}+\left(2 A_{1}+2 B_{0}\right) x+2 B_{1} x^{2}\right] \cos 2 x } \\
& +\left[B_{0}+\left(2 B_{1}-2 A_{0}\right) x-2 A_{1} x^{2}\right] \sin 2 x \\
\text { and } \quad u_{p}^{\prime \prime}= & {\left[2 A_{1}+4 B_{0}-\left(4 A_{0}-8 B_{1}\right) x-4 A_{1} x^{2}\right] \cos 2 x } \\
& +\left[2 B_{1}-4 A_{0}-\left(4 B_{0}+8 A_{1}\right) x-4 B_{1} x^{2}\right] \sin 2 x,
\end{aligned}
$$

so

$$
u_{p}^{\prime \prime}+4 u_{p}=\left(2 A_{1}+4 B_{0}+8 B_{1} x\right) \cos 2 x+\left(2 B_{1}-4 A_{0}-8 A_{1} x\right) \sin 2 x
$$

Equating the coefficients of $x \cos 2 x, x \sin 2 x, \cos 2 x$, and $\sin 2 x$ here with the corresponding coefficients on the right side of (5.5.21) shows that $u_{p}$ is a solution of (5.5.21) if

$$
\begin{align*}
8 B_{1} & = & -16 \\
-8 A_{1} & = & -8  \tag{5.5.22}\\
4 B_{0}+2 A_{1} & = & 6 \\
-4 A_{0}+2 B_{1} & = & -8 .
\end{align*}
$$

The solution of this system is $A_{1}=1, B_{1}=-2, B_{0}=1, A_{0}=1$. Therefore

$$
u_{p}=x[(1+x) \cos 2 x+(1-2 x) \sin 2 x]
$$

is a particular solution of (5.5.21), and

$$
y_{p}=x e^{-x}[(1+x) \cos 2 x+(1-2 x) \sin 2 x]
$$

is a particular solution of (5.5.20).
You can also find a particular solution of (5.5.20) by substituting

$$
y_{p}=x e^{-x}\left[\left(A_{0}+A_{1} x\right) \cos 2 x+\left(B_{0}+B_{1} x\right) \sin 2 x\right]
$$

for $y$ in (5.5.20) and equating the coefficients of $x e^{-x} \cos 2 x, x e^{-x} \sin 2 x, e^{-x} \cos 2 x$, and $e^{-x} \sin 2 x$ in the resulting expression for

$$
y_{p}^{\prime \prime}+2 y_{p}^{\prime}+5 y_{p}
$$

with the corresponding coefficients on the right side of (5.5.20). (See Exercise 38). This leads to the same system (5.5.22) of equations for $A_{0}, A_{1}, B_{0}$, and $B_{1}$ that we obtained in Example 5.5.6. However, if you try this approach you'll see that deriving (5.5.22) this way is much more tedious than the way we did it in Example 5.5.6.

### 5.5 Exercises

In Exercises 1-17 find a particular solution.

1. $y^{\prime \prime}+3 y^{\prime}+2 y=7 \cos x-\sin x$
2. $y^{\prime \prime}+3 y^{\prime}+y=(2-6 x) \cos x-9 \sin x$
3. $y^{\prime \prime}+2 y^{\prime}+y=e^{x}(6 \cos x+17 \sin x)$
4. $y^{\prime \prime}+3 y^{\prime}-2 y=-e^{2 x}(5 \cos 2 x+9 \sin 2 x)$
5. $y^{\prime \prime}-y^{\prime}+y=e^{x}(2+x) \sin x$
6. $y^{\prime \prime}+3 y^{\prime}-2 y=e^{-2 x}[(4+20 x) \cos 3 x+(26-32 x) \sin 3 x]$
7. $y^{\prime \prime}+4 y=-12 \cos 2 x-4 \sin 2 x$
8. $y^{\prime \prime}+y=(-4+8 x) \cos x+(8-4 x) \sin x$
9. $4 y^{\prime \prime}+y=-4 \cos x / 2-8 x \sin x / 2$
10. $y^{\prime \prime}+2 y^{\prime}+2 y=e^{-x}(8 \cos x-6 \sin x)$
11. $y^{\prime \prime}-2 y^{\prime}+5 y=e^{x}[(6+8 x) \cos 2 x+(6-8 x) \sin 2 x]$
12. $y^{\prime \prime}+2 y^{\prime}+y=8 x^{2} \cos x-4 x \sin x$
13. $y^{\prime \prime}+3 y^{\prime}+2 y=\left(12+20 x+10 x^{2}\right) \cos x+8 x \sin x$
14. $y^{\prime \prime}+3 y^{\prime}+2 y=\left(1-x-4 x^{2}\right) \cos 2 x-\left(1+7 x+2 x^{2}\right) \sin 2 x$
15. $y^{\prime \prime}-5 y^{\prime}+6 y=-e^{x}\left[\left(4+6 x-x^{2}\right) \cos x-\left(2-4 x+3 x^{2}\right) \sin x\right]$
16. $y^{\prime \prime}-2 y^{\prime}+y=-e^{x}\left[\left(3+4 x-x^{2}\right) \cos x+\left(3-4 x-x^{2}\right) \sin x\right]$
17. $y^{\prime \prime}-2 y^{\prime}+2 y=e^{x}\left[\left(2-2 x-6 x^{2}\right) \cos x+\left(2-10 x+6 x^{2}\right) \sin x\right]$

In Exercises 1-17 find a particular solution and graph it.
18. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+2 y^{\prime}+y=e^{-x}[(5-2 x) \cos x-(3+3 x) \sin x]$
19. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+9 y=-6 \cos 3 x-12 \sin 3 x$
20. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+3 y^{\prime}+2 y=\left(1-x-4 x^{2}\right) \cos 2 x-\left(1+7 x+2 x^{2}\right) \sin 2 x$
21. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y^{\prime}+3 y=e^{-x}\left[\left(2+x+x^{2}\right) \cos x+\left(5+4 x+2 x^{2}\right) \sin x\right]$

In Exercises 22-26 solve the initial value problem.
22. $y^{\prime \prime}-7 y^{\prime}+6 y=-e^{x}(17 \cos x-7 \sin x), \quad y(0)=4, y^{\prime}(0)=2$
23. $y^{\prime \prime}-2 y^{\prime}+2 y=-e^{x}(6 \cos x+4 \sin x), \quad y(0)=1, y^{\prime}(0)=4$
24. $y^{\prime \prime}+6 y^{\prime}+10 y=-40 e^{x} \sin x, \quad y(0)=2, \quad y^{\prime}(0)=-3$
25. $y^{\prime \prime}-6 y^{\prime}+10 y=-e^{3 x}(6 \cos x+4 \sin x), \quad y(0)=2, \quad y^{\prime}(0)=7$
26. $y^{\prime \prime}-3 y^{\prime}+2 y=e^{3 x}[21 \cos x-(11+10 x) \sin x], y(0)=0, \quad y^{\prime}(0)=6$

In Exercises 27-32 use the principle of superposition to find a particular solution. Where indicated, solve the initial value problem.
27. $y^{\prime \prime}-2 y^{\prime}-3 y=4 e^{3 x}+e^{x}(\cos x-2 \sin x)$
28. $y^{\prime \prime}+y=4 \cos x-2 \sin x+x e^{x}+e^{-x}$
29. $y^{\prime \prime}-3 y^{\prime}+2 y=x e^{x}+2 e^{2 x}+\sin x$
30. $y^{\prime \prime}-2 y^{\prime}+2 y=4 x e^{x} \cos x+x e^{-x}+1+x^{2}$
31. $y^{\prime \prime}-4 y^{\prime}+4 y=e^{2 x}(1+x)+e^{2 x}(\cos x-\sin x)+3 e^{3 x}+1+x$
32. $y^{\prime \prime}-4 y^{\prime}+4 y=6 e^{2 x}+25 \sin x, \quad y(0)=5, y^{\prime}(0)=3$

In Exercises 33-35 solve the initial value problem and graph the solution.
33. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y=-e^{-2 x}[(4-7 x) \cos x+(2-4 x) \sin x], y(0)=3, \quad y^{\prime}(0)=1$
34. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y^{\prime}+4 y=2 \cos 2 x+3 \sin 2 x+e^{-x}, \quad y(0)=-1, y^{\prime}(0)=2$
35. C/G $y^{\prime \prime}+4 y=e^{x}(11+15 x)+8 \cos 2 x-12 \sin 2 x, \quad y(0)=3, y^{\prime}(0)=5$
36. (a) Verify that if

$$
y_{p}=A(x) \cos \omega x+B(x) \sin \omega x
$$

where $A$ and $B$ are twice differentiable, then

$$
\begin{aligned}
y_{p}^{\prime} & =\left(A^{\prime}+\omega B\right) \cos \omega x+\left(B^{\prime}-\omega A\right) \sin \omega x \text { and } \\
y_{p}^{\prime \prime} & =\left(A^{\prime \prime}+2 \omega B^{\prime}-\omega^{2} A\right) \cos \omega x+\left(B^{\prime \prime}-2 \omega A^{\prime}-\omega^{2} B\right) \sin \omega x
\end{aligned}
$$

(b) Use the results of (a) to verify that

$$
\begin{aligned}
a y_{p}^{\prime \prime}+b y_{p}^{\prime}+c y_{p}= & {\left[\left(c-a \omega^{2}\right) A+b \omega B+2 a \omega B^{\prime}+b A^{\prime}+a A^{\prime \prime}\right] \cos \omega x+} \\
& {\left[-b \omega A+\left(c-a \omega^{2}\right) B-2 a \omega A^{\prime}+b B^{\prime}+a B^{\prime \prime}\right] \sin \omega x . }
\end{aligned}
$$

(c) Use the results of (a) to verify that

$$
y_{p}^{\prime \prime}+\omega^{2} y_{p}=\left(A^{\prime \prime}+2 \omega B^{\prime}\right) \cos \omega x+\left(B^{\prime \prime}-2 \omega A^{\prime}\right) \sin \omega x
$$

(d) Prove Theorem 5.5.2.
37. Let $a, b, c$, and $\omega$ be constants, with $a \neq 0$ and $\omega>0$, and let

$$
P(x)=p_{0}+p_{1} x+\cdots+p_{k} x^{k} \quad \text { and } \quad Q(x)=q_{0}+q_{1} x+\cdots+q_{k} x^{k}
$$

where at least one of the coefficients $p_{k}, q_{k}$ is nonzero, so $k$ is the larger of the degrees of $P$ and $Q$.
(a) Show that if $\cos \omega x$ and $\sin \omega x$ are not solutions of the complementary equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=0
$$

then there are polynomials

$$
\begin{equation*}
A(x)=A_{0}+A_{1} x+\cdots+A_{k} x^{k} \quad \text { and } \quad B(x)=B_{0}+B_{1} x+\cdots+B_{k} x^{k} \tag{A}
\end{equation*}
$$

such that

$$
\begin{aligned}
\left(c-a \omega^{2}\right) A+b \omega B+2 a \omega B^{\prime}+b A^{\prime}+a A^{\prime \prime} & =P \\
-b \omega A+\left(c-a \omega^{2}\right) B-2 a \omega A^{\prime}+b B^{\prime}+a B^{\prime \prime} & =Q
\end{aligned}
$$

where $\left(A_{k}, B_{k}\right),\left(A_{k-1}, B_{k-1}\right), \ldots,\left(A_{0}, B_{0}\right)$ can be computed successively by solving the systems

$$
\begin{aligned}
\left(c-a \omega^{2}\right) A_{k}+b \omega B_{k} & =p_{k} \\
-b \omega A_{k}+\left(c-a \omega^{2}\right) B_{k} & =q_{k}
\end{aligned}
$$

and, if $1 \leq r \leq k$,

$$
\begin{aligned}
\left(c-a \omega^{2}\right) A_{k-r}+b \omega B_{k-r} & =p_{k-r}+\cdots \\
-b \omega A_{k-r}+\left(c-a \omega^{2}\right) B_{k-r} & =q_{k-r}+\cdots
\end{aligned}
$$

where the terms indicated by ". . ." depend upon the previously computed coefficients with subscripts greater than $k-r$. Conclude from this and Exercise 36(b) that

$$
\begin{equation*}
y_{p}=A(x) \cos \omega x+B(x) \sin \omega x \tag{B}
\end{equation*}
$$

is a particular solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=P(x) \cos \omega x+Q(x) \sin \omega x .
$$

(b) Conclude from Exercise 36(c) that the equation

$$
\begin{equation*}
a\left(y^{\prime \prime}+\omega^{2} y\right)=P(x) \cos \omega x+Q(x) \sin \omega x \tag{C}
\end{equation*}
$$

does not have a solution of the form (B) with $A$ and $B$ as in (A). Then show that there are polynomials

$$
A(x)=A_{0} x+A_{1} x^{2}+\cdots+A_{k} x^{k+1} \quad \text { and } \quad B(x)=B_{0} x+B_{1} x^{2}+\cdots+B_{k} x^{k+1}
$$

such that

$$
\begin{aligned}
& a\left(A^{\prime \prime}+2 \omega B^{\prime}\right)=P \\
& a\left(B^{\prime \prime}-2 \omega A^{\prime}\right)=Q
\end{aligned}
$$

where the pairs $\left(A_{k}, B_{k}\right),\left(A_{k-1}, B_{k-1}\right), \ldots,\left(A_{0}, B_{0}\right)$ can be computed successively as follows:

$$
\begin{aligned}
A_{k} & =-\frac{q_{k}}{2 a \omega(k+1)} \\
B_{k} & =\frac{p_{k}}{2 a \omega(k+1)}
\end{aligned}
$$

and, if $k \geq 1$,

$$
\begin{aligned}
A_{k-j} & =-\frac{1}{2 \omega}\left[\frac{q_{k-j}}{a(k-j+1)}-(k-j+2) B_{k-j+1}\right] \\
B_{k-j} & =\frac{1}{2 \omega}\left[\frac{p_{k-j}}{a(k-j+1)}-(k-j+2) A_{k-j+1}\right]
\end{aligned}
$$

for $1 \leq j \leq k$. Conclude that (B) with this choice of the polynomials $A$ and $B$ is a particular solution of (C).
38. Show that Theorem 5.5.1 implies the next theorem: Suppose $w$ is a positive number and $P$ and $Q$ are polynomials. Let $k$ be the larger of the degrees of $P$ and $Q$. Then the equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=e^{\lambda x}(P(x) \cos \omega x+Q(x) \sin \omega x)
$$

has a particular solution

$$
\begin{equation*}
y_{p}=e^{\lambda x}(A(x) \cos \omega x+B(x) \sin \omega x) \tag{A}
\end{equation*}
$$

where

$$
A(x)=A_{0}+A_{1} x+\cdots+A_{k} x^{k} \quad \text { and } \quad B(x)=B_{0}+B_{1} x+\cdots+B_{k} x^{k}
$$

provided that $e^{\lambda x} \cos \omega x$ and $e^{\lambda x} \sin \omega x$ are not solutions of the complementary equation. The equation

$$
a\left[y^{\prime \prime}-2 \lambda y^{\prime}+\left(\lambda^{2}+\omega^{2}\right) y\right]=e^{\lambda x}(P(x) \cos \omega x+Q(x) \sin \omega x)
$$

(for which $e^{\lambda x} \cos \omega x$ and $e^{\lambda x} \sin \omega x$ are solutions of the complementary equation) has a particular solution of the form (A), where

$$
A(x)=A_{0} x+A_{1} x^{2}+\cdots+A_{k} x^{k+1} \quad \text { and } \quad B(x)=B_{0} x+B_{1} x^{2}+\cdots+B_{k} x^{k+1}
$$

39. This exercise presents a method for evaluating the integral

$$
y=\int e^{\lambda x}(P(x) \cos \omega x+Q(x) \sin \omega x) d x
$$

where $\omega \neq 0$ and

$$
P(x)=p_{0}+p_{1} x+\cdots+p_{k} x^{k}, \quad Q(x)=q_{0}+q_{1} x+\cdots+q_{k} x^{k}
$$

(a) Show that $y=e^{\lambda x} u$, where

$$
\begin{equation*}
u^{\prime}+\lambda u=P(x) \cos \omega x+Q(x) \sin \omega x \tag{A}
\end{equation*}
$$

(b) Show that (A) has a particular solution of the form

$$
u_{p}=A(x) \cos \omega x+B(x) \sin \omega x
$$

where

$$
A(x)=A_{0}+A_{1} x+\cdots+A_{k} x^{k}, \quad B(x)=B_{0}+B_{1} x+\cdots+B_{k} x^{k}
$$

and the pairs of coefficients $\left(A_{k}, B_{k}\right),\left(A_{k-1}, B_{k-1}\right), \ldots,\left(A_{0}, B_{0}\right)$ can be computed successively as the solutions of pairs of equations obtained by equating the coefficients of $x^{r} \cos \omega x$ and $x^{r} \sin \omega x$ for $r=k, k-1, \ldots, 0$.
(c) Conclude that

$$
\int e^{\lambda x}(P(x) \cos \omega x+Q(x) \sin \omega x) d x=e^{\lambda x}(A(x) \cos \omega x+B(x) \sin \omega x)+c
$$

where $c$ is a constant of integration.
40. Use the method of Exercise 39 to evaluate the integral.
(a) $\int x^{2} \cos x d x$
(b) $\int x^{2} e^{x} \cos x d x$
(c) $\int x e^{-x} \sin 2 x d x$
(d) $\int x^{2} e^{-x} \sin x d x$
(e) $\int x^{3} e^{x} \sin x d x$
(f) $\int e^{x}[x \cos x-(1+3 x) \sin x] d x$
(g) $\int e^{-x}\left[\left(1+x^{2}\right) \cos x+\left(1-x^{2}\right) \sin x\right] d x$

### 5.7 VARIATION OF PARAMETERS

In this section we give a method called variation of parameters for finding a particular solution of

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F(x) \tag{5.7.1}
\end{equation*}
$$

if we know a fundamental set $\left\{y_{1}, y_{2}\right\}$ of solutions of the complementary equation

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0 \tag{5.7.2}
\end{equation*}
$$

Having found a particular solution $y_{p}$ by this method, we can write the general solution of (5.7.1) as

$$
y=y_{p}+c_{1} y_{1}+c_{2} y_{2} .
$$

Since we need only one nontrivial solution of (5.7.2) to find the general solution of (5.7.1) by reduction of order, it's natural to ask why we're interested in variation of parameters, which requires two linearly independent solutions of (5.7.2) to achieve the same goal. Here's the answer:

- If we already know two linearly independent solutions of (5.7.2) then variation of parameters will probably be simpler than reduction of order.
- Variation of parameters generalizes naturally to a method for finding particular solutions of higher order linear equations (Section 9.4) and linear systems of equations (Section 10.7), while reduction of order doesn't.
- Variation of parameters is a powerful theoretical tool used by researchers in differential equations. Although a detailed discussion of this is beyond the scope of this book, you can get an idea of what it means from Exercises 37-39.

We'll now derive the method. As usual, we consider solutions of (5.7.1) and (5.7.2) on an interval $(a, b)$ where $P_{0}, P_{1}, P_{2}$, and $F$ are continuous and $P_{0}$ has no zeros. Suppose that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of the complementary equation (5.7.2). We look for a particular solution of (5.7.1) in the form

$$
\begin{equation*}
y_{p}=u_{1} y_{1}+u_{2} y_{2} \tag{5.7.3}
\end{equation*}
$$

where $u_{1}$ and $u_{2}$ are functions to be determined so that $y_{p}$ satisfies (5.7.1). You may not think this is a good idea, since there are now two unknown functions to be determined, rather than one. However, since $u_{1}$ and $u_{2}$ have to satisfy only one condition (that $y_{p}$ is a solution of (5.7.1)), we can impose a second condition that produces a convenient simplification, as follows.

Differentiating (5.7.3) yields

$$
\begin{equation*}
y_{p}^{\prime}=u_{1} y_{1}^{\prime}+u_{2} y_{2}^{\prime}+u_{1}^{\prime} y_{1}+u_{2}^{\prime} y_{2} \tag{5.7.4}
\end{equation*}
$$

As our second condition on $u_{1}$ and $u_{2}$ we require that

$$
\begin{equation*}
u_{1}^{\prime} y_{1}+u_{2}^{\prime} y_{2}=0 \tag{5.7.5}
\end{equation*}
$$

Then (5.7.4) becomes

$$
\begin{equation*}
y_{p}^{\prime}=u_{1} y_{1}^{\prime}+u_{2} y_{2}^{\prime} \tag{5.7.6}
\end{equation*}
$$

that is, (5.7.5) permits us to differentiate $y_{p}$ (once!) as if $u_{1}$ and $u_{2}$ are constants. Differentiating (5.7.4) yields

$$
\begin{equation*}
y_{p}^{\prime \prime}=u_{1} y_{1}^{\prime \prime}+u_{2} y_{2}^{\prime \prime}+u_{1}^{\prime} y_{1}^{\prime}+u_{2}^{\prime} y_{2}^{\prime} . \tag{5.7.7}
\end{equation*}
$$

(There are no terms involving $u_{1}^{\prime \prime}$ and $u_{2}^{\prime \prime}$ here, as there would be if we hadn't required (5.7.5).) Substituting (5.7.3), (5.7.6), and (5.7.7) into (5.7.1) and collecting the coefficients of $u_{1}$ and $u_{2}$ yields

$$
u_{1}\left(P_{0} y_{1}^{\prime \prime}+P_{1} y_{1}^{\prime}+P_{2} y_{1}\right)+u_{2}\left(P_{0} y_{2}^{\prime \prime}+P_{1} y_{2}^{\prime}+P_{2} y_{2}\right)+P_{0}\left(u_{1}^{\prime} y_{1}^{\prime}+u_{2}^{\prime} y_{2}^{\prime}\right)=F .
$$

As in the derivation of the method of reduction of order, the coefficients of $u_{1}$ and $u_{2}$ here are both zero because $y_{1}$ and $y_{2}$ satisfy the complementary equation. Hence, we can rewrite the last equation as

$$
\begin{equation*}
P_{0}\left(u_{1}^{\prime} y_{1}^{\prime}+u_{2}^{\prime} y_{2}^{\prime}\right)=F \tag{5.7.8}
\end{equation*}
$$

Therefore $y_{p}$ in (5.7.3) satisfies (5.7.1) if

$$
\begin{align*}
u_{1}^{\prime} y_{1}+u_{2}^{\prime} y_{2} & =0 \\
u_{1}^{\prime} y_{1}^{\prime}+u_{2}^{\prime} y_{2}^{\prime} & =\frac{F}{P_{0}} \tag{5.7.9}
\end{align*}
$$

where the first equation is the same as (5.7.5) and the second is from (5.7.8).
We'll now show that you can always solve (5.7.9) for $u_{1}^{\prime}$ and $u_{2}^{\prime}$. (The method that we use here will always work, but simpler methods usually work when you're dealing with specific equations.) To obtain $u_{1}^{\prime}$, multiply the first equation in (5.7.9) by $y_{2}^{\prime}$ and the second equation by $y_{2}$. This yields

$$
\begin{aligned}
u_{1}^{\prime} y_{1} y_{2}^{\prime}+u_{2}^{\prime} y_{2} y_{2}^{\prime} & =0 \\
u_{1}^{\prime} y_{1}^{\prime} y_{2}+u_{2}^{\prime} y_{2}^{\prime} y_{2} & =\frac{F y_{2}}{P_{0}} .
\end{aligned}
$$

Subtracting the second equation from the first yields

$$
\begin{equation*}
u_{1}^{\prime}\left(y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}\right)=-\frac{F y_{2}}{P_{0}} \tag{5.7.10}
\end{equation*}
$$

Since $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (5.7.2) on $(a, b)$, Theorem 5.1.6 implies that the Wronskian $y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}$ has no zeros on $(a, b)$. Therefore we can solve (5.7.10) for $u_{1}^{\prime}$, to obtain

$$
\begin{equation*}
u_{1}^{\prime}=-\frac{F y_{2}}{P_{0}\left(y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}\right)} . \tag{5.7.11}
\end{equation*}
$$

We leave it to you to start from (5.7.9) and show by a similar argument that

$$
\begin{equation*}
u_{2}^{\prime}=\frac{F y_{1}}{P_{0}\left(y_{1} y_{2}^{\prime}-y_{1}^{\prime} y_{2}\right)} . \tag{5.7.12}
\end{equation*}
$$

We can now obtain $u_{1}$ and $u_{2}$ by integrating $u_{1}^{\prime}$ and $u_{2}^{\prime}$. The constants of integration can be taken to be zero, since any choice of $u_{1}$ and $u_{2}$ in (5.7.3) will suffice.

You should not memorize (5.7.11) and (5.7.12). On the other hand, you don't want to rederive the whole procedure for every specific problem. We recommend the a compromise:
(a) Write

$$
\begin{equation*}
y_{p}=u_{1} y_{1}+u_{2} y_{2} \tag{5.7.13}
\end{equation*}
$$

to remind yourself of what you're doing.
(b) Write the system

$$
\begin{align*}
u_{1}^{\prime} y_{1}+u_{2}^{\prime} y_{2} & =0 \\
u_{1}^{\prime} y_{1}^{\prime}+u_{2}^{\prime} y_{2}^{\prime} & =\frac{F}{P_{0}} \tag{5.7.14}
\end{align*}
$$

for the specific problem you're trying to solve.
(c) Solve (5.7.14) for $u_{1}^{\prime}$ and $u_{2}^{\prime}$ by any convenient method.
(d) Obtain $u_{1}$ and $u_{2}$ by integrating $u_{1}^{\prime}$ and $u_{2}^{\prime}$, taking the constants of integration to be zero.
(e) Substitute $u_{1}$ and $u_{2}$ into (5.7.13) to obtain $y_{p}$.

Example 5.7.1 Find a particular solution $y_{p}$ of

$$
\begin{equation*}
x^{2} y^{\prime \prime}-2 x y^{\prime}+2 y=x^{9 / 2} \tag{5.7.15}
\end{equation*}
$$

given that $y_{1}=x$ and $y_{2}=x^{2}$ are solutions of the complementary equation

$$
x^{2} y^{\prime \prime}-2 x y^{\prime}+2 y=0
$$

Then find the general solution of (5.7.15).

Solution We set

$$
y_{p}=u_{1} x+u_{2} x^{2}
$$

where

$$
\begin{aligned}
u_{1}^{\prime} x+u_{2}^{\prime} x^{2} & =0 \\
u_{1}^{\prime}+2 u_{2}^{\prime} x & =\frac{x^{9 / 2}}{x^{2}}=x^{5 / 2}
\end{aligned}
$$

From the first equation, $u_{1}^{\prime}=-u_{2}^{\prime} x$. Substituting this into the second equation yields $u_{2}^{\prime} x=x^{5 / 2}$, so $u_{2}^{\prime}=x^{3 / 2}$ and therefore $u_{1}^{\prime}=-u_{2}^{\prime} x=-x^{5 / 2}$. Integrating and taking the constants of integration to be zero yields

$$
u_{1}=-\frac{2}{7} x^{7 / 2} \quad \text { and } \quad u_{2}=\frac{2}{5} x^{5 / 2}
$$

Therefore

$$
y_{p}=u_{1} x+u_{2} x^{2}=-\frac{2}{7} x^{7 / 2} x+\frac{2}{5} x^{5 / 2} x^{2}=\frac{4}{35} x^{9 / 2}
$$

and the general solution of (5.7.15) is

$$
y=\frac{4}{35} x^{9 / 2}+c_{1} x+c_{2} x^{2} .
$$

Example 5.7.2 Find a particular solution $y_{p}$ of

$$
\begin{equation*}
(x-1) y^{\prime \prime}-x y^{\prime}+y=(x-1)^{2}, \tag{5.7.16}
\end{equation*}
$$

given that $y_{1}=x$ and $y_{2}=e^{x}$ are solutions of the complementary equation

$$
(x-1) y^{\prime \prime}-x y^{\prime}+y=0
$$

Then find the general solution of (5.7.16).

Solution We set

$$
y_{p}=u_{1} x+u_{2} e^{x}
$$

where

$$
\begin{aligned}
u_{1}^{\prime} x+u_{2}^{\prime} e^{x} & =0 \\
u_{1}^{\prime}+u_{2}^{\prime} e^{x} & =\frac{(x-1)^{2}}{x-1}=x-1
\end{aligned}
$$

Subtracting the first equation from the second yields $-u_{1}^{\prime}(x-1)=x-1$, so $u_{1}^{\prime}=-1$. From this and the first equation, $u_{2}^{\prime}=-x e^{-x} u_{1}^{\prime}=x e^{-x}$. Integrating and taking the constants of integration to be zero yields

$$
u_{1}=-x \quad \text { and } \quad u_{2}=-(x+1) e^{-x} .
$$

Therefore

$$
y_{p}=u_{1} x+u_{2} e^{x}=(-x) x+\left(-(x+1) e^{-x}\right) e^{x}=-x^{2}-x-1
$$

so the general solution of (5.7.16) is

$$
\begin{equation*}
y=y_{p}+c_{1} x+c_{2} e^{x}=-x^{2}-x-1+c_{1} x+c_{2} e^{x}=-x^{2}-1+\left(c_{1}-1\right) x+c_{2} e^{x} . \tag{5.7.17}
\end{equation*}
$$

However, since $c_{1}$ is an arbitrary constant, so is $c_{1}-1$; therefore, we improve the appearance of this result by renaming the constant and writing the general solution as

$$
\begin{equation*}
y=-x^{2}-1+c_{1} x+c_{2} e^{x} \tag{5.7.18}
\end{equation*}
$$

There's nothing wrong with leaving the general solution of (5.7.16) in the form (5.7.17); however, we think you'll agree that (5.7.18) is preferable. We can also view the transition from (5.7.17) to (5.7.18) differently. In this example the particular solution $y_{p}=-x^{2}-x-1$ contained the term $-x$, which satisfies the complementary equation. We can drop this term and redefine $y_{p}=-x^{2}-1$, since $-x^{2}-x-1$ is a solution of (5.7.16) and $x$ is a solution of the complementary equation; hence, $-x^{2}-1=\left(-x^{2}-x-1\right)+x$ is also a solution of (5.7.16). In general, it's always legitimate to drop linear combinations of $\left\{y_{1}, y_{2}\right\}$ from particular solutions obtained by variation of parameters. (See Exercise 36 for a general discussion of this question.) We'll do this in the following examples and in the answers to exercises that ask for a particular solution. Therefore, don't be concerned if your answer to such an exercise differs from ours only by a solution of the complementary equation.

Example 5.7.3 Find a particular solution of

$$
\begin{equation*}
y^{\prime \prime}+3 y^{\prime}+2 y=\frac{1}{1+e^{x}} \tag{5.7.19}
\end{equation*}
$$

Then find the general solution.

## Solution

The characteristic polynomial of the complementary equation

$$
\begin{equation*}
y^{\prime \prime}+3 y^{\prime}+2 y=0 \tag{5.7.20}
\end{equation*}
$$

is $p(r)=r^{2}+3 r+2=(r+1)(r+2)$, so $y_{1}=e^{-x}$ and $y_{2}=e^{-2 x}$ form a fundamental set of solutions of (5.7.20). We look for a particular solution of (5.7.19) in the form

$$
y_{p}=u_{1} e^{-x}+u_{2} e^{-2 x}
$$

where

$$
\begin{aligned}
u_{1}^{\prime} e^{-x}+u_{2}^{\prime} e^{-2 x} & =0 \\
-u_{1}^{\prime} e^{-x}-2 u_{2}^{\prime} e^{-2 x} & =\frac{1}{1+e^{x}}
\end{aligned}
$$

Adding these two equations yields

$$
-u_{2}^{\prime} e^{-2 x}=\frac{1}{1+e^{x}}, \quad \text { so } \quad u_{2}^{\prime}=-\frac{e^{2 x}}{1+e^{x}}
$$

From the first equation,

$$
u_{1}^{\prime}=-u_{2}^{\prime} e^{-x}=\frac{e^{x}}{1+e^{x}}
$$

Integrating by means of the substitution $v=e^{x}$ and taking the constants of integration to be zero yields

$$
u_{1}=\int \frac{e^{x}}{1+e^{x}} d x=\int \frac{d v}{1+v}=\ln (1+v)=\ln \left(1+e^{x}\right)
$$

and

$$
\begin{aligned}
u_{2} & =-\int \frac{e^{2 x}}{1+e^{x}} d x=-\int \frac{v}{1+v} d v=\int\left[\frac{1}{1+v}-1\right] d v \\
& =\ln (1+v)-v=\ln \left(1+e^{x}\right)-e^{x}
\end{aligned}
$$

Therefore

$$
\begin{aligned}
y_{p} & =u_{1} e^{-x}+u_{2} e^{-2 x} \\
& =\left[\ln \left(1+e^{x}\right)\right] e^{-x}+\left[\ln \left(1+e^{x}\right)-e^{x}\right] e^{-2 x}
\end{aligned}
$$

so

$$
y_{p}=\left(e^{-x}+e^{-2 x}\right) \ln \left(1+e^{x}\right)-e^{-x} .
$$

Since the last term on the right satisfies the complementary equation, we drop it and redefine

$$
y_{p}=\left(e^{-x}+e^{-2 x}\right) \ln \left(1+e^{x}\right) .
$$

The general solution of (5.7.19) is

$$
y=y_{p}+c_{1} e^{-x}+c_{2} e^{-2 x}=\left(e^{-x}+e^{-2 x}\right) \ln \left(1+e^{x}\right)+c_{1} e^{-x}+c_{2} e^{-2 x} .
$$

Example 5.7.4 Solve the initial value problem

$$
\begin{equation*}
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=\frac{2}{x+1}, \quad y(0)=-1, \quad y^{\prime}(0)=-5, \tag{5.7.21}
\end{equation*}
$$

given that

$$
y_{1}=\frac{1}{x-1} \quad \text { and } \quad y_{2}=\frac{1}{x+1}
$$

are solutions of the complementary equation

$$
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0 .
$$

Solution We first use variation of parameters to find a particular solution of

$$
\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=\frac{2}{x+1}
$$

on $(-1,1)$ in the form

$$
y_{p}=\frac{u_{1}}{x-1}+\frac{u_{2}}{x+1},
$$

where

$$
\begin{align*}
\frac{u_{1}^{\prime}}{x-1}+\frac{u_{2}^{\prime}}{x+1} & =0  \tag{5.7.22}\\
-\frac{u_{1}^{\prime}}{(x-1)^{2}}-\frac{u_{2}^{\prime}}{(x+1)^{2}} & =\frac{2}{(x+1)\left(x^{2}-1\right)}
\end{align*}
$$

Multiplying the first equation by $1 /(x-1)$ and adding the result to the second equation yields

$$
\begin{equation*}
\left[\frac{1}{x^{2}-1}-\frac{1}{(x+1)^{2}}\right] u_{2}^{\prime}=\frac{2}{(x+1)\left(x^{2}-1\right)} . \tag{5.7.23}
\end{equation*}
$$

Since

$$
\left[\frac{1}{x^{2}-1}-\frac{1}{(x+1)^{2}}\right]=\frac{(x+1)-(x-1)}{(x+1)\left(x^{2}-1\right)}=\frac{2}{(x+1)\left(x^{2}-1\right)},
$$

(5.7.23) implies that $u_{2}^{\prime}=1$. From (5.7.22),

$$
u_{1}^{\prime}=-\frac{x-1}{x+1} u_{2}^{\prime}=-\frac{x-1}{x+1} .
$$

Integrating and taking the constants of integration to be zero yields

$$
\begin{aligned}
u_{1} & =-\int \frac{x-1}{x+1} d x=-\int \frac{x+1-2}{x+1} d x \\
& =\int\left[\frac{2}{x+1}-1\right] d x=2 \ln (x+1)-x
\end{aligned}
$$

and

$$
u_{2}=\int d x=x .
$$

Therefore

$$
\begin{aligned}
y_{p} & =\frac{u_{1}}{x-1}+\frac{u_{2}}{x+1}=[2 \ln (x+1)-x] \frac{1}{x-1}+x \frac{1}{x+1} \\
& =\frac{2 \ln (x+1)}{x-1}+x\left[\frac{1}{x+1}-\frac{1}{x-1}\right]=\frac{2 \ln (x+1)}{x-1}-\frac{2 x}{(x+1)(x-1)} .
\end{aligned}
$$

However, since

$$
\frac{2 x}{(x+1)(x-1)}=\left[\frac{1}{x+1}+\frac{1}{x-1}\right]
$$

is a solution of the complementary equation, we redefine

$$
y_{p}=\frac{2 \ln (x+1)}{x-1} .
$$

Therefore the general solution of (5.7.24) is

$$
\begin{equation*}
y=\frac{2 \ln (x+1)}{x-1}+\frac{c_{1}}{x-1}+\frac{c_{2}}{x+1} . \tag{5.7.24}
\end{equation*}
$$

Differentiating this yields

$$
y^{\prime}=\frac{2}{x^{2}-1}-\frac{2 \ln (x+1)}{(x-1)^{2}}-\frac{c_{1}}{(x-1)^{2}}-\frac{c_{2}}{(x+1)^{2}} .
$$

Setting $x=0$ in the last two equations and imposing the initial conditions $y(0)=-1$ and $y^{\prime}(0)=-5$ yields the system

$$
\begin{aligned}
-c_{1}+c_{2} & =-1 \\
-2-c_{1}-c_{2} & =-5 .
\end{aligned}
$$

The solution of this system is $c_{1}=2, c_{2}=1$. Substituting these into (5.7.24) yields

$$
\begin{aligned}
y & =\frac{2 \ln (x+1)}{x-1}+\frac{2}{x-1}+\frac{1}{x+1} \\
& =\frac{2 \ln (x+1)}{x-1}+\frac{3 x+1}{x^{2}-1}
\end{aligned}
$$

as the solution of (5.7.21). Figure 5.7.1 is a graph of the solution.
Comparison of Methods
We've now considered three methods for solving nonhomogeneous linear equations: undetermined coefficients, reduction of order, and variation of parameters. It's natural to ask which method is best for a given problem. The method of undetermined coefficients should be used for constant coefficient equations with forcing functions that are linear combinations of polynomials multiplied by functions of the form $e^{\alpha x}, e^{\lambda x} \cos \omega x$, or $e^{\lambda x} \sin \omega x$. Although the other two methods can be used to solve such problems, they will be more difficult except in the most trivial cases, because of the integrations involved.

If the equation isn't a constant coefficient equation or the forcing function isn't of the form just specified, the method of undetermined coefficients does not apply and the choice is necessarily between the other two methods. The case could be made that reduction of order is better because it requires only one solution of the complementary equation while variation of parameters requires two. However, variation of parameters will probably be easier if you already know a fundamental set of solutions of the complementary equation.


Figure 5.7.1 $y=\frac{2 \ln (x+1)}{x-1}+\frac{3 x+1}{x^{2}-1}$

### 5.7 Exercises

In Exercises 1-6 use variation of parameters to find a particular solution.

1. $y^{\prime \prime}+9 y=\tan 3 x$
2. $y^{\prime \prime}+4 y=\sin 2 x \sec ^{2} 2 x$
3. $y^{\prime \prime}-3 y^{\prime}+2 y=\frac{4}{1+e^{-x}}$
4. $y^{\prime \prime}-2 y^{\prime}+2 y=3 e^{x} \sec x$
5. $y^{\prime \prime}-2 y^{\prime}+y=14 x^{3 / 2} e^{x}$
6. $y^{\prime \prime}-y=\frac{4 e^{-x}}{1-e^{-2 x}}$

In Exercises 7-29 use variation of parameters to find a particular solution, given the solutions $y_{1}, y_{2}$ of the complementary equation.
7. $x^{2} y^{\prime \prime}+x y^{\prime}-y=2 x^{2}+2 ; \quad y_{1}=x, \quad y_{2}=\frac{1}{x}$
8. $x y^{\prime \prime}+(2-2 x) y^{\prime}+(x-2) y=e^{2 x} ; \quad y_{1}=e^{x}, \quad y_{2}=\frac{e^{x}}{x}$
9. $4 x^{2} y^{\prime \prime}+\left(4 x-8 x^{2}\right) y^{\prime}+\left(4 x^{2}-4 x-1\right) y=4 x^{1 / 2} e^{x}, \quad x>0$;
$y_{1}=x^{1 / 2} e^{x}, y_{2}=x^{-1 / 2} e^{x}$
10. $y^{\prime \prime}+4 x y^{\prime}+\left(4 x^{2}+2\right) y=4 e^{-x(x+2)} ; \quad y_{1}=e^{-x^{2}}, \quad y_{2}=x e^{-x^{2}}$
11. $x^{2} y^{\prime \prime}-4 x y^{\prime}+6 y=x^{5 / 2}, x>0 ; \quad y_{1}=x^{2}, y_{2}=x^{3}$
12. $x^{2} y^{\prime \prime}-3 x y^{\prime}+3 y=2 x^{4} \sin x ; \quad y_{1}=x, y_{2}=x^{3}$
13. $(2 x+1) y^{\prime \prime}-2 y^{\prime}-(2 x+3) y=(2 x+1)^{2} e^{-x} ; \quad y_{1}=e^{-x}, \quad y_{2}=x e^{x}$
14. $4 x y^{\prime \prime}+2 y^{\prime}+y=\sin \sqrt{x} ; \quad y_{1}=\cos \sqrt{x}, \quad y_{2}=\sin \sqrt{x}$
15. $x y^{\prime \prime}-(2 x+2) y^{\prime}+(x+2) y=6 x^{3} e^{x} ; \quad y_{1}=e^{x}, \quad y_{2}=x^{3} e^{x}$
16. $x^{2} y^{\prime \prime}-(2 a-1) x y^{\prime}+a^{2} y=x^{a+1} ; \quad y_{1}=x^{a}, \quad y_{2}=x^{a} \ln x$
17. $x^{2} y^{\prime \prime}-2 x y^{\prime}+\left(x^{2}+2\right) y=x^{3} \cos x ; \quad y_{1}=x \cos x, \quad y_{2}=x \sin x$
18. $x y^{\prime \prime}-y^{\prime}-4 x^{3} y=8 x^{5} ; \quad y_{1}=e^{x^{2}}, y_{2}=e^{-x^{2}}$
19. $(\sin x) y^{\prime \prime}+(2 \sin x-\cos x) y^{\prime}+(\sin x-\cos x) y=e^{-x} ; \quad y_{1}=e^{-x}, \quad y_{2}=e^{-x} \cos x$
20. $4 x^{2} y^{\prime \prime}-4 x y^{\prime}+\left(3-16 x^{2}\right) y=8 x^{5 / 2} ; \quad y_{1}=\sqrt{x} e^{2 x}, y_{2}=\sqrt{x} e^{-2 x}$
21. $4 x^{2} y^{\prime \prime}-4 x y^{\prime}+\left(4 x^{2}+3\right) y=x^{7 / 2} ; \quad y_{1}=\sqrt{x} \sin x, y_{2}=\sqrt{x} \cos x$
22. $x^{2} y^{\prime \prime}-2 x y^{\prime}-\left(x^{2}-2\right) y=3 x^{4} ; \quad y_{1}=x e^{x}, y_{2}=x e^{-x}$
23. $x^{2} y^{\prime \prime}-2 x(x+1) y^{\prime}+\left(x^{2}+2 x+2\right) y=x^{3} e^{x} ; \quad y_{1}=x e^{x}, \quad y_{2}=x^{2} e^{x}$
24. $\quad x^{2} y^{\prime \prime}-x y^{\prime}-3 y=x^{3 / 2} ; \quad y_{1}=1 / x, \quad y_{2}=x^{3}$
25. $\quad x^{2} y^{\prime \prime}-x(x+4) y^{\prime}+2(x+3) y=x^{4} e^{x} ; \quad y_{1}=x^{2}, \quad y_{2}=x^{2} e^{x}$
26. $\quad x^{2} y^{\prime \prime}-2 x(x+2) y^{\prime}+\left(x^{2}+4 x+6\right) y=2 x e^{x} ; \quad y_{1}=x^{2} e^{x}, \quad y_{2}=x^{3} e^{x}$
27. $x^{2} y^{\prime \prime}-4 x y^{\prime}+\left(x^{2}+6\right) y=x^{4} ; \quad y_{1}=x^{2} \cos x, \quad y_{2}=x^{2} \sin x$
28. $(x-1) y^{\prime \prime}-x y^{\prime}+y=2(x-1)^{2} e^{x} ; \quad y_{1}=x, \quad y_{2}=e^{x}$
29. $4 x^{2} y^{\prime \prime}-4 x(x+1) y^{\prime}+(2 x+3) y=x^{5 / 2} e^{x} ; \quad y_{1}=\sqrt{x}, \quad y_{2}=\sqrt{x} e^{x}$

In Exercises 30-32 use variation of parameters to solve the initial value problem, given $y_{1}, y_{2}$ are solutions of the complementary equation.
30. $(3 x-1) y^{\prime \prime}-(3 x+2) y^{\prime}-(6 x-8) y=(3 x-1)^{2} e^{2 x}, \quad y(0)=1, y^{\prime}(0)=2$;
$y_{1}=e^{2 x}, y_{2}=x e^{-x}$
31. $(x-1)^{2} y^{\prime \prime}-2(x-1) y^{\prime}+2 y=(x-1)^{2}, \quad y(0)=3, \quad y^{\prime}(0)=-6$; $y_{1}=x-1, y_{2}=x^{2}-1$
32. $(x-1)^{2} y^{\prime \prime}-\left(x^{2}-1\right) y^{\prime}+(x+1) y=(x-1)^{3} e^{x}, \quad y(0)=4, \quad y^{\prime}(0)=-6$; $y_{1}=(x-1) e^{x}, \quad y_{2}=x-1$

In Exercises 33-35 use variation of parameters to solve the initial value problem and graph the solution, given that $y_{1}, y_{2}$ are solutions of the complementary equation.
33. $\mathrm{C} / \mathrm{G}\left(x^{2}-1\right) y^{\prime \prime}+4 x y^{\prime}+2 y=2 x, \quad y(0)=0, y^{\prime}(0)=-2 ; \quad y_{1}=\frac{1}{x-1}, y_{2}=\frac{1}{x+1}$
34. $\mathrm{C} / \mathrm{G} x^{2} y^{\prime \prime}+2 x y^{\prime}-2 y=-2 x^{2}, \quad y(1)=1, y^{\prime}(1)=-1 ; \quad y_{1}=x, y_{2}=\frac{1}{x^{2}}$
35. $\mathrm{C} / \mathrm{G}(x+1)(2 x+3) y^{\prime \prime}+2(x+2) y^{\prime}-2 y=(2 x+3)^{2}, \quad y(0)=0, \quad y^{\prime}(0)=0$; $y_{1}=x+2, \quad y_{2}=\frac{1}{x+1}$
36. Suppose

$$
y_{p}=\bar{y}+a_{1} y_{1}+a_{2} y_{2}
$$

is a particular solution of

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=F(x), \tag{A}
\end{equation*}
$$

where $y_{1}$ and $y_{2}$ are solutions of the complementary equation

$$
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0
$$

Show that $\bar{y}$ is also a solution of (A).
37. Suppose $p, q$, and $f$ are continuous on $(a, b)$ and let $x_{0}$ be in $(a, b)$. Let $y_{1}$ and $y_{2}$ be the solutions of

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=0
$$

such that

$$
y_{1}\left(x_{0}\right)=1, \quad y_{1}^{\prime}\left(x_{0}\right)=0, \quad y_{2}\left(x_{0}\right)=0, \quad y_{2}^{\prime}\left(x_{0}\right)=1 .
$$

Use variation of parameters to show that the solution of the initial value problem

$$
y^{\prime \prime}+p(x) y^{\prime}+q(x) y=f(x), \quad y\left(x_{0}\right)=k_{0}, y^{\prime}\left(x_{0}\right)=k_{1}
$$

is

$$
\begin{aligned}
y(x)= & k_{0} y_{1}(x)+k_{1} y_{2}(x) \\
& \quad+\int_{x_{0}}^{x}\left(y_{1}(t) y_{2}(x)-y_{1}(x) y_{2}(t)\right) f(t) \exp \left(\int_{x_{0}}^{t} p(s) d s\right) d t
\end{aligned}
$$

Hint: Use Abel's formula for the Wronskian of $\left\{y_{1}, y_{2}\right\}$, and integrate $u_{1}^{\prime}$ and $u_{2}^{\prime}$ from $x_{0}$ to $x$. Show also that

$$
\begin{aligned}
y^{\prime}(x)= & k_{0} y_{1}^{\prime}(x)+k_{1} y_{2}^{\prime}(x) \\
& \quad+\int_{x_{0}}^{x}\left(y_{1}(t) y_{2}^{\prime}(x)-y_{1}^{\prime}(x) y_{2}(t)\right) f(t) \exp \left(\int_{x_{0}}^{t} p(s) d s\right) d t
\end{aligned}
$$

38. Suppose $f$ is continuous on an open interval that contains $x_{0}=0$. Use variation of parameters to find a formula for the solution of the initial value problem

$$
y^{\prime \prime}-y=f(x), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} .
$$

39. Suppose $f$ is continuous on $(a, \infty)$, where $a<0$, so $x_{0}=0$ is in $(a, \infty)$.
(a) Use variation of parameters to find a formula for the solution of the initial value problem

$$
y^{\prime \prime}+y=f(x), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} .
$$

Hint: You will need the addition formulas for the sine and cosine:

$$
\begin{aligned}
\sin (A+B) & =\sin A \cos B+\cos A \sin B \\
\cos (A+B) & =\cos A \cos B-\sin A \sin B
\end{aligned}
$$

For the rest of this exercise assume that the improper integral $\int_{0}^{\infty} f(t) d t$ is absolutely convergent.
(b) Show that if $y$ is a solution of

$$
\begin{equation*}
y^{\prime \prime}+y=f(x) \tag{A}
\end{equation*}
$$

on $(a, \infty)$, then

$$
\begin{equation*}
\lim _{x \rightarrow \infty}\left(y(x)-A_{0} \cos x-A_{1} \sin x\right)=0 \tag{B}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{x \rightarrow \infty}\left(y^{\prime}(x)+A_{0} \sin x-A_{1} \cos x\right)=0 \tag{C}
\end{equation*}
$$

where

$$
A_{0}=k_{0}-\int_{0}^{\infty} f(t) \sin t d t \quad \text { and } \quad A_{1}=k_{1}+\int_{0}^{\infty} f(t) \cos t d t
$$

HINT: Recall from calculus that if $\int_{0}^{\infty} f(t) d t$ converges absolutely, then $\lim _{x \rightarrow \infty} \int_{x}^{\infty}|f(t)| d t=0$.
(c) Show that if $A_{0}$ and $A_{1}$ are arbitrary constants, then there's a unique solution of $y^{\prime \prime}+y=$ $f(x)$ on $(a, \infty)$ that satisfies (B) and (C).

# CHAPTER 6 <br> Applications of Linear Second Order Equations 

[^0]
### 6.1 SPRING PROBLEMS I

We consider the motion of an object of mass $m$, suspended from a spring of negligible mass. We say that the spring-mass system is in equilibrium when the object is at rest and the forces acting on it sum to zero. The position of the object in this case is the equilibrium position. We define $y$ to be the displacement of the object from its equilibrium position (Figure 6.1.1), measured positive upward.


Figure 6.1.1 (a) $y>0$ (b) $y=0$, (c) $y<0$
Figure 6.1.2 A spring - mass system with damping
Our model accounts for the following kinds of forces acting on the object:

- The force $-m g$, due to gravity.
- A force $F_{s}$ exerted by the spring resisting change in its length. The natural length of the spring is its length with no mass attached. We assume that the spring obeys Hooke's law: If the length of the spring is changed by an amount $\Delta L$ from its natural length, then the spring exerts a force $F_{s}=k \Delta L$, where $k$ is a positive number called the spring constant. If the spring is stretched then $\Delta L>0$ and $F_{s}>0$, so the spring force is upward, while if the spring is compressed then $\Delta L<0$ and $F_{s}<0$, so the spring force is downward.
- A damping force $F_{d}=-c y^{\prime}$ that resists the motion with a force proportional to the velocity of the object. It may be due to air resistance or friction in the spring. However, a convenient way to visualize a damping force is to assume that the object is rigidly attached to a piston with negligible mass immersed in a cylinder (called a dashpot) filled with a viscous liquid (Figure 6.1.2). As the piston moves, the liquid exerts a damping force. We say that the motion is undamped if $c=0$, or damped if $c>0$.
- An external force $F$, other than the force due to gravity, that may vary with $t$, but is independent of displacement and velocity. We say that the motion is free if $F \equiv 0$, or forced if $F \not \equiv 0$.

From Newton's second law of motion,

$$
\begin{equation*}
m y^{\prime \prime}=-m g+F_{d}+F_{s}+F=-m g-c y^{\prime}+F_{s}+F . \tag{6.1.1}
\end{equation*}
$$

We must now relate $F_{s}$ to $y$. In the absence of external forces the object stretches the spring by an amount $\Delta l$ to assume its equilibrium position (Figure 6.1.3). Since the sum of the forces acting on the object is then zero, Hooke's Law implies that $m g=k \Delta l$. If the object is displaced $y$ units from its equilibrium position, the total change in the length of the spring is $\Delta L=\Delta l-y$, so Hooke's law implies that

$$
F_{s}=k \Delta L=k \Delta l-k y .
$$

Substituting this into (6.1.1) yields

$$
m y^{\prime \prime}=-m g-c y^{\prime}+k \Delta l-k y+F
$$



Figure 6.1.3 (a) Natural length of spring (b) Spring stretched by mass

Since $m g=k \Delta l$ this can be written as

$$
\begin{equation*}
m y^{\prime \prime}+c y^{\prime}+k y=F . \tag{6.1.2}
\end{equation*}
$$

We call this the equation of motion.
Simple Harmonic Motion
Throughout the rest of this section we'll consider spring-mass systems without damping; that is, $c=0$. We'll consider systems with damping in the next section.

We first consider the case where the motion is also free; that is, $F=0$. We begin with an example.
Example 6.1.1 An object stretches a spring 6 inches in equilibrium.
(a) Set up the equation of motion and find its general solution.
(b) Find the displacement of the object for $t>0$ if it's initially displaced 18 inches above equilibrium and given a downward velocity of $3 \mathrm{ft} / \mathrm{s}$.
$\underline{\text { SOLUTION(a) }}$ Setting $c=0$ and $F=0$ in (6.1.2) yields the equation of motion

$$
m y^{\prime \prime}+k y=0
$$

which we rewrite as

$$
\begin{equation*}
y^{\prime \prime}+\frac{k}{m} y=0 . \tag{6.1.3}
\end{equation*}
$$

Although we would need the weight of the object to obtain $k$ from the equation $m g=k \Delta l$ we can obtain $k / m$ from $\Delta l$ alone; thus, $k / m=g / \Delta l$. Consistent with the units used in the problem statement, we take $g=32 \mathrm{ft} / \mathrm{s}^{2}$. Although $\Delta l$ is stated in inches, we must convert it to feet to be consistent with this choice of $g$; that is, $\Delta l=1 / 2 \mathrm{ft}$. Therefore

$$
\frac{k}{m}=\frac{32}{1 / 2}=64
$$

and (6.1.3) becomes

$$
\begin{equation*}
y^{\prime \prime}+64 y=0 . \tag{6.1.4}
\end{equation*}
$$

The characteristic equation of (6.1.4) is

$$
r^{2}+64=0
$$



Figure 6.1.4 $y=\frac{3}{2} \cos 8 t-\frac{3}{8} \sin 8 t$
which has the zeros $r= \pm 8 i$. Therefore the general solution of (6.1.4) is

$$
\begin{equation*}
y=c_{1} \cos 8 t+c_{2} \sin 8 t \tag{6.1.5}
\end{equation*}
$$

SOLUTION(b) The initial upward displacement of 18 inches is positive and must be expressed in feet. The initial downward velocity is negative; thus,

$$
y(0)=\frac{3}{2} \quad \text { and } \quad y^{\prime}(0)=-3 .
$$

Differentiating (6.1.5) yields

$$
\begin{equation*}
y^{\prime}=-8 c_{1} \sin 8 t+8 c_{2} \cos 8 t \tag{6.1.6}
\end{equation*}
$$

Setting $t=0$ in (6.1.5) and (6.1.6) and imposing the initial conditions shows that $c_{1}=3 / 2$ and $c_{2}=$ $-3 / 8$. Therefore

$$
y=\frac{3}{2} \cos 8 t-\frac{3}{8} \sin 8 t,
$$

where $y$ is in feet (Figure 6.1.4).
We'll now consider the equation

$$
m y^{\prime \prime}+k y=0
$$

where $m$ and $k$ are arbitrary positive numbers. Dividing through by $m$ and defining $\omega_{0}=\sqrt{k / m}$ yields

$$
y^{\prime \prime}+\omega_{0}^{2} y=0 .
$$

The general solution of this equation is

$$
\begin{equation*}
y=c_{1} \cos \omega_{0} t+c_{2} \sin \omega_{0} t . \tag{6.1.7}
\end{equation*}
$$

We can rewrite this in a more useful form by defining

$$
\begin{equation*}
R=\sqrt{c_{1}^{2}+c_{2}^{2}} \tag{6.1.8}
\end{equation*}
$$

and

$$
\begin{equation*}
c_{1}=R \cos \phi \quad \text { and } \quad c_{2}=R \sin \phi \tag{6.1.9}
\end{equation*}
$$



Figure 6.1.5 $R=\sqrt{c_{1}^{2}+c_{2}^{2}} ; \quad c_{1}=R \cos \phi ; \quad c_{2}=R \sin \phi$

Substituting from (6.1.9) into (6.1.7) and applying the identity

$$
\cos \omega_{0} t \cos \phi+\sin \omega_{0} t \sin \phi=\cos \left(\omega_{0} t-\phi\right)
$$

yields

$$
\begin{equation*}
y=R \cos \left(\omega_{0} t-\phi\right) \tag{6.1.10}
\end{equation*}
$$

From (6.1.8) and (6.1.9) we see that the $R$ and $\phi$ can be interpreted as polar coordinates of the point with rectangular coordinates $\left(c_{1}, c_{2}\right)$ (Figure 6.1.5). Given $c_{1}$ and $c_{2}$, we can compute $R$ from (6.1.8). From (6.1.8) and (6.1.9), we see that $\phi$ is related to $c_{1}$ and $c_{2}$ by

$$
\cos \phi=\frac{c_{1}}{\sqrt{c_{1}^{2}+c_{2}^{2}}} \quad \text { and } \quad \sin \phi=\frac{c_{2}}{\sqrt{c_{1}^{2}+c_{2}^{2}}}
$$

There are infinitely many angles $\phi$, differing by integer multiples of $2 \pi$, that satisfy these equations. We will always choose $\phi$ so that $-\pi \leq \phi<\pi$.

The motion described by (6.1.7) or (6.1.10) is simple harmonic motion. We see from either of these equations that the motion is periodic, with period

$$
T=2 \pi / \omega_{0}
$$

This is the time required for the object to complete one full cycle of oscillation (for example, to move from its highest position to its lowest position and back to its highest position). Since the highest and lowest positions of the object are $y=R$ and $y=-R$, we say that $R$ is the amplitude of the oscillation. The angle $\phi$ in (6.1.10) is the phase angle. It's measured in radians. Equation (6.1.10) is the amplitude-phase form of the displacement. If $t$ is in seconds then $\omega_{0}$ is in radians per second ( $\mathrm{rad} / \mathrm{s}$ ); it's the frequency of the motion. It is also called the natural frequency of the spring-mass system without damping.

Example 6.1.2 We found the displacement of the object in Example 6.1.1 to be

$$
y=\frac{3}{2} \cos 8 t-\frac{3}{8} \sin 8 t .
$$

Find the frequency, period, amplitude, and phase angle of the motion.

Solution The frequency is $\omega_{0}=8 \mathrm{rad} / \mathrm{s}$, and the period is $T=2 \pi / \omega_{0}=\pi / 4 \mathrm{~s}$. Since $c_{1}=3 / 2$ and $c_{2}=-3 / 8$, the amplitude is

$$
R=\sqrt{c_{1}^{2}+c_{2}^{2}}=\sqrt{\left(\frac{3}{2}\right)^{2}+\left(\frac{3}{8}\right)^{2}}=\frac{3}{8} \sqrt{17}
$$

The phase angle is determined by

$$
\begin{equation*}
\cos \phi=\frac{\frac{3}{2}}{\frac{3}{8} \sqrt{17}}=\frac{4}{\sqrt{17}} \tag{6.1.11}
\end{equation*}
$$

and

$$
\begin{equation*}
\sin \phi=\frac{-\frac{3}{8}}{\frac{3}{8} \sqrt{17}}=-\frac{1}{\sqrt{17}} . \tag{6.1.12}
\end{equation*}
$$

Using a calculator, we see from (6.1.11) that

$$
\phi \approx \pm .245 \mathrm{rad}
$$

Since $\sin \phi<0$ (see (6.1.12)), the minus sign applies here; that is,

$$
\phi \approx-.245 \mathrm{rad}
$$

Example 6.1.3 The natural length of a spring is 1 m . An object is attached to it and the length of the spring increases to 102 cm when the object is in equilibrium. Then the object is initially displaced downward 1 cm and given an upward velocity of $14 \mathrm{~cm} / \mathrm{s}$. Find the displacement for $t>0$. Also, find the natural frequency, period, amplitude, and phase angle of the resulting motion. Express the answers in cgs units.

Solution In cgs units $g=980 \mathrm{~cm} / \mathrm{s}^{2}$. Since $\Delta l=2 \mathrm{~cm}, \omega_{0}^{2}=g / \Delta l=490$. Therefore

$$
y^{\prime \prime}+490 y=0, \quad y(0)=-1, \quad y^{\prime}(0)=14
$$

The general solution of the differential equation is

$$
y=c_{1} \cos 7 \sqrt{10} t+c_{2} \sin 7 \sqrt{10} t
$$

so

$$
y^{\prime}=7 \sqrt{10}\left(-c_{1} \sin 7 \sqrt{10} t+c_{2} \cos 7 \sqrt{10} t\right)
$$

Substituting the initial conditions into the last two equations yields $c_{1}=-1$ and $c_{2}=2 / \sqrt{10}$. Hence,

$$
y=-\cos 7 \sqrt{10} t+\frac{2}{\sqrt{10}} \sin 7 \sqrt{10} t
$$

The frequency is $7 \sqrt{10} \mathrm{rad} / \mathrm{s}$, and the period is $T=2 \pi /(7 \sqrt{10}) \mathrm{s}$. The amplitude is

$$
R=\sqrt{c_{1}^{2}+c_{2}^{2}}=\sqrt{(-1)^{2}+\left(\frac{2}{\sqrt{10}}\right)^{2}}=\sqrt{\frac{7}{5}} \mathrm{~cm}
$$

The phase angle is determined by

$$
\cos \phi=\frac{c_{1}}{R}=-\sqrt{\frac{5}{7}} \quad \text { and } \quad \sin \phi=\frac{c_{2}}{R}=\sqrt{\frac{2}{7}}
$$

Therefore $\phi$ is in the second quadrant and

$$
\phi=\cos ^{-1}\left(-\sqrt{\frac{5}{7}}\right) \approx 2.58 \mathrm{rad}
$$

## Undamped Forced Oscillation

In many mechanical problems a device is subjected to periodic external forces. For example, soldiers marching in cadence on a bridge cause periodic disturbances in the bridge, and the engines of a propeller driven aircraft cause periodic disturbances in its wings. In the absence of sufficient damping forces, such disturbances - even if small in magnitude - can cause structural breakdown if they are at certain critical frequencies. To illustrate, this we'll consider the motion of an object in a spring-mass system without damping, subject to an external force

$$
F(t)=F_{0} \cos \omega t
$$

where $F_{0}$ is a constant. In this case the equation of motion (6.1.2) is

$$
m y^{\prime \prime}+k y=F_{0} \cos \omega t
$$

which we rewrite as

$$
\begin{equation*}
y^{\prime \prime}+\omega_{0}^{2} y=\frac{F_{0}}{m} \cos \omega t \tag{6.1.13}
\end{equation*}
$$

with $\omega_{0}=\sqrt{k / m}$. We'll see from the next two examples that the solutions of (6.1.13) with $\omega \neq \omega_{0}$ behave very differently from the solutions with $\omega=\omega_{0}$.

Example 6.1.4 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+\omega_{0}^{2} y=\frac{F_{0}}{m} \cos \omega t, \quad y(0)=0, \quad y^{\prime}(0)=0 \tag{6.1.14}
\end{equation*}
$$

given that $\omega \neq \omega_{0}$.

Solution We first obtain a particular solution of (6.1.13) by the method of undetermined coefficients. Since $\omega \neq \omega_{0}$, $\cos \omega t$ isn't a solution of the complementary equation

$$
y^{\prime \prime}+\omega_{0}^{2} y=0
$$

Therefore (6.1.13) has a particular solution of the form

$$
y_{p}=A \cos \omega t+B \sin \omega t
$$

Since

$$
\begin{gathered}
y_{p}^{\prime \prime}=-\omega^{2}(A \cos \omega t+B \sin \omega t), \\
y_{p}^{\prime \prime}+\omega_{0}^{2} y_{p}=\frac{F_{0}}{m} \cos \omega t
\end{gathered}
$$

if and only if

$$
\left(\omega_{0}^{2}-\omega^{2}\right)(A \cos \omega t+B \sin \omega t)=\frac{F_{0}}{m} \cos \omega t .
$$

This holds if and only if

$$
A=\frac{F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)} \quad \text { and } \quad B=0
$$

so

$$
y_{p}=\frac{F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)} \cos \omega t .
$$

The general solution of (6.1.13) is

$$
\begin{equation*}
y=\frac{F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)} \cos \omega t+c_{1} \cos \omega_{0} t+c_{2} \sin \omega_{0} t \tag{6.1.15}
\end{equation*}
$$

so

$$
y^{\prime}=\frac{-\omega F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)} \sin \omega t+\omega_{0}\left(-c_{1} \sin \omega_{0} t+c_{2} \cos \omega_{0} t\right)
$$

The initial conditions $y(0)=0$ and $y^{\prime}(0)=0$ in (6.1.14) imply that

$$
c_{1}=-\frac{F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)} \quad \text { and } \quad c_{2}=0
$$

Substituting these into (6.1.15) yields

$$
\begin{equation*}
y=\frac{F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)}\left(\cos \omega t-\cos \omega_{0} t\right) \tag{6.1.16}
\end{equation*}
$$

It is revealing to write this in a different form. We start with the trigonometric identities

$$
\begin{aligned}
& \cos (\alpha-\beta)=\cos \alpha \cos \beta+\sin \alpha \sin \beta \\
& \cos (\alpha+\beta)=\cos \alpha \cos \beta-\sin \alpha \sin \beta
\end{aligned}
$$



Figure 6.1.6 Undamped oscillation with beats

Subtracting the second identity from the first yields

$$
\begin{equation*}
\cos (\alpha-\beta)-\cos (\alpha+\beta)=2 \sin \alpha \sin \beta \tag{6.1.17}
\end{equation*}
$$

Now let

$$
\begin{equation*}
\alpha-\beta=\omega t \quad \text { and } \quad \alpha+\beta=\omega_{0} t, \tag{6.1.18}
\end{equation*}
$$

so that

$$
\begin{equation*}
\alpha=\frac{\left(\omega_{0}+\omega\right) t}{2} \quad \text { and } \quad \beta=\frac{\left(\omega_{0}-\omega\right) t}{2} . \tag{6.1.19}
\end{equation*}
$$

Substituting (6.1.18) and (6.1.19) into (6.1.17) yields

$$
\cos \omega t-\cos \omega_{0} t=2 \sin \frac{\left(\omega_{0}-\omega\right) t}{2} \sin \frac{\left(\omega_{0}+\omega\right) t}{2}
$$

and substituting this into (6.1.16) yields

$$
\begin{equation*}
y=R(t) \sin \frac{\left(\omega_{0}+\omega\right) t}{2}, \tag{6.1.20}
\end{equation*}
$$

where

$$
\begin{equation*}
R(t)=\frac{2 F_{0}}{m\left(\omega_{0}^{2}-\omega^{2}\right)} \sin \frac{\left(\omega_{0}-\omega\right) t}{2} . \tag{6.1.21}
\end{equation*}
$$

From (6.1.20) we can regard $y$ as a sinusoidal variation with frequency $\left(\omega_{0}+\omega\right) / 2$ and variable amplitude $|R(t)|$. In Figure 6.1.6 the dashed curve above the $t$ axis is $y=|R(t)|$, the dashed curve below the $t$ axis is $y=-|R(t)|$, and the displacement $y$ appears as an oscillation bounded by them. The oscillation of $y$ for $t$ on an interval between successive zeros of $R(t)$ is called a beat.

You can see from (6.1.20) and (6.1.21) that

$$
|y(t)| \leq \frac{2\left|F_{0}\right|}{m\left|\omega_{0}^{2}-\omega^{2}\right|}
$$

moreover, if $\omega+\omega_{0}$ is sufficiently large compared with $\omega-\omega_{0}$, then $|y|$ assumes values close to (perhaps equal to) this upper bound during each beat. However, the oscillation remains bounded for all $t$. (This assumes that the spring can withstand deflections of this size and continue to obey Hooke's law.) The next example shows that this isn't so if $\omega=\omega_{0}$.

Example 6.1.5 Find the general solution of

$$
\begin{equation*}
y^{\prime \prime}+\omega_{0}^{2} y=\frac{F_{0}}{m} \cos \omega_{0} t \tag{6.1.22}
\end{equation*}
$$

Solution We first obtain a particular solution $y_{p}$ of (6.1.22). Since $\cos \omega_{0} t$ is a solution of the complementary equation, the form for $y_{p}$ is

$$
\begin{equation*}
y_{p}=t\left(A \cos \omega_{0} t+B \sin \omega_{0} t\right) . \tag{6.1.23}
\end{equation*}
$$

Then

$$
y_{p}^{\prime}=A \cos \omega_{0} t+B \sin \omega_{0} t+\omega_{0} t\left(-A \sin \omega_{0} t+B \cos \omega_{0} t\right)
$$

and

$$
\begin{equation*}
y_{p}^{\prime \prime}=2 \omega_{0}\left(-A \sin \omega_{0} t+B \cos \omega_{0} t\right)-\omega_{0}^{2} t\left(A \cos \omega_{0} t+B \sin \omega_{0} t\right) . \tag{6.1.24}
\end{equation*}
$$

From (6.1.23) and (6.1.24), we see that $y_{p}$ satisfies (6.1.22) if

$$
-2 A \omega_{0} \sin \omega_{0} t+2 B \omega_{0} \cos \omega_{0} t=\frac{F_{0}}{m} \cos \omega_{0} t
$$

that is, if

$$
A=0 \quad \text { and } \quad B=\frac{F_{0}}{2 m \omega_{0}} .
$$

Therefore

$$
y_{p}=\frac{F_{0} t}{2 m \omega_{0}} \sin \omega_{0} t
$$

is a particular solution of (6.1.22). The general solution of (6.1.22) is

$$
y=\frac{F_{0} t}{2 m \omega_{0}} \sin \omega_{0} t+c_{1} \cos \omega_{0} t+c_{2} \sin \omega_{0} t .
$$

The graph of $y_{p}$ is shown in Figure 6.1.7, where it can be seen that $y_{p}$ oscillates between the dashed lines

$$
y=\frac{F_{0} t}{2 m \omega_{0}} \quad \text { and } \quad y=-\frac{F_{0} t}{2 m \omega_{0}}
$$

with increasing amplitude that approaches $\infty$ as $t \rightarrow \infty$. Of course, this means that the spring must eventually fail to obey Hooke's law or break.

This phenomenon of unbounded displacements of a spring-mass system in response to a periodic forcing function at its natural frequency is called resonance. More complicated mechanical structures can also exhibit resonance-like phenomena. For example, rhythmic oscillations of a suspension bridge by wind forces or of an airplane wing by periodic vibrations of reciprocating engines can cause damage or even failure if the frequencies of the disturbances are close to critical frequencies determined by the parameters of the mechanical system in question.


Figure 6.1.7 Unbounded displacement due to resonance

### 6.1 Exercises

In the following exercises assume that there's no damping.

1. C/G An object stretches a spring 4 inches in equilibrium. Find and graph its displacement for $t>0$ if it's initially displaced 36 inches above equilibrium and given a downward velocity of 2 $\mathrm{ft} / \mathrm{s}$.
2. An object stretches a string 1.2 inches in equilibrium. Find its displacement for $t>0$ if it's initially displaced 3 inches below equilibrium and given a downward velocity of $2 \mathrm{ft} / \mathrm{s}$.
3. A spring with natural length .5 m has length 50.5 cm with a mass of 2 gm suspended from it. The mass is initially displaced 1.5 cm below equilibrium and released with zero velocity. Find its displacement for $t>0$.
4. An object stretches a spring 6 inches in equilibrium. Find its displacement for $t>0$ if it's initially displaced 3 inches above equilibrium and given a downward velocity of 6 inches/s. Find the frequency, period, amplitude and phase angle of the motion.
5. C/G An object stretches a spring 5 cm in equilibrium. It is initially displaced 10 cm above equilibrium and given an upward velocity of $.25 \mathrm{~m} / \mathrm{s}$. Find and graph its displacement for $t>0$. Find the frequency, period, amplitude, and phase angle of the motion.
6. A 10 kg mass stretches a spring 70 cm in equilibrium. Suppose a 2 kg mass is attached to the spring, initially displaced 25 cm below equilibrium, and given an upward velocity of $2 \mathrm{~m} / \mathrm{s}$. Find its displacement for $t>0$. Find the frequency, period, amplitude, and phase angle of the motion.
7. A weight stretches a spring 1.5 inches in equilibrium. The weight is initially displaced 8 inches above equilibrium and given a downward velocity of $4 \mathrm{ft} / \mathrm{s}$. Find its displacement for $t>0$.
8. A weight stretches a spring 6 inches in equilibrium. The weight is initially displaced 6 inches above equilibrium and given a downward velocity of $3 \mathrm{ft} / \mathrm{s}$. Find its displacement for $t>0$.
9. A spring-mass system has natural frequency $7 \sqrt{10} \mathrm{rad} / \mathrm{s}$. The natural length of the spring is .7 m . What is the length of the spring when the mass is in equilibrium?
10. A 64 lb weight is attached to a spring with constant $k=8 \mathrm{lb} / \mathrm{ft}$ and subjected to an external force $F(t)=2 \sin t$. The weight is initially displaced 6 inches above equilibrium and given an upward velocity of $2 \mathrm{ft} / \mathrm{s}$. Find its displacement for $t>0$.
11. A unit mass hangs in equilibrium from a spring with constant $k=1 / 16$. Starting at $t=0$, a force $F(t)=3 \sin t$ is applied to the mass. Find its displacement for $t>0$.
12. $\mathrm{C} / \mathrm{G} \mathrm{A} 4 \mathrm{lb}$ weight stretches a spring 1 ft in equilibrium. An external force $F(t)=.25 \sin 8 t$ lb is applied to the weight, which is initially displaced 4 inches above equilibrium and given a downward velocity of $1 \mathrm{ft} / \mathrm{s}$. Find and graph its displacement for $t>0$.
13. A 2 lb weight stretches a spring 6 inches in equilibrium. An external force $F(t)=\sin 8 t \mathrm{lb}$ is applied to the weight, which is released from rest 2 inches below equilibrium. Find its displacement for $t>0$.
14. A 10 gm mass suspended on a spring moves in simple harmonic motion with period 4 s . Find the period of the simple harmonic motion of a 20 gm mass suspended from the same spring.
15. A 6 lb weight stretches a spring 6 inches in equilibrium. Suppose an external force $F(t)=$ $\frac{3}{16} \sin \omega t+\frac{3}{8} \cos \omega t \mathrm{lb}$ is applied to the weight. For what value of $\omega$ will the displacement be unbounded? Find the displacement if $\omega$ has this value. Assume that the motion starts from equilibrium with zero initial velocity.
16. $\mathrm{C} / \mathrm{G}$ A 6 lb weight stretches a spring 4 inches in equilibrium. Suppose an external force $F(t)=$ $4 \sin \omega t-6 \cos \omega t \mathrm{lb}$ is applied to the weight. For what value of $\omega$ will the displacement be unbounded? Find and graph the displacement if $\omega$ has this value. Assume that the motion starts from equilibrium with zero initial velocity.
17. A mass of one kg is attached to a spring with constant $k=4 \mathrm{~N} / \mathrm{m}$. An external force $F(t)=$ $-\cos \omega t-2 \sin \omega t \mathrm{n}$ is applied to the mass. Find the displacement $y$ for $t>0$ if $\omega$ equals the natural frequency of the spring-mass system. Assume that the mass is initially displaced 3 m above equilibrium and given an upward velocity of $450 \mathrm{~cm} / \mathrm{s}$.
18. An object is in simple harmonic motion with frequency $\omega_{0}$, with $y(0)=y_{0}$ and $y^{\prime}(0)=v_{0}$. Find its displacement for $t>0$. Also, find the amplitude of the oscillation and give formulas for the sine and cosine of the initial phase angle.
19. Two objects suspended from identical springs are set into motion. The period of one object is twice the period of the other. How are the weights of the two objects related?
20. Two objects suspended from identical springs are set into motion. The weight of one object is twice the weight of the other. How are the periods of the resulting motions related?
21. Two identical objects suspended from different springs are set into motion. The period of one motion is 3 times the period of the other. How are the two spring constants related?

### 6.2 SPRING PROBLEMS II

## Free Vibrations With Damping

In this section we consider the motion of an object in a spring-mass system with damping. We start with unforced motion, so the equation of motion is

$$
\begin{equation*}
m y^{\prime \prime}+c y^{\prime}+k y=0 \tag{6.2.1}
\end{equation*}
$$

Now suppose the object is displaced from equilibrium and given an initial velocity. Intuition suggests that if the damping force is sufficiently weak the resulting motion will be oscillatory, as in the undamped case considered in the previous section, while if it's sufficiently strong the object may just move slowly toward the equilibrium position without ever reaching it. We'll now confirm these intuitive ideas mathematically. The characteristic equation of (6.2.1) is

$$
m r^{2}+c r+k=0
$$

The roots of this equation are

$$
\begin{equation*}
r_{1}=\frac{-c-\sqrt{c^{2}-4 m k}}{2 m} \quad \text { and } \quad r_{2}=\frac{-c+\sqrt{c^{2}-4 m k}}{2 m} . \tag{6.2.2}
\end{equation*}
$$



Figure 6.2.1 Underdamped motion

In Section 5.2 we saw that the form of the solution of (6.2.1) depends upon whether $c^{2}-4 m k$ is positive, negative, or zero. We'll now consider these three cases.
Underdamped Motion
We say the motion is underdamped if $c<\sqrt{4 m k}$. In this case $r_{1}$ and $r_{2}$ in (6.2.2) are complex conjugates, which we write as

$$
r_{1}=-\frac{c}{2 m}-i \omega_{1} \quad \text { and } \quad r_{2}=-\frac{c}{2 m}+i \omega_{1}
$$

where

$$
\omega_{1}=\frac{\sqrt{4 m k-c^{2}}}{2 m}
$$

The general solution of (6.2.1) in this case is

$$
y=e^{-c t / 2 m}\left(c_{1} \cos \omega_{1} t+c_{2} \sin \omega_{1} t\right)
$$

By the method used in Section 6.1 to derive the amplitude-phase form of the displacement of an object in simple harmonic motion, we can rewrite this equation as

$$
\begin{equation*}
y=R e^{-c t / 2 m} \cos \left(\omega_{1} t-\phi\right) \tag{6.2.3}
\end{equation*}
$$

where

$$
R=\sqrt{c_{1}^{2}+c_{2}^{2}}, \quad R \cos \phi=c_{1}, \quad \text { and } \quad R \sin \phi=c_{2}
$$

The factor $R e^{-c t / 2 m}$ in (6.2.3) is called the time-varying amplitude of the motion, the quantity $\omega_{1}$ is called the frequency, and $T=2 \pi / \omega_{1}$ (which is the period of the cosine function in (6.2.3) is called the quasi-period. A typical graph of (6.2.3) is shown in Figure 6.2.1. As illustrated in that figure, the graph of $y$ oscillates between the dashed exponential curves $y= \pm R e^{-c t / 2 m}$.

Overdamped Motion
We say the motion is overdamped if $c>\sqrt{4 m k}$. In this case the zeros $r_{1}$ and $r_{2}$ of the characteristic polynomial are real, with $r_{1}<r_{2}<0$ (see (6.2.2)), and the general solution of (6.2.1) is

$$
y=c_{1} e^{r_{1} t}+c_{2} e^{r_{2} t} .
$$

Again $\lim _{t \rightarrow \infty} y(t)=0$ as in the underdamped case, but the motion isn't oscillatory, since $y$ can't equal zero for more than one value of $t$ unless $c_{1}=c_{2}=0$. (Exercise 23.)

## Critically Damped Motion

We say the motion is critically damped if $c=\sqrt{4 m k}$. In this case $r_{1}=r_{2}=-c / 2 m$ and the general solution of (6.2.1) is

$$
y=e^{-c t / 2 m}\left(c_{1}+c_{2} t\right) .
$$

Again $\lim _{t \rightarrow \infty} y(t)=0$ and the motion is nonoscillatory, since $y$ can't equal zero for more than one value of $t$ unless $c_{1}=c_{2}=0$. (Exercise 22).

Example 6.2.1 Suppose a 64 lb weight stretches a spring 6 inches in equilibrium and a dashpot provides a damping force of $c \mathrm{lb}$ for each $\mathrm{ft} / \mathrm{sec}$ of velocity.
(a) Write the equation of motion of the object and determine the value of $c$ for which the motion is critically damped.
(b) Find the displacement $y$ for $t>0$ if the motion is critically damped and the initial conditions are $y(0)=1$ and $y^{\prime}(0)=20$.
(c) Find the displacement $y$ for $t>0$ if the motion is critically damped and the initial conditions are $y(0)=1$ and $y^{\prime}(0)=-20$.
$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ Here $m=2$ slugs and $k=64 / .5=128 \mathrm{lb} / \mathrm{ft}$. Therefore the equation of motion (6.2.1) is

$$
\begin{equation*}
2 y^{\prime \prime}+c y^{\prime}+128 y=0 . \tag{6.2.4}
\end{equation*}
$$

The characteristic equation is

$$
2 r^{2}+c r+128=0
$$

which has roots

$$
r=\frac{-c \pm \sqrt{c^{2}-8 \cdot 128}}{4}
$$

Therefore the damping is critical if

$$
c=\sqrt{8 \cdot 128}=32 \mathrm{lb}-\mathrm{sec} / \mathrm{ft} .
$$

$\underline{\text { SOLUTION(b) }}$ Setting $c=32$ in (6.2.4) and cancelling the common factor 2 yields

$$
y^{\prime \prime}+16 y+64 y=0 .
$$

The characteristic equation is

$$
r^{2}+16 r+64 y=(r+8)^{2}=0
$$

Hence, the general solution is

$$
\begin{equation*}
y=e^{-8 t}\left(c_{1}+c_{2} t\right) \tag{6.2.5}
\end{equation*}
$$

Differentiating this yields

$$
\begin{equation*}
y^{\prime}=-8 y+c_{2} e^{-8 t} \tag{6.2.6}
\end{equation*}
$$

Imposing the initial conditions $y(0)=1$ and $y^{\prime}(0)=20$ in the last two equations shows that $1=c_{1}$ and $20=-8+c_{2}$. Hence, the solution of the initial value problem is

$$
y=e^{-8 t}(1+28 t)
$$

Therefore the object approaches equilibrium from above as $t \rightarrow \infty$. There's no oscillation.
SOLUTION(c) Imposing the initial conditions $y(0)=1$ and $y^{\prime}(0)=-20$ in (6.2.5) and (6.2.6) yields $1=c_{1}$ and $-20=-8+c_{2}$. Hence, the solution of this initial value problem is

$$
y=e^{-8 t}(1-12 t)
$$

Therefore the object moves downward through equilibrium just once, and then approaches equilibrium from below as $t \rightarrow \infty$. Again, there's no oscillation. The solutions of these two initial value problems are graphed in Figure 6.2.2.

Example 6.2.2 Find the displacement of the object in Example 6.2.1 if the damping constant is $c=4$ $\mathrm{lb}-\mathrm{sec} / \mathrm{ft}$ and the initial conditions are $y(0)=1.5 \mathrm{ft}$ and $y^{\prime}(0)=-3 \mathrm{ft} / \mathrm{sec}$.


Figure 6.2.2 (a) $y=e^{-8 t}(1+28 t) \quad$ (b) $y=e^{-8 t}(1-12 t)$

Solution With $c=4$, the equation of motion (6.2.4) becomes

$$
\begin{equation*}
y^{\prime \prime}+2 y^{\prime}+64 y=0 \tag{6.2.7}
\end{equation*}
$$

after cancelling the common factor 2 . The characteristic equation

$$
r^{2}+2 r+64=0
$$

has complex conjugate roots

$$
r=\frac{-2 \pm \sqrt{4-4 \cdot 64}}{2}=-1 \pm 3 \sqrt{7} i
$$

Therefore the motion is underdamped and the general solution of (6.2.7) is

$$
y=e^{-t}\left(c_{1} \cos 3 \sqrt{7} t+c_{2} \sin 3 \sqrt{7} t\right) .
$$

Differentiating this yields

$$
y^{\prime}=-y+3 \sqrt{7} e^{-t}\left(-c_{1} \sin 3 \sqrt{7} t+c_{2} \cos 3 \sqrt{7} t\right) .
$$

Imposing the initial conditions $y(0)=1.5$ and $y^{\prime}(0)=-3$ in the last two equations yields $1.5=c_{1}$ and $-3=-1.5+3 \sqrt{7} c_{2}$. Hence, the solution of the initial value problem is

$$
\begin{equation*}
y=e^{-t}\left(\frac{3}{2} \cos 3 \sqrt{7} t-\frac{1}{2 \sqrt{7}} \sin 3 \sqrt{7} t\right) . \tag{6.2.8}
\end{equation*}
$$

The amplitude of the function in parentheses is

$$
R=\sqrt{\left(\frac{3}{2}\right)^{2}+\left(\frac{1}{2 \sqrt{7}}\right)^{2}}=\sqrt{\frac{9}{4}+\frac{1}{4 \cdot 7}}=\sqrt{\frac{64}{4 \cdot 7}}=\frac{4}{\sqrt{7}} .
$$

Therefore we can rewrite (6.2.8) as

$$
y=\frac{4}{\sqrt{7}} e^{-t} \cos (3 \sqrt{7} t-\phi),
$$

where

$$
\cos \phi=\frac{3}{2 R}=\frac{3 \sqrt{7}}{8} \quad \text { and } \quad \sin \phi=-\frac{1}{2 \sqrt{7} R}=-\frac{1}{8} .
$$

Therefore $\phi \cong-.125$ radians.


Figure 6.2.3 $y=\frac{17}{12} e^{-4 t}-\frac{5}{12} e^{-16 t}$

Example 6.2.3 Let the damping constant in Example 1 be $c=40 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. Find the displacement $y$ for $t>0$ if $y(0)=1$ and $y^{\prime}(0)=1$.

Solution With $c=40$, the equation of motion (6.2.4) reduces to

$$
\begin{equation*}
y^{\prime \prime}+20 y^{\prime}+64 y=0 \tag{6.2.9}
\end{equation*}
$$

after cancelling the common factor 2 . The characteristic equation

$$
r^{2}+20 r+64=(r+16)(r+4)=0
$$

has the roots $r_{1}=-4$ and $r_{2}=-16$. Therefore the general solution of (6.2.9) is

$$
\begin{equation*}
y=c_{1} e^{-4 t}+c_{2} e^{-16 t} \tag{6.2.10}
\end{equation*}
$$

Differentiating this yields

$$
y^{\prime}=-4 e^{-4 t}-16 c_{2} e^{-16 t}
$$

The last two equations and the initial conditions $y(0)=1$ and $y^{\prime}(0)=1$ imply that

$$
\begin{aligned}
c_{1} & +\quad c_{2}
\end{aligned}=1.1 .
$$

The solution of this system is $c_{1}=17 / 12, c_{2}=-5 / 12$. Substituting these into (6.2.10) yields

$$
y=\frac{17}{12} e^{-4 t}-\frac{5}{12} e^{-16 t}
$$

as the solution of the given initial value problem (Figure 6.2.3).
Forced Vibrations With Damping
Now we consider the motion of an object in a spring-mass system with damping, under the influence of a periodic forcing function $F(t)=F_{0} \cos \omega t$, so that the equation of motion is

$$
\begin{equation*}
m y^{\prime \prime}+c y^{\prime}+k y=F_{0} \cos \omega t . \tag{6.2.11}
\end{equation*}
$$

In Section 6.1 we considered this equation with $c=0$ and found that the resulting displacement $y$ assumed arbitrarily large values in the case of resonance (that is, when $\omega=\omega_{0}=\sqrt{k / m}$ ). Here we'll see that in
the presence of damping the displacement remains bounded for all $t$, and the initial conditions have little effect on the motion as $t \rightarrow \infty$. In fact, we'll see that for large $t$ the displacement is closely approximated by a function of the form

$$
\begin{equation*}
y=R \cos (\omega t-\phi) \tag{6.2.12}
\end{equation*}
$$

where the amplitude $R$ depends upon $m, c, k, F_{0}$, and $\omega$. We're interested in the following question:

QUestion:Assuming that $m, c, k$, and $F_{0}$ are held constant, what value of $\omega$ produces the largest amplitude $R$ in (6.2.12), and what is this largest amplitude?

To answer this question, we must solve (6.2.11) and determine $R$ in terms of $F_{0}, \omega_{0}, \omega$, and $c$. We can obtain a particular solution of (6.2.11) by the method of undetermined coefficients. Since $\cos \omega t$ does not satisfy the complementary equation

$$
m y^{\prime \prime}+c y^{\prime}+k y=0
$$

we can obtain a particular solution of (6.2.11) in the form

$$
\begin{equation*}
y_{p}=A \cos \omega t+B \sin \omega t . \tag{6.2.13}
\end{equation*}
$$

Differentiating this yields

$$
y_{p}^{\prime}=\omega(-A \sin \omega t+B \cos \omega t)
$$

and

$$
y_{p}^{\prime \prime}=-\omega^{2}(A \cos \omega t+B \sin \omega t)
$$

From the last three equations,

$$
m y_{p}^{\prime \prime}+c y_{p}^{\prime}+k y_{p}=\left(-m \omega^{2} A+c \omega B+k A\right) \cos \omega t+\left(-m \omega^{2} B-c \omega A+k B\right) \sin \omega t
$$

so $y_{p}$ satisfies (6.2.11) if

$$
\begin{aligned}
& \left(k-m \omega^{2}\right) A+c \omega B=F_{0} \\
& -c \omega A+\left(k-m \omega^{2}\right) B=0
\end{aligned}
$$

Solving for $A$ and $B$ and substituting the results into (6.2.13) yields

$$
y_{p}=\frac{F_{0}}{\left(k-m \omega^{2}\right)^{2}+c^{2} \omega^{2}}\left[\left(k-m \omega^{2}\right) \cos \omega t+c \omega \sin \omega t\right],
$$

which can be written in amplitude-phase form as

$$
\begin{equation*}
y_{p}=\frac{F_{0}}{\sqrt{\left(k-m \omega^{2}\right)^{2}+c^{2} \omega^{2}}} \cos (\omega t-\phi) \tag{6.2.14}
\end{equation*}
$$

where

$$
\begin{equation*}
\cos \phi=\frac{k-m \omega^{2}}{\sqrt{\left(k-m \omega^{2}\right)^{2}+c^{2} \omega^{2}}} \quad \text { and } \quad \sin \phi=\frac{c \omega}{\sqrt{\left(k-m \omega^{2}\right)^{2}+c^{2} \omega^{2}}} \tag{6.2.15}
\end{equation*}
$$

To compare this with the undamped forced vibration that we considered in Section 6.1 it 's useful to write

$$
\begin{equation*}
k-m \omega^{2}=m\left(\frac{k}{m}-\omega^{2}\right)=m\left(\omega_{0}^{2}-\omega^{2}\right) \tag{6.2.16}
\end{equation*}
$$

where $\omega_{0}=\sqrt{k / m}$ is the natural angular frequency of the undamped simple harmonic motion of an object with mass $m$ on a spring with constant $k$. Substituting (6.2.16) into (6.2.14) yields

$$
\begin{equation*}
y_{p}=\frac{F_{0}}{\sqrt{m^{2}\left(\omega_{0}^{2}-\omega^{2}\right)^{2}+c^{2} \omega^{2}}} \cos (\omega t-\phi) . \tag{6.2.17}
\end{equation*}
$$

The solution of an initial value problem

$$
m y^{\prime \prime}+c y^{\prime}+k y=F_{0} \cos \omega t, \quad y(0)=y_{0}, \quad y^{\prime}(0)=v_{0}
$$

is of the form $y=y_{c}+y_{p}$, where $y_{c}$ has one of the three forms

$$
\begin{aligned}
& y_{c}=e^{-c t / 2 m}\left(c_{1} \cos \omega_{1} t+c_{2} \sin \omega_{1} t\right) \\
& y_{c}=e^{-c t / 2 m}\left(c_{1}+c_{2} t\right) \\
& y_{c}=c_{1} e^{r_{1} t}+c_{2} e^{r_{2} t}\left(r_{1}, r_{2}<0\right)
\end{aligned}
$$

In all three cases $\lim _{t \rightarrow \infty} y_{c}(t)=0$ for any choice of $c_{1}$ and $c_{2}$. For this reason we say that $y_{c}$ is the transient component of the solution $y$. The behavior of $y$ for large $t$ is determined by $y_{p}$, which we call the steady state component of $y$. Thus, for large $t$ the motion is like simple harmonic motion at the frequency of the external force.

The amplitude of $y_{p}$ in (6.2.17) is

$$
\begin{equation*}
R=\frac{F_{0}}{\sqrt{m^{2}\left(\omega_{0}^{2}-\omega^{2}\right)^{2}+c^{2} \omega^{2}}}, \tag{6.2.18}
\end{equation*}
$$

which is finite for all $\omega$; that is, the presence of damping precludes the phenomenon of resonance that we encountered in studying undamped vibrations under a periodic forcing function. We'll now find the value $\omega_{\max }$ of $\omega$ for which $R$ is maximized. This is the value of $\omega$ for which the function

$$
\rho(\omega)=m^{2}\left(\omega_{0}^{2}-\omega^{2}\right)^{2}+c^{2} \omega^{2}
$$

in the denominator of (6.2.18) attains its minimum value. By rewriting this as

$$
\begin{equation*}
\rho(\omega)=m^{2}\left(\omega_{0}^{4}+\omega^{4}\right)+\left(c^{2}-2 m^{2} \omega_{0}^{2}\right) \omega^{2} \tag{6.2.19}
\end{equation*}
$$

you can see that $\rho$ is a strictly increasing function of $\omega^{2}$ if

$$
c \geq \sqrt{2 m^{2} \omega_{0}^{2}}=\sqrt{2 m k}
$$

(Recall that $\omega_{0}^{2}=k / m$ ). Therefore $\omega_{\max }=0$ if this inequality holds. From (6.2.15), you can see that $\phi=0$ if $\omega=0$. In this case, (6.2.14) reduces to

$$
y_{p}=\frac{F_{0}}{\sqrt{m^{2} \omega_{0}^{4}}}=\frac{F_{0}}{k},
$$

which is consistent with Hooke's law: if the mass is subjected to a constant force $F_{0}$, its displacement should approach a constant $y_{p}$ such that $k y_{p}=F_{0}$. Now suppose $c<\sqrt{2 m k}$. Then, from (6.2.19),

$$
\rho^{\prime}(\omega)=2 \omega\left(2 m^{2} \omega^{2}+c^{2}-2 m^{2} \omega_{0}^{2}\right),
$$

and $\omega_{\max }$ is the value of $\omega$ for which the expression in parentheses equals zero; that is,

$$
\omega_{\max }=\sqrt{\omega_{0}^{2}-\frac{c^{2}}{2 m^{2}}}=\sqrt{\frac{k}{m}\left(1-\frac{c^{2}}{2 k m}\right)} .
$$

(To see that $\rho\left(\omega_{\max }\right)$ is the minimum value of $\rho(\omega)$, note that $\rho^{\prime}(\omega)<0$ if $\omega<\omega_{\max }$ and $\rho^{\prime}(\omega)>0$ if $\omega>\omega_{\max }$.) Substituting $\omega=\omega_{\max }$ in (6.2.18) and simplifying shows that the maximum amplitude $R_{\text {max }}$ is

$$
R_{\max }=\frac{2 m F_{0}}{c \sqrt{4 m k-c^{2}}} \quad \text { if } \quad c<\sqrt{2 m k} .
$$

We summarize our results as follows.
Theorem 6.2.1 Suppose we consider the amplitude $R$ of the steady state component of the solution of

$$
m y^{\prime \prime}+c y^{\prime}+k y=F_{0} \cos \omega t
$$

as a function of $\omega$.
(a) If $c \geq \sqrt{2 m k}$, the maximum amplitude is $R_{\max }=F_{0} / k$ and it's attained when $\omega=\omega_{\max }=0$.
(b) If $c<\sqrt{2 m k}$, the maximum amplitude is

$$
\begin{equation*}
R_{\max }=\frac{2 m F_{0}}{c \sqrt{4 m k-c^{2}}} \tag{6.2.20}
\end{equation*}
$$

and it's attained when

$$
\begin{equation*}
\omega=\omega_{\max }=\sqrt{\frac{k}{m}\left(1-\frac{c^{2}}{2 k m}\right)} . \tag{6.2.21}
\end{equation*}
$$

Note that $R_{\max }$ and $\omega_{\max }$ are continuous functions of $c$, for $c \geq 0$, since (6.2.20) and (6.2.21) reduce to $R_{\max }=F_{0} / k$ and $\omega_{\max }=0$ if $c=\sqrt{2 k m}$.

### 6.2 Exercises

1. A 64 lb object stretches a spring 4 ft in equilibrium. It is attached to a dashpot with damping constant $c=8 \mathrm{lb}$-sec/ft. The object is initially displaced 18 inches above equilibrium and given a downward velocity of $4 \mathrm{ft} / \mathrm{sec}$. Find its displacement and time-varying amplitude for $t>0$.
2. C/G A 16 lb weight is attached to a spring with natural length 5 ft . With the weight attached, the spring measures 8.2 ft . The weight is initially displaced 3 ft below equilibrium and given an upward velocity of $2 \mathrm{ft} / \mathrm{sec}$. Find and graph its displacement for $t>0$ if the medium resists the motion with a force of one lb for each $\mathrm{ft} / \mathrm{sec}$ of velocity. Also, find its time-varying amplitude.
3. $\mathrm{C} / \mathrm{G}$ An 8 lb weight stretches a spring 1.5 inches. It is attached to a dashpot with damping constant $c=8 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. The weight is initially displaced 3 inches above equilibrium and given an upward velocity of $6 \mathrm{ft} / \mathrm{sec}$. Find and graph its displacement for $t>0$.
4. A 96 lb weight stretches a spring 3.2 ft in equilibrium. It is attached to a dashpot with damping constant $c=18 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. The weight is initially displaced 15 inches below equilibrium and given a downward velocity of $12 \mathrm{ft} / \mathrm{sec}$. Find its displacement for $t>0$.
5. A 16 lb weight stretches a spring 6 inches in equilibrium. It is attached to a damping mechanism with constant $c$. Find all values of $c$ such that the free vibration of the weight has infinitely many oscillations.
6. An 8 lb weight stretches a spring .32 ft . The weight is initially displaced 6 inches above equilibrium and given an upward velocity of $4 \mathrm{ft} / \mathrm{sec}$. Find its displacement for $t>0$ if the medium exerts a damping force of 1.5 lb for each $\mathrm{ft} / \mathrm{sec}$ of velocity.
7. A 32 lb weight stretches a spring 2 ft in equilibrium. It is attached to a dashpot with constant $c=8$ $\mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. The weight is initially displaced 8 inches below equilibrium and released from rest. Find its displacement for $t>0$.
8. A mass of 20 gm stretches a spring 5 cm . The spring is attached to a dashpot with damping constant 400 dyne sec $/ \mathrm{cm}$. Determine the displacement for $t>0$ if the mass is initially displaced 9 cm above equilibrium and released from rest.
9. A 64 lb weight is suspended from a spring with constant $k=25 \mathrm{lb} / \mathrm{ft}$. It is initially displaced 18 inches above equilibrium and released from rest. Find its displacement for $t>0$ if the medium resists the motion with 6 lb of force for each $\mathrm{ft} / \mathrm{sec}$ of velocity.
10. A 32 lb weight stretches a spring 1 ft in equilibrium. The weight is initially displaced 6 inches above equilibrium and given a downward velocity of $3 \mathrm{ft} / \mathrm{sec}$. Find its displacement for $t>0$ if the medium resists the motion with a force equal to 3 times the speed in $\mathrm{ft} / \mathrm{sec}$.
11. An 8 lb weight stretches a spring 2 inches. It is attached to a dashpot with damping constant $c=4 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. The weight is initially displaced 3 inches above equilibrium and given a downward velocity of $4 \mathrm{ft} / \mathrm{sec}$. Find its displacement for $t>0$.
12. $\mathrm{C} / \mathrm{G}$ A 2 lb weight stretches a spring .32 ft . The weight is initially displaced 4 inches below equilibrium and given an upward velocity of $5 \mathrm{ft} / \mathrm{sec}$. The medium provides damping with constant $c=1 / 8 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. Find and graph the displacement for $t>0$.
13. An 8 lb weight stretches a spring 8 inches in equilibrium. It is attached to a dashpot with damping constant $c=.5 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$ and subjected to an external force $F(t)=4 \cos 2 t \mathrm{lb}$. Determine the steady state component of the displacement for $t>0$.
14. A 32 lb weight stretches a spring 1 ft in equilibrium. It is attached to a dashpot with constant $c=12 \mathrm{lb}-\mathrm{sec} / \mathrm{ft}$. The weight is initially displaced 8 inches above equilibrium and released from rest. Find its displacement for $t>0$.
15. A mass of one kg stretches a spring 49 cm in equilibrium. A dashpot attached to the spring supplies a damping force of 4 N for each $\mathrm{m} / \mathrm{sec}$ of speed. The mass is initially displaced 10 cm above equilibrium and given a downward velocity of $1 \mathrm{~m} / \mathrm{sec}$. Find its displacement for $t>0$.
16. A mass of 100 grams stretches a spring 98 cm in equilibrium. A dashpot attached to the spring supplies a damping force of 600 dynes for each $\mathrm{cm} / \mathrm{sec}$ of speed. The mass is initially displaced 10 cm above equilibrium and given a downward velocity of $1 \mathrm{~m} / \mathrm{sec}$. Find its displacement for $t>0$.
17. A 192 lb weight is suspended from a spring with constant $k=6 \mathrm{lb} / \mathrm{ft}$ and subjected to an external force $F(t)=8 \cos 3 t \mathrm{lb}$. Find the steady state component of the displacement for $t>0$ if the medium resists the motion with a force equal to 8 times the speed in $\mathrm{ft} / \mathrm{sec}$.
18. A 2 gm mass is attached to a spring with constant 20 dyne/cm. Find the steady state component of the displacement if the mass is subjected to an external force $F(t)=3 \cos 4 t-5 \sin 4 t$ dynes and a dashpot supplies 4 dynes of damping for each $\mathrm{cm} / \mathrm{sec}$ of velocity.
19. C/G A 96 lb weight is attached to a spring with constant $12 \mathrm{lb} / \mathrm{ft}$. Find and graph the steady state component of the displacement if the mass is subjected to an external force $F(t)=18 \cos t-9 \sin t$ lb and a dashpot supplies 24 lb of damping for each $\mathrm{ft} / \mathrm{sec}$ of velocity.
20. A mass of one kg stretches a spring 49 cm in equilibrium. It is attached to a dashpot that supplies a damping force of 4 N for each $\mathrm{m} / \mathrm{sec}$ of speed. Find the steady state component of its displacement if it's subjected to an external force $F(t)=8 \sin 2 t-6 \cos 2 t \mathrm{~N}$.
21. A mass $m$ is suspended from a spring with constant $k$ and subjected to an external force $F(t)=$ $\alpha \cos \omega_{0} t+\beta \sin \omega_{0} t$, where $\omega_{0}$ is the natural frequency of the spring-mass system without damping. Find the steady state component of the displacement if a dashpot with constant $c$ supplies damping.
22. Show that if $c_{1}$ and $c_{2}$ are not both zero then

$$
y=e^{r_{1} t}\left(c_{1}+c_{2} t\right)
$$

can't equal zero for more than one value of $t$.
23. Show that if $c_{1}$ and $c_{2}$ are not both zero then

$$
y=c_{1} e^{r_{1} t}+c_{2} e^{r_{2} t}
$$

can't equal zero for more than one value of $t$.
24. Find the solution of the initial value problem

$$
m y^{\prime \prime}+c y^{\prime}+k y=0, \quad y(0)=y_{0}, y^{\prime}(0)=v_{0}
$$

given that the motion is underdamped, so the general solution of the equation is

$$
y=e^{-c t / 2 m}\left(c_{1} \cos \omega_{1} t+c_{2} \sin \omega_{1} t\right)
$$

25. Find the solution of the initial value problem

$$
m y^{\prime \prime}+c y^{\prime}+k y=0, \quad y(0)=y_{0}, y^{\prime}(0)=v_{0}
$$

given that the motion is overdamped, so the general solution of the equation is

$$
y=c_{1} e^{r_{1} t}+c_{2} e^{r_{2} t}\left(r_{1}, r_{2}<0\right) .
$$

26. Find the solution of the initial value problem

$$
m y^{\prime \prime}+c y^{\prime}+k y=0, \quad y(0)=y_{0}, y^{\prime}(0)=v_{0}
$$

given that the motion is critically damped, so that the general solution of the equation is of the form

$$
y=e^{r_{1} t}\left(c_{1}+c_{2} t\right)\left(r_{1}<0\right)
$$

# CHAPTER 7 <br> Series Solutions of Linear Second Order Equations 

IN THIS CHAPTER we study a class of second order differential equations that occur in many applications, but can't be solved in closed form in terms of elementary functions. Here are some examples:
(1) Bessel's equation

$$
x^{2} y^{\prime \prime}+x y^{\prime}+\left(x^{2}-\nu^{2}\right) y=0
$$

which occurs in problems displaying cylindrical symmetry, such as diffraction of light through a circular aperture, propagation of electromagnetic radiation through a coaxial cable, and vibrations of a circular drum head.
(2) Airy's equation,

$$
y^{\prime \prime}-x y=0
$$

which occurs in astronomy and quantum physics.
(3) Legendre's equation

$$
\left(1-x^{2}\right) y^{\prime \prime}-2 x y^{\prime}+\alpha(\alpha+1) y=0
$$

which occurs in problems displaying spherical symmetry, particularly in electromagnetism.
These equations and others considered in this chapter can be written in the form

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0 \tag{A}
\end{equation*}
$$

where $P_{0}, P_{1}$, and $P_{2}$ are polynomials with no common factor. For most equations that occur in applications, these polynomials are of degree two or less. We'll impose this restriction, although the methods that we'll develop can be extended to the case where the coefficient functions are polynomials of arbitrary degree, or even power series that converge in some circle around the origin in the complex plane.

Since (A) does not in general have closed form solutions, we seek series representations for solutions. We'll see that if $P_{0}(0) \neq 0$ then solutions of $(\mathrm{A})$ can be written as power series

$$
y=\sum_{n=0}^{\infty} a_{n} x^{n}
$$

that converge in an open interval centered at $x=0$.

SECTION 7.1 reviews the properties of power series.
SECTIONS 7.2 AND 7.3 are devoted to finding power series solutions of $(\mathrm{A})$ in the case where $P_{0}(0) \neq 0$. The situation is more complicated if $P_{0}(0)=0$; however, if $P_{1}$ and $P_{2}$ satisfy assumptions that apply to most equations of interest, then we're able to use a modified series method to obtain solutions of (A).

SECTION 7.4 introduces the appropriate assumptions on $P_{1}$ and $P_{2}$ in the case where $P_{0}(0)=0$, and deals with Euler's equation

$$
a x^{2} y^{\prime \prime}+b x y^{\prime}+c y=0
$$

where $a, b$, and $c$ are constants. This is the simplest equation that satisfies these assumptions.
SECTIONS 7.5-7.7 deal with three distinct cases satisfying the assumptions introduced in Section 7.4. In all three cases, (A) has at least one solution of the form

$$
y_{1}=x^{r} \sum_{n=0}^{\infty} a_{n} x^{n}
$$

where $r$ need not be an integer. The problem is that there are three possibilities - each requiring a different approach - for the form of a second solution $y_{2}$ such that $\left\{y_{1}, y_{2}\right\}$ is a fundamental pair of solutions of (A).

### 7.1 REVIEW OF POWER SERIES

Many applications give rise to differential equations with solutions that can't be expressed in terms of elementary functions such as polynomials, rational functions, exponential and logarithmic functions, and trigonometric functions. The solutions of some of the most important of these equations can be expressed in terms of power series. We'll study such equations in this chapter. In this section we review relevant properties of power series. We'll omit proofs, which can be found in any standard calculus text.

Definition 7.1.1 An infinite series of the form

$$
\begin{equation*}
\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n} \tag{7.1.1}
\end{equation*}
$$

where $x_{0}$ and $a_{0}, a_{1}, \ldots, a_{n}, \ldots$ are constants, is called a power series in $x-x_{0}$. We say that the power series (7.1.1) converges for a given $x$ if the limit

$$
\lim _{N \rightarrow \infty} \sum_{n=0}^{N} a_{n}\left(x-x_{0}\right)^{n}
$$

exists; otherwise, we say that the power series diverges for the given $x$.
A power series in $x-x_{0}$ must converge if $x=x_{0}$, since the positive powers of $x-x_{0}$ are all zero in this case. This may be the only value of $x$ for which the power series converges. However, the next theorem shows that if the power series converges for some $x \neq x_{0}$ then the set of all values of $x$ for which it converges forms an interval.

Theorem 7.1.2 For any power series

$$
\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

exactly one of the these statements is true:
(i) The power series converges only for $x=x_{0}$.
(ii) The power series converges for all values of $x$.
(iii) There's a positive number $R$ such that the power series converges if $\left|x-x_{0}\right|<R$ and diverges if $\left|x-x_{0}\right|>R$.

In case (iii) we say that $R$ is the radius of convergence of the power series. For convenience, we include the other two cases in this definition by defining $R=0$ in case (i) and $R=\infty$ in case (ii). We define the open interval of convergence of $\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}$ to be

$$
\left(x_{0}-R, x_{0}+R\right) \quad \text { if } \quad 0<R<\infty, \quad \text { or } \quad(-\infty, \infty) \quad \text { if } \quad R=\infty
$$

If $R$ is finite, no general statement can be made concerning convergence at the endpoints $x=x_{0} \pm R$ of the open interval of convergence; the series may converge at one or both points, or diverge at both.

Recall from calculus that a series of constants $\sum_{n=0}^{\infty} \alpha_{n}$ is said to converge absolutely if the series of absolute values $\sum_{n=0}^{\infty}\left|\alpha_{n}\right|$ converges. It can be shown that a power series $\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}$ with a positive radius of convergence $R$ converges absolutely in its open interval of convergence; that is, the series

$$
\sum_{n=0}^{\infty}\left|a_{n}\right|\left|x-x_{0}\right|^{n}
$$

of absolute values converges if $\left|x-x_{0}\right|<R$. However, if $R<\infty$, the series may fail to converge absolutely at an endpoint $x_{0} \pm R$, even if it converges there.

The next theorem provides a useful method for determining the radius of convergence of a power series. It's derived in calculus by applying the ratio test to the corresponding series of absolute values. For related theorems see Exercises 2 and 4.
Theorem 7.1.3 Suppose there's an integer $N$ such that $a_{n} \neq 0$ if $n \geq N$ and

$$
\lim _{n \rightarrow \infty}\left|\frac{a_{n+1}}{a_{n}}\right|=L
$$

where $0 \leq L \leq \infty$. Then the radius of convergence of $\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}$ is $R=1 / L$, which should be interpreted to mean that $R=0$ if $L=\infty$, or $R=\infty$ if $L=0$.

Example 7.1.1 Find the radius of convergence of the series:
(a) $\sum_{n=0}^{\infty} n!x^{n}$
(b) $\quad \sum_{n=10}^{\infty}(-1)^{n} \frac{x^{n}}{n!}$
(c) $\quad \sum_{n=0}^{\infty} 2^{n} n^{2}(x-1)^{n}$.
$\underline{\text { SOLUTION(a) Here } a_{n}=n \text { !, so }}$

$$
\lim _{n \rightarrow \infty}\left|\frac{a_{n+1}}{a_{n}}\right|=\lim _{n \rightarrow \infty} \frac{(n+1)!}{n!}=\lim _{n \rightarrow \infty}(n+1)=\infty
$$

Hence, $R=0$.
$\underline{\text { SOLUTION(b) Here } a_{n}=(1)^{n} / n!\text { for } n \geq N=10 \text {, so }}$

$$
\lim _{n \rightarrow \infty}\left|\frac{a_{n+1}}{a_{n}}\right|=\lim _{n \rightarrow \infty} \frac{n!}{(n+1)!}=\lim _{n \rightarrow \infty} \frac{1}{n+1}=0
$$

Hence, $R=\infty$.
$\underline{\text { SOLUTION(c) Here } a_{n}=2^{n} n^{2} \text {, so }}$

$$
\lim _{n \rightarrow \infty}\left|\frac{a_{n+1}}{a_{n}}\right|=\lim _{n \rightarrow \infty} \frac{2^{n+1}(n+1)^{2}}{2^{n} n^{2}}=2 \lim _{n \rightarrow \infty}\left(1+\frac{1}{n}\right)^{2}=2 .
$$

Hence, $R=1 / 2$.
Taylor Series
If a function $f$ has derivatives of all orders at a point $x=x_{0}$, then the Taylor series of $f$ about $x_{0}$ is defined by

$$
\sum_{n=0}^{\infty} \frac{f^{(n)}\left(x_{0}\right)}{n!}\left(x-x_{0}\right)^{n}
$$

In the special case where $x_{0}=0$, this series is also called the Maclaurin series of $f$.
Taylor series for most of the common elementary functions converge to the functions on their open intervals of convergence. For example, you are probably familiar with the following Maclaurin series:

$$
\begin{align*}
e^{x} & =\sum_{n=0}^{\infty} \frac{x^{n}}{n!}, \quad-\infty<x<\infty  \tag{7.1.2}\\
\sin x & =\sum_{n=0}^{\infty}(-1)^{n} \frac{x^{2 n+1}}{(2 n+1)!}, \quad-\infty<x<\infty  \tag{7.1.3}\\
\cos x & =\sum_{n=0}^{\infty}(-1)^{n} \frac{x^{2 n}}{(2 n)!}, \quad-\infty<x<\infty  \tag{7.1.4}\\
\frac{1}{1-x} & =\sum_{n=0}^{\infty} x^{n}, \quad-1<x<1 \tag{7.1.5}
\end{align*}
$$

Differentiation of Power Series
A power series with a positive radius of convergence defines a function

$$
f(x)=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

on its open interval of convergence. We say that the series represents $f$ on the open interval of convergence. A function $f$ represented by a power series may be a familiar elementary function as in (7.1.2)(7.1.5); however, it often happens that $f$ isn't a familiar function, so the series actually defines $f$.

The next theorem shows that a function represented by a power series has derivatives of all orders on the open interval of convergence of the power series, and provides power series representations of the derivatives.

Theorem 7.1.4 A power series

$$
f(x)=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

with positive radius of convergence $R$ has derivatives of all orders in its open interval of convergence, and successive derivatives can be obtained by repeatedly differentiating term by term; that is,

$$
\begin{align*}
f^{\prime}(x) & =\sum_{n=1}^{\infty} n a_{n}\left(x-x_{0}\right)^{n-1},  \tag{7.1.6}\\
f^{\prime \prime}(x) & =\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2},  \tag{7.1.7}\\
& \vdots \\
f^{(k)}(x) & =\sum_{n=k}^{\infty} n(n-1) \cdots(n-k+1) a_{n}\left(x-x_{0}\right)^{n-k} . \tag{7.1.8}
\end{align*}
$$

Moreover, all of these series have the same radius of convergence $R$.

Example 7.1.2 Let $f(x)=\sin x$. From (7.1.3),

$$
f(x)=\sum_{n=0}^{\infty}(-1)^{n} \frac{x^{2 n+1}}{(2 n+1)!} .
$$

From (7.1.6),

$$
f^{\prime}(x)=\sum_{n=0}^{\infty}(-1)^{n} \frac{d}{d x}\left[\frac{x^{2 n+1}}{(2 n+1)!}\right]=\sum_{n=0}^{\infty}(-1)^{n} \frac{x^{2 n}}{(2 n)!},
$$

which is the series (7.1.4) for $\cos x$.

Uniqueness of Power Series
The next theorem shows that if $f$ is defined by a power series in $x-x_{0}$ with a positive radius of convergence, then the power series is the Taylor series of $f$ about $x_{0}$.

Theorem 7.1.5 If the power series

$$
f(x)=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

has a positive radius of convergence, then

$$
\begin{equation*}
a_{n}=\frac{f^{(n)}\left(x_{0}\right)}{n!} \tag{7.1.9}
\end{equation*}
$$

that is, $\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}$ is the Taylor series of $f$ about $x_{0}$.
This result can be obtained by setting $x=x_{0}$ in (7.1.8), which yields

$$
f^{(k)}\left(x_{0}\right)=k(k-1) \cdots 1 \cdot a_{k}=k!a_{k}
$$

This implies that

$$
a_{k}=\frac{f^{(k)}\left(x_{0}\right)}{k!} .
$$

Except for notation, this is the same as (7.1.9).
The next theorem lists two important properties of power series that follow from Theorem 7.1.5.

## Theorem 7.1.6

(a) If

$$
\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}=\sum_{n=0}^{\infty} b_{n}\left(x-x_{0}\right)^{n}
$$

for all $x$ in an open interval that contains $x_{0}$, then $a_{n}=b_{n}$ for $n=0,1,2, \ldots$.
(b) $I f$

$$
\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}=0
$$

for all $x$ in an open interval that contains $x_{0}$, then $a_{n}=0$ for $n=0,1,2, \ldots$.
To obtain (a) we observe that the two series represent the same function $f$ on the open interval; hence, Theorem 7.1.5 implies that

$$
a_{n}=b_{n}=\frac{f^{(n)}\left(x_{0}\right)}{n!}, \quad n=0,1,2, \ldots
$$

(b) can be obtained from (a) by taking $b_{n}=0$ for $n=0,1,2, \ldots$.

Taylor Polynomials
If $f$ has $N$ derivatives at a point $x_{0}$, we say that

$$
T_{N}(x)=\sum_{n=0}^{N} \frac{f^{(n)}\left(x_{0}\right)}{n!}\left(x-x_{0}\right)^{n}
$$

is the $N$-th Taylor polynomial of $f$ about $x_{0}$. This definition and Theorem 7.1.5 imply that if

$$
f(x)=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n},
$$

where the power series has a positive radius of convergence, then the Taylor polynomials of $f$ about $x_{0}$ are given by

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n}\left(x-x_{0}\right)^{n} .
$$

In numerical applications, we use the Taylor polynomials to approximate $f$ on subintervals of the open interval of convergence of the power series. For example, (7.1.2) implies that the Taylor polynomial $T_{N}$ of $f(x)=e^{x}$ is

$$
T_{N}(x)=\sum_{n=0}^{N} \frac{x^{n}}{n!}
$$

The solid curve in Figure 7.1.1 is the graph of $y=e^{x}$ on the interval [0,5]. The dotted curves in Figure 7.1.1 are the graphs of the Taylor polynomials $T_{1}, \ldots, T_{6}$ of $y=e^{x}$ about $x_{0}=0$. From this figure, we conclude that the accuracy of the approximation of $y=e^{x}$ by its Taylor polynomial $T_{N}$ improves as $N$ increases.
Shifting the Summation Index
In Definition 7.1.1 of a power series in $x-x_{0}$, the $n$-th term is a constant multiple of $\left(x-x_{0}\right)^{n}$. This isn't true in (7.1.6), (7.1.7), and (7.1.8), where the general terms are constant multiples of $\left(x-x_{0}\right)^{n-1}$, $\left(x-x_{0}\right)^{n-2}$, and $\left(x-x_{0}\right)^{n-k}$, respectively. However, these series can all be rewritten so that their $n$-th terms are constant multiples of $\left(x-x_{0}\right)^{n}$. For example, letting $n=k+1$ in the series in (7.1.6) yields

$$
\begin{equation*}
f^{\prime}(x)=\sum_{k=0}^{\infty}(k+1) a_{k+1}\left(x-x_{0}\right)^{k} \tag{7.1.10}
\end{equation*}
$$

where we start the new summation index $k$ from zero so that the first term in (7.1.10) (obtained by setting $k=0$ ) is the same as the first term in (7.1.6) (obtained by setting $n=1$ ). However, the sum of a series is


Figure 7.1.1 Approximation of $y=e^{x}$ by Taylor polynomials about $x=0$
independent of the symbol used to denote the summation index, just as the value of a definite integral is independent of the symbol used to denote the variable of integration. Therefore we can replace $k$ by $n$ in (7.1.10) to obtain

$$
\begin{equation*}
f^{\prime}(x)=\sum_{n=0}^{\infty}(n+1) a_{n+1}\left(x-x_{0}\right)^{n} \tag{7.1.11}
\end{equation*}
$$

where the general term is a constant multiple of $\left(x-x_{0}\right)^{n}$.
It isn't really necessary to introduce the intermediate summation index $k$. We can obtain (7.1.11) directly from (7.1.6) by replacing $n$ by $n+1$ in the general term of (7.1.6) and subtracting 1 from the lower limit of (7.1.6). More generally, we use the following procedure for shifting indices.

## Shifting the Summation Index in a Power Series

For any integer $k$, the power series

$$
\sum_{n=n_{0}}^{\infty} b_{n}\left(x-x_{0}\right)^{n-k}
$$

can be rewritten as

$$
\sum_{n=n_{0}-k}^{\infty} b_{n+k}\left(x-x_{0}\right)^{n}
$$

that is, replacing $n$ by $n+k$ in the general term and subtracting $k$ from the lower limit of summation leaves the series unchanged.

Example 7.1.3 Rewrite the following power series from (7.1.7) and (7.1.8) so that the general term in each is a constant multiple of $\left(x-x_{0}\right)^{n}$ :
(a) $\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}$
(b) $\sum_{n=k}^{\infty} n(n-1) \cdots(n-k+1) a_{n}\left(x-x_{0}\right)^{n-k}$.
$\underline{\text { Solution(a) }}$ Replacing $n$ by $n+2$ in the general term and subtracting 2 from the lower limit of
summation yields

$$
\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}=\sum_{n=0}^{\infty}(n+2)(n+1) a_{n+2}\left(x-x_{0}\right)^{n}
$$

SOLUTION(b) Replacing $n$ by $n+k$ in the general term and subtracting $k$ from the lower limit of summation yields

$$
\sum_{n=k}^{\infty} n(n-1) \cdots(n-k+1) a_{n}\left(x-x_{0}\right)^{n-k}=\sum_{n=0}^{\infty}(n+k)(n+k-1) \cdots(n+1) a_{n+k}\left(x-x_{0}\right)^{n} .
$$

Example 7.1.4 Given that

$$
f(x)=\sum_{n=0}^{\infty} a_{n} x^{n}
$$

write the function $x f^{\prime \prime}$ as a power series in which the general term is a constant multiple of $x^{n}$.

Solution From Theorem 7.1.4 with $x_{0}=0$,

$$
f^{\prime \prime}(x)=\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2} .
$$

Therefore

$$
x f^{\prime \prime}(x)=\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-1}
$$

Replacing $n$ by $n+1$ in the general term and subtracting 1 from the lower limit of summation yields

$$
x f^{\prime \prime}(x)=\sum_{n=1}^{\infty}(n+1) n a_{n+1} x^{n}
$$

We can also write this as

$$
x f^{\prime \prime}(x)=\sum_{n=0}^{\infty}(n+1) n a_{n+1} x^{n}
$$

since the first term in this last series is zero. (We'll see later that sometimes it's useful to include zero terms at the beginning of a series.)
Linear Combinations of Power Series
If a power series is multiplied by a constant, then the constant can be placed inside the summation; that is,

$$
c \sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}=\sum_{n=0}^{\infty} c a_{n}\left(x-x_{0}\right)^{n} .
$$

Two power series

$$
f(x)=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n} \quad \text { and } \quad g(x)=\sum_{n=0}^{\infty} b_{n}\left(x-x_{0}\right)^{n}
$$

with positive radii of convergence can be added term by term at points common to their open intervals of convergence; thus, if the first series converges for $\left|x-x_{0}\right|<R_{1}$ and the second converges for $\left|x-x_{0}\right|<$ $R_{2}$, then

$$
f(x)+g(x)=\sum_{n=0}^{\infty}\left(a_{n}+b_{n}\right)\left(x-x_{0}\right)^{n}
$$

for $\left|x-x_{0}\right|<R$, where $R$ is the smaller of $R_{1}$ and $R_{2}$. More generally, linear combinations of power series can be formed term by term; for example,

$$
c_{1} f(x)+c_{2} g(x)=\sum_{n=0}^{\infty}\left(c_{1} a_{n}+c_{2} b_{n}\right)\left(x-x_{0}\right)^{n} .
$$

Example 7.1.5 Find the Maclaurin series for $\cosh x$ as a linear combination of the Maclaurin series for $e^{x}$ and $e^{-x}$.

Solution By definition,

$$
\cosh x=\frac{1}{2} e^{x}+\frac{1}{2} e^{-x} .
$$

Since

$$
e^{x}=\sum_{n=0}^{\infty} \frac{x^{n}}{n!} \quad \text { and } \quad e^{-x}=\sum_{n=0}^{\infty}(-1)^{n} \frac{x^{n}}{n!}
$$

it follows that

$$
\begin{equation*}
\cosh x=\sum_{n=0}^{\infty} \frac{1}{2}\left[1+(-1)^{n}\right] \frac{x^{n}}{n!} . \tag{7.1.12}
\end{equation*}
$$

Since

$$
\frac{1}{2}\left[1+(-1)^{n}\right]= \begin{cases}1 & \text { if } n=2 m, \text { an even integer } \\ 0 & \text { if } n=2 m+1, \text { an odd integer }\end{cases}
$$

we can rewrite (7.1.12) more simply as

$$
\cosh x=\sum_{m=0}^{\infty} \frac{x^{2 m}}{(2 m)!}
$$

This result is valid on $(-\infty, \infty)$, since this is the open interval of convergence of the Maclaurin series for $e^{x}$ and $e^{-x}$.

Example 7.1.6 Suppose

$$
y=\sum_{n=0}^{\infty} a_{n} x^{n}
$$

on an open interval $I$ that contains the origin.
(a) Express

$$
(2-x) y^{\prime \prime}+2 y
$$

as a power series in $x$ on $I$.
(b) Use the result of (a) to find necessary and sufficient conditions on the coefficients $\left\{a_{n}\right\}$ for $y$ to be a solution of the homogeneous equation

$$
\begin{equation*}
(2-x) y^{\prime \prime}+2 y=0 \tag{7.1.13}
\end{equation*}
$$

on $I$.
$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ From (7.1.7) with $x_{0}=0$,

$$
y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2} .
$$

Therefore

$$
\begin{align*}
(2-x) y^{\prime \prime}+2 y & =2 y^{\prime \prime}-x y^{\prime \prime}+2 y \\
& =\sum_{n=2}^{\infty} 2 n(n-1) a_{n} x^{n-2}-\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-1}+\sum_{n=0}^{\infty} 2 a_{n} x^{n} . \tag{7.1.14}
\end{align*}
$$

To combine the three series we shift indices in the first two to make their general terms constant multiples of $x^{n}$; thus,

$$
\begin{equation*}
\sum_{n=2}^{\infty} 2 n(n-1) a_{n} x^{n-2}=\sum_{n=0}^{\infty} 2(n+2)(n+1) a_{n+2} x^{n} \tag{7.1.15}
\end{equation*}
$$

and

$$
\begin{equation*}
\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-1}=\sum_{n=1}^{\infty}(n+1) n a_{n+1} x^{n}=\sum_{n=0}^{\infty}(n+1) n a_{n+1} x^{n} \tag{7.1.16}
\end{equation*}
$$

where we added a zero term in the last series so that when we substitute from (7.1.15) and (7.1.16) into (7.1.14) all three series will start with $n=0$; thus,

$$
\begin{equation*}
(2-x) y^{\prime \prime}+2 y=\sum_{n=0}^{\infty}\left[2(n+2)(n+1) a_{n+2}-(n+1) n a_{n+1}+2 a_{n}\right] x^{n} \tag{7.1.17}
\end{equation*}
$$

$\underline{\text { SOLUTION(b) From (7.1.17) we see that } y \text { satisfies (7.1.13) on } I \text { if }}$

$$
\begin{equation*}
2(n+2)(n+1) a_{n+2}-(n+1) n a_{n+1}+2 a_{n}=0, \quad n=0,1,2, \ldots \tag{7.1.18}
\end{equation*}
$$

Conversely, Theorem 7.1.6 (b) implies that if $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ satisfies (7.1.13) on $I$, then (7.1.18) holds.
Example 7.1.7 Suppose

$$
y=\sum_{n=0}^{\infty} a_{n}(x-1)^{n}
$$

on an open interval $I$ that contains $x_{0}=1$. Express the function

$$
\begin{equation*}
(1+x) y^{\prime \prime}+2(x-1)^{2} y^{\prime}+3 y \tag{7.1.19}
\end{equation*}
$$

as a power series in $x-1$ on $I$.

Solution Since we want a power series in $x-1$, we rewrite the coefficient of $y^{\prime \prime}$ in (7.1.19) as $1+x=$ $2+(x-1)$, so (7.1.19) becomes

$$
2 y^{\prime \prime}+(x-1) y^{\prime \prime}+2(x-1)^{2} y^{\prime}+3 y .
$$

From (7.1.6) and (7.1.7) with $x_{0}=1$,

$$
y^{\prime}=\sum_{n=1}^{\infty} n a_{n}(x-1)^{n-1} \quad \text { and } \quad y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n}(x-1)^{n-2} .
$$

Therefore

$$
\begin{aligned}
2 y^{\prime \prime} & =\sum_{n=2}^{\infty} 2 n(n-1) a_{n}(x-1)^{n-2} \\
(x-1) y^{\prime \prime} & =\sum_{n=2}^{\infty} n(n-1) a_{n}(x-1)^{n-1} \\
2(x-1)^{2} y^{\prime} & =\sum_{n=1}^{\infty} 2 n a_{n}(x-1)^{n+1} \\
3 y & =\sum_{n=0}^{\infty} 3 a_{n}(x-1)^{n} .
\end{aligned}
$$

Before adding these four series we shift indices in the first three so that their general terms become constant multiples of $(x-1)^{n}$. This yields

$$
\begin{align*}
2 y^{\prime \prime} & =\sum_{n=0}^{\infty} 2(n+2)(n+1) a_{n+2}(x-1)^{n}  \tag{7.1.20}\\
(x-1) y^{\prime \prime} & =\sum_{n=0}^{\infty}(n+1) n a_{n+1}(x-1)^{n}  \tag{7.1.21}\\
2(x-1)^{2} y^{\prime} & =\sum_{n=1}^{\infty} 2(n-1) a_{n-1}(x-1)^{n},  \tag{7.1.22}\\
3 y & =\sum_{n=0}^{\infty} 3 a_{n}(x-1)^{n} \tag{7.1.23}
\end{align*}
$$

where we added initial zero terms to the series in (7.1.21) and (7.1.22). Adding (7.1.20)-(7.1.23) yields

$$
\begin{aligned}
(1+x) y^{\prime \prime}+2(x-1)^{2} y^{\prime}+3 y & =2 y^{\prime \prime}+(x-1) y^{\prime \prime}+2(x-1)^{2} y^{\prime}+3 y \\
& =\sum_{n=0}^{\infty} b_{n}(x-1)^{n}
\end{aligned}
$$

where

$$
\begin{align*}
& b_{0}=4 a_{2}+3 a_{0}  \tag{7.1.24}\\
& b_{n}=2(n+2)(n+1) a_{n+2}+(n+1) n a_{n+1}+2(n-1) a_{n-1}+3 a_{n}, n \geq 1 \tag{7.1.25}
\end{align*}
$$

The formula (7.1.24) for $b_{0}$ can't be obtained by setting $n=0$ in (7.1.25), since the summation in (7.1.22) begins with $n=1$, while those in (7.1.20), (7.1.21), and (7.1.23) begin with $n=0$.

### 7.1 Exercises

1. For each power series use Theorem 7.1.3 to find the radius of convergence $R$. If $R>0$, find the open interval of convergence.
(a) $\sum_{n=0}^{\infty} \frac{(-1)^{n}}{2^{n} n}(x-1)^{n}$
(b) $\sum_{n=0}^{\infty} 2^{n} n(x-2)^{n}$
(c) $\sum_{n=0}^{\infty} \frac{n!}{9^{n}} x^{n}$
(d) $\sum_{n=0}^{\infty} \frac{n(n+1)}{16^{n}}(x-2)^{n}$
(e) $\sum_{n=0}^{\infty}(-1)^{n} \frac{7^{n}}{n!} x^{n}$
(f) $\sum_{n=0}^{\infty} \frac{3^{n}}{4^{n+1}(n+1)^{2}}(x+7)^{n}$
2. Suppose there's an integer $M$ such that $b_{m} \neq 0$ for $m \geq M$, and

$$
\lim _{m \rightarrow \infty}\left|\frac{b_{m+1}}{b_{m}}\right|=L
$$

where $0 \leq L \leq \infty$. Show that the radius of convergence of

$$
\sum_{m=0}^{\infty} b_{m}\left(x-x_{0}\right)^{2 m}
$$

is $R=1 / \sqrt{L}$, which is interpreted to mean that $R=0$ if $L=\infty$ or $R=\infty$ if $L=0$. Hint: Apply Theorem 7.1.3 to the series $\sum_{m=0}^{\infty} b_{m} z^{m}$ and then let $z=\left(x-x_{0}\right)^{2}$.
3. For each power series, use the result of Exercise 2 to find the radius of convergence $R$. If $R>0$, find the open interval of convergence.
(a) $\sum_{m=0}^{\infty}(-1)^{m}(3 m+1)(x-1)^{2 m+1}$
(b) $\sum_{m=0}^{\infty}(-1)^{m} \frac{m(2 m+1)}{2^{m}}(x+2)^{2 m}$
(c) $\sum_{m=0}^{\infty} \frac{m!}{(2 m)!}(x-1)^{2 m}$
(d) $\sum_{m=0}^{\infty}(-1)^{m} \frac{m!}{9^{m}}(x+8)^{2 m}$
(e) $\sum_{m=0}^{\infty}(-1)^{m} \frac{(2 m-1)}{3^{m}} x^{2 m+1}$
(f) $\sum_{m=0}^{\infty}(x-1)^{2 m}$
4. Suppose there's an integer $M$ such that $b_{m} \neq 0$ for $m \geq M$, and

$$
\lim _{m \rightarrow \infty}\left|\frac{b_{m+1}}{b_{m}}\right|=L
$$

where $0 \leq L \leq \infty$. Let $k$ be a positive integer. Show that the radius of convergence of

$$
\sum_{m=0}^{\infty} b_{m}\left(x-x_{0}\right)^{k m}
$$

is $R=1 / \sqrt[k]{L}$, which is interpreted to mean that $R=0$ if $L=\infty$ or $R=\infty$ if $L=0$. Hint: Apply Theorem 7.1.3 to the series $\sum_{m=0}^{\infty} b_{m} z^{m}$ and then let $z=\left(x-x_{0}\right)^{k}$.
5. For each power series use the result of Exercise 4 to find the radius of convergence $R$. If $R>0$, find the open interval of convergence.
(a) $\sum_{\substack{m=0 \\ \infty}}^{\infty} \frac{(-1)^{m}}{(27)^{m}}(x-3)^{3 m+2}$
(b) $\sum_{\substack{m=0 \\ \infty}}^{\infty} \frac{x^{7 m+6}}{m}$
(c) $\sum_{m=0}^{\infty} \frac{9^{m}(m+1)}{(m+2)}(x-3)^{4 m+2}$
(d) $\sum_{m=0}^{\infty}(-1)^{m} \frac{2^{m}}{m!} x^{4 m+3}$
(e) $\sum_{m=0}^{\infty} \frac{m!}{(26)^{m}}(x+1)^{4 m+3}$
(f) $\sum_{m=0}^{m=0} \frac{(-1)^{m}}{8^{m} m(m+1)}(x-1)^{3 m+1}$
6. L Graph $y=\sin x$ and the Taylor polynomial

$$
T_{2 M+1}(x)=\sum_{n=0}^{M} \frac{(-1)^{n} x^{2 n+1}}{(2 n+1)!}
$$

on the interval $(-2 \pi, 2 \pi)$ for $M=1,2,3, \ldots$, until you find a value of $M$ for which there's no perceptible difference between the two graphs.
7. L Graph $y=\cos x$ and the Taylor polynomial

$$
T_{2 M}(x)=\sum_{n=0}^{M} \frac{(-1)^{n} x^{2 n}}{(2 n)!}
$$

on the interval $(-2 \pi, 2 \pi)$ for $M=1,2,3, \ldots$, until you find a value of $M$ for which there's no perceptible difference between the two graphs.
8. L Graph $y=1 /(1-x)$ and the Taylor polynomial

$$
T_{N}(x)=\sum_{n=0}^{N} x^{n}
$$

on the interval $[0, .95]$ for $N=1,2,3, \ldots$, until you find a value of $N$ for which there's no perceptible difference between the two graphs. Choose the scale on the $y$-axis so that $0 \leq y \leq 20$.
9. L Graph $y=\cosh x$ and the Taylor polynomial

$$
T_{2 M}(x)=\sum_{n=0}^{M} \frac{x^{2 n}}{(2 n)!}
$$

on the interval $(-5,5)$ for $M=1,2,3, \ldots$, until you find a value of $M$ for which there's no perceptible difference between the two graphs. Choose the scale on the $y$-axis so that $0 \leq y \leq 75$.
10. L Graph $y=\sinh x$ and the Taylor polynomial

$$
T_{2 M+1}(x)=\sum_{n=0}^{M} \frac{x^{2 n+1}}{(2 n+1)!}
$$

on the interval $(-5,5)$ for $M=0,1,2, \ldots$, until you find a value of $M$ for which there's no perceptible difference between the two graphs. Choose the scale on the $y$-axis so that $-75 \leq y \leq 75$.

In Exercises 11-15 find a power series solution $y(x)=\sum_{n=0}^{\infty} a_{n} x^{n}$.
11. $(2+x) y^{\prime \prime}+x y^{\prime}+3 y$
13. $\left(1+2 x^{2}\right) y^{\prime \prime}+(2-3 x) y^{\prime}+4 y$
14. $\left(1+x^{2}\right) y^{\prime \prime}+(2-x) y^{\prime}+3 y$
15. $\left(1+3 x^{2}\right) y^{\prime \prime}-2 x y^{\prime}+4 y$
16. Suppose $y(x)=\sum_{n=0}^{\infty} a_{n}(x+1)^{n}$ on an open interval that contains $x_{0}=-1$. Find a power series in $x+1$ for

$$
x y^{\prime \prime}+(4+2 x) y^{\prime}+(2+x) y
$$

17. Suppose $y(x)=\sum_{n=0}^{\infty} a_{n}(x-2)^{n}$ on an open interval that contains $x_{0}=2$. Find a power series in $x-2$ for

$$
x^{2} y^{\prime \prime}+2 x y^{\prime}-3 x y
$$

18. L Do the following experiment for various choices of real numbers $a_{0}$ and $a_{1}$.
(a) Use differential equations software to solve the initial value problem

$$
(2-x) y^{\prime \prime}+2 y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}
$$

numerically on $(-1.95,1.95)$. Choose the most accurate method your software package provides. (See Section 10.1 for a brief discussion of one such method.)
(b) For $N=2,3,4, \ldots$, compute $a_{2}, \ldots, a_{N}$ from Eqn.(7.1.18) and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on the same axes. Continue increasing $N$ until it's obvious that there's no point in continuing. (This sounds vague, but you'll know when to stop.)
19. L Follow the directions of Exercise 18 for the initial value problem

$$
(1+x) y^{\prime \prime}+2(x-1)^{2} y^{\prime}+3 y=0, \quad y(1)=a_{0}, \quad y^{\prime}(1)=a_{1}
$$

on the interval $(0,2)$. Use Eqns. (7.1.24) and (7.1.25) to compute $\left\{a_{n}\right\}$.
20. Suppose the series $\sum_{n=0}^{\infty} a_{n} x^{n}$ converges on an open interval $(-R, R)$, let $r$ be an arbitrary real number, and define

$$
y(x)=x^{r} \sum_{n=0}^{\infty} a_{n} x^{n}=\sum_{n=0}^{\infty} a_{n} x^{n+r}
$$

on $(0, R)$. Use Theorem 7.1.4 and the rule for differentiating the product of two functions to show that

$$
\begin{aligned}
y^{\prime}(x) & =\sum_{n=0}^{\infty}(n+r) a_{n} x^{n+r-1}, \\
y^{\prime \prime}(x) & =\sum_{n=0}^{\infty}(n+r)(n+r-1) a_{n} x^{n+r-2}, \\
& \vdots \\
y^{(k)}(x) & =\sum_{n=0}^{\infty}(n+r)(n+r-1) \cdots(n+r-k) a_{n} x^{n+r-k}
\end{aligned}
$$

on $(0, R)$
In Exercises 21-26 let y be as defined in Exercise 20, and write the given expression in the form $x^{r} \sum_{n=0}^{\infty} b_{n} x^{n}$.
21. $x^{2}(1-x) y^{\prime \prime}+x(4+x) y^{\prime}+(2-x) y$
22. $x^{2}(1+x) y^{\prime \prime}+x(1+2 x) y^{\prime}-(4+6 x) y$
23. $x^{2}(1+x) y^{\prime \prime}-x\left(1-6 x-x^{2}\right) y^{\prime}+\left(1+6 x+x^{2}\right) y$
24. $x^{2}(1+3 x) y^{\prime \prime}+x\left(2+12 x+x^{2}\right) y^{\prime}+2 x(3+x) y$
25. $x^{2}\left(1+2 x^{2}\right) y^{\prime \prime}+x\left(4+2 x^{2}\right) y^{\prime}+2\left(1-x^{2}\right) y$
26. $x^{2}\left(2+x^{2}\right) y^{\prime \prime}+2 x\left(5+x^{2}\right) y^{\prime}+2\left(3-x^{2}\right) y$

### 7.2 SERIES SOLUTIONS NEAR AN ORDINARY POINT I

Many physical applications give rise to second order homogeneous linear differential equations of the form

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0 \tag{7.2.1}
\end{equation*}
$$

where $P_{0}, P_{1}$, and $P_{2}$ are polynomials. Usually the solutions of these equations can't be expressed in terms of familiar elementary functions. Therefore we'll consider the problem of representing solutions of (7.2.1) with series.

We assume throughout that $P_{0}, P_{1}$ and $P_{2}$ have no common factors. Then we say that $x_{0}$ is an ordinary point of (7.2.1) if $P_{0}\left(x_{0}\right) \neq 0$, or a singular point if $P_{0}\left(x_{0}\right)=0$. For Legendre's equation,

$$
\begin{equation*}
\left(1-x^{2}\right) y^{\prime \prime}-2 x y^{\prime}+\alpha(\alpha+1) y=0, \tag{7.2.2}
\end{equation*}
$$

$x_{0}=1$ and $x_{0}=-1$ are singular points and all other points are ordinary points. For Bessel's equation,

$$
x^{2} y^{\prime \prime}+x y^{\prime}+\left(x^{2}-\nu^{2}\right) y=0
$$

$x_{0}=0$ is a singular point and all other points are ordinary points. If $P_{0}$ is a nonzero constant as in Airy's equation,

$$
\begin{equation*}
y^{\prime \prime}-x y=0, \tag{7.2.3}
\end{equation*}
$$

then every point is an ordinary point.
Since polynomials are continuous everywhere, $P_{1} / P_{0}$ and $P_{2} / P_{0}$ are continuous at any point $x_{0}$ that isn't a zero of $P_{0}$. Therefore, if $x_{0}$ is an ordinary point of (7.2.1) and $a_{0}$ and $a_{1}$ are arbitrary real numbers, then the initial value problem

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0, \quad y\left(x_{0}\right)=a_{0}, \quad y^{\prime}\left(x_{0}\right)=a_{1} \tag{7.2.4}
\end{equation*}
$$

has a unique solution on the largest open interval that contains $x_{0}$ and does not contain any zeros of $P_{0}$. To see this, we rewrite the differential equation in (7.2.4) as

$$
y^{\prime \prime}+\frac{P_{1}(x)}{P_{0}(x)} y^{\prime}+\frac{P_{2}(x)}{P_{0}(x)} y=0
$$

and apply Theorem 5.1.1 with $p=P_{1} / P_{0}$ and $q=P_{2} / P_{0}$. In this section and the next we consider the problem of representing solutions of (7.2.1) by power series that converge for values of $x$ near an ordinary point $x_{0}$.

We state the next theorem without proof.
Theorem 7.2.1 Suppose $P_{0}, P_{1}$, and $P_{2}$ are polynomials with no common factor and $P_{0}$ isn't identically zero. Let $x_{0}$ be a point such that $P_{0}\left(x_{0}\right) \neq 0$, and let $\rho$ be the distance from $x_{0}$ to the nearest zero of $P_{0}$ in the complex plane. (If $P_{0}$ is constant, then $\rho=\infty$.) Then every solution of

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0 \tag{7.2.5}
\end{equation*}
$$

can be represented by a power series

$$
\begin{equation*}
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n} \tag{7.2.6}
\end{equation*}
$$

that converges at least on the open interval $\left(x_{0}-\rho, x_{0}+\rho\right)$. (If $P_{0}$ is nonconstant, so that $\rho$ is necessarily finite, then the open interval of convergence of (7.2.6) may be larger than $\left(x_{0}-\rho, x_{0}+\rho\right)$. If $P_{0}$ is constant then $\rho=\infty$ and $\left(x_{0}-\rho, x_{0}+\rho\right)=(-\infty, \infty)$.)

We call (7.2.6) a power series solution in $x-x_{0}$ of (7.2.5). We'll now develop a method for finding power series solutions of (7.2.5). For this purpose we write (7.2.5) as $L y=0$, where

$$
\begin{equation*}
L y=P_{0} y^{\prime \prime}+P_{1} y^{\prime}+P_{2} y . \tag{7.2.7}
\end{equation*}
$$

Theorem 7.2.1 implies that every solution of $L y=0$ on $\left(x_{0}-\rho, x_{0}+\rho\right)$ can be written as

$$
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n} .
$$

Setting $x=x_{0}$ in this series and in the series

$$
y^{\prime}=\sum_{n=1}^{\infty} n a_{n}\left(x-x_{0}\right)^{n-1}
$$

shows that $y\left(x_{0}\right)=a_{0}$ and $y^{\prime}\left(x_{0}\right)=a_{1}$. Since every initial value problem (7.2.4) has a unique solution, this means that $a_{0}$ and $a_{1}$ can be chosen arbitrarily, and $a_{2}, a_{3}, \ldots$ are uniquely determined by them.

To find $a_{2}, a_{3}, \ldots$, we write $P_{0}, P_{1}$, and $P_{2}$ in powers of $x-x_{0}$, substitute

$$
\begin{gathered}
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}, \\
y^{\prime}=\sum_{n=1}^{\infty} n a_{n}\left(x-x_{0}\right)^{n-1}, \\
y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}
\end{gathered}
$$

into (7.2.7), and collect the coefficients of like powers of $x-x_{0}$. This yields

$$
\begin{equation*}
L y=\sum_{n=0}^{\infty} b_{n}\left(x-x_{0}\right)^{n}, \tag{7.2.8}
\end{equation*}
$$

where $\left\{b_{0}, b_{1}, \ldots, b_{n}, \ldots\right\}$ are expressed in terms of $\left\{a_{0}, a_{1}, \ldots, a_{n}, \ldots\right\}$ and the coefficients of $P_{0}, P_{1}$, and $P_{2}$, written in powers of $x-x_{0}$. Since (7.2.8) and (a) of Theorem 7.1.6 imply that $L y=0$ if and only if $b_{n}=0$ for $n \geq 0$, all power series solutions in $x-x_{0}$ of $L y=0$ can be obtained by choosing $a_{0}$ and $a_{1}$ arbitrarily and computing $a_{2}, a_{3}, \ldots$, successively so that $b_{n}=0$ for $n \geq 0$. For simplicity, we call the power series obtained this way the power series in $x-x_{0}$ for the general solution of $L y=0$, without explicitly identifying the open interval of convergence of the series.

Example 7.2.1 Let $x_{0}$ be an arbitrary real number. Find the power series in $x-x_{0}$ for the general solution of

$$
\begin{equation*}
y^{\prime \prime}+y=0 \tag{7.2.9}
\end{equation*}
$$

## Solution Here

$$
L y=y^{\prime \prime}+y .
$$

If

$$
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

then

$$
y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2},
$$

so

$$
L y=\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}+\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n} .
$$

To collect coefficients of like powers of $x-x_{0}$, we shift the summation index in the first sum. This yields

$$
L y=\sum_{n=0}^{\infty}(n+2)(n+1) a_{n+2}\left(x-x_{0}\right)^{n}+\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}=\sum_{n=0}^{\infty} b_{n}\left(x-x_{0}\right)^{n}
$$

with

$$
b_{n}=(n+2)(n+1) a_{n+2}+a_{n} .
$$

Therefore $L y=0$ if and only if

$$
\begin{equation*}
a_{n+2}=\frac{-a_{n}}{(n+2)(n+1)}, \quad n \geq 0 \tag{7.2.10}
\end{equation*}
$$

where $a_{0}$ and $a_{1}$ are arbitrary. Since the indices on the left and right sides of (7.2.10) differ by two, we write (7.2.10) separately for $n$ even $(n=2 m)$ and $n$ odd $(n=2 m+1)$. This yields

$$
\begin{equation*}
a_{2 m+2}=\frac{-a_{2 m}}{(2 m+2)(2 m+1)}, \quad m \geq 0, \tag{7.2.11}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{2 m+3}=\frac{-a_{2 m+1}}{(2 m+3)(2 m+2)}, \quad m \geq 0 \tag{7.2.12}
\end{equation*}
$$

Computing the coefficients of the even powers of $x-x_{0}$ from (7.2.11) yields

$$
\begin{aligned}
& a_{2}=-\frac{a_{0}}{2 \cdot 1} \\
& a_{4}=-\frac{a_{2}}{4 \cdot 3}=-\frac{1}{4 \cdot 3}\left(-\frac{a_{0}}{2 \cdot 1}\right)=\frac{a_{0}}{4 \cdot 3 \cdot 2 \cdot 1}, \\
& a_{6}=-\frac{a_{4}}{6 \cdot 5}=-\frac{1}{6 \cdot 5}\left(\frac{a_{0}}{4 \cdot 3 \cdot 2 \cdot 1}\right)=-\frac{a_{0}}{6 \cdot 5 \cdot 4 \cdot 3 \cdot 2 \cdot 1},
\end{aligned}
$$

and, in general,

$$
\begin{equation*}
a_{2 m}=(-1)^{m} \frac{a_{0}}{(2 m)!}, \quad m \geq 0 \tag{7.2.13}
\end{equation*}
$$

Computing the coefficients of the odd powers of $x-x_{0}$ from (7.2.12) yields

$$
\begin{aligned}
& a_{3}=-\frac{a_{1}}{3 \cdot 2} \\
& a_{5}=-\frac{a_{3}}{5 \cdot 4}=-\frac{1}{5 \cdot 4}\left(-\frac{a_{1}}{3 \cdot 2}\right)=\frac{a_{1}}{5 \cdot 4 \cdot 3 \cdot 2}, \\
& a_{7}=-\frac{a_{5}}{7 \cdot 6}=-\frac{1}{7 \cdot 6}\left(\frac{a_{1}}{5 \cdot 4 \cdot 3 \cdot 2}\right)=-\frac{a_{1}}{7 \cdot 6 \cdot 5 \cdot 4 \cdot 3 \cdot 2}
\end{aligned}
$$

and, in general,

$$
\begin{equation*}
a_{2 m+1}=\frac{(-1)^{m} a_{1}}{(2 m+1)!} \quad m \geq 0 \tag{7.2.14}
\end{equation*}
$$

Thus, the general solution of (7.2.9) can be written as

$$
y=\sum_{m=0}^{\infty} a_{2 m}\left(x-x_{0}\right)^{2 m}+\sum_{m=0}^{\infty} a_{2 m+1}\left(x-x_{0}\right)^{2 m+1}
$$

or, from (7.2.13) and (7.2.14), as

$$
\begin{equation*}
y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{\left(x-x_{0}\right)^{2 m}}{(2 m)!}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{\left(x-x_{0}\right)^{2 m+1}}{(2 m+1)!} . \tag{7.2.15}
\end{equation*}
$$

If we recall from calculus that

$$
\sum_{m=0}^{\infty}(-1)^{m} \frac{\left(x-x_{0}\right)^{2 m}}{(2 m)!}=\cos \left(x-x_{0}\right) \quad \text { and } \quad \sum_{m=0}^{\infty}(-1)^{m} \frac{\left(x-x_{0}\right)^{2 m+1}}{(2 m+1)!}=\sin \left(x-x_{0}\right)
$$

then (7.2.15) becomes

$$
y=a_{0} \cos \left(x-x_{0}\right)+a_{1} \sin \left(x-x_{0}\right)
$$

which should look familiar.
Equations like (7.2.10), (7.2.11), and (7.2.12), which define a given coefficient in the sequence $\left\{a_{n}\right\}$ in terms of one or more coefficients with lesser indices are called recurrence relations. When we use a recurrence relation to compute terms of a sequence we're computing recursively.

In the remainder of this section we consider the problem of finding power series solutions in $x-x_{0}$ for equations of the form

$$
\begin{equation*}
\left(1+\alpha\left(x-x_{0}\right)^{2}\right) y^{\prime \prime}+\beta\left(x-x_{0}\right) y^{\prime}+\gamma y=0 \tag{7.2.16}
\end{equation*}
$$

Many important equations that arise in applications are of this form with $x_{0}=0$, including Legendre's equation (7.2.2), Airy's equation (7.2.3), Chebyshev's equation,

$$
\left(1-x^{2}\right) y^{\prime \prime}-x y^{\prime}+\alpha^{2} y=0
$$

and Hermite's equation,

$$
y^{\prime \prime}-2 x y^{\prime}+2 \alpha y=0
$$

Since

$$
P_{0}(x)=1+\alpha\left(x-x_{0}\right)^{2}
$$

in (7.2.16), the point $x_{0}$ is an ordinary point of (7.2.16), and Theorem 7.2.1 implies that the solutions of (7.2.16) can be written as power series in $x-x_{0}$ that converge on the interval ( $\left.x_{0}-1 / \sqrt{\mid} \alpha\left|, x_{0}+1 / \sqrt{ }\right| \alpha \mid\right)$ if $\alpha \neq 0$, or on $(-\infty, \infty)$ if $\alpha=0$. We'll see that the coefficients in these power series can be obtained by methods similar to the one used in Example 7.2.1.

To simplify finding the coefficients, we introduce some notation for products:

$$
\prod_{j=r}^{s} b_{j}=b_{r} b_{r+1} \cdots b_{s} \quad \text { if } \quad s \geq r
$$

Thus,

$$
\begin{gathered}
\prod_{j=2}^{7} b_{j}=b_{2} b_{3} b_{4} b_{5} b_{6} b_{7} \\
\prod_{j=0}^{4}(2 j+1)=(1)(3)(5)(7)(9)=945
\end{gathered}
$$

and

$$
\prod_{j=2}^{2} j^{2}=2^{2}=4
$$

We define

$$
\prod_{j=r}^{s} b_{j}=1 \quad \text { if } \quad s<r
$$

no matter what the form of $b_{j}$.
Example 7.2.2 Find the power series in $x$ for the general solution of

$$
\begin{equation*}
\left(1+2 x^{2}\right) y^{\prime \prime}+6 x y^{\prime}+2 y=0 \tag{7.2.17}
\end{equation*}
$$

## Solution Here

$$
L y=\left(1+2 x^{2}\right) y^{\prime \prime}+6 x y^{\prime}+2 y
$$

If

$$
y=\sum_{n=0}^{\infty} a_{n} x^{n}
$$

then

$$
y^{\prime}=\sum_{n=1}^{\infty} n a_{n} x^{n-1} \quad \text { and } \quad y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}
$$

so

$$
\begin{aligned}
L y & =\left(1+2 x^{2}\right) \sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}+6 x \sum_{n=1}^{\infty} n a_{n} x^{n-1}+2 \sum_{n=0}^{\infty} a_{n} x^{n} \\
& =\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}+\sum_{n=0}^{\infty}[2 n(n-1)+6 n+2] a_{n} x^{n} \\
& =\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}+2 \sum_{n=0}^{\infty}(n+1)^{2} a_{n} x^{n} .
\end{aligned}
$$

To collect coefficients of $x^{n}$, we shift the summation index in the first sum. This yields

$$
L y=\sum_{n=0}^{\infty}(n+2)(n+1) a_{n+2} x^{n}+2 \sum_{n=0}^{\infty}(n+1)^{2} a_{n} x^{n}=\sum_{n=0}^{\infty} b_{n} x^{n}
$$

with

$$
b_{n}=(n+2)(n+1) a_{n+2}+2(n+1)^{2} a_{n}, \quad n \geq 0
$$

To obtain solutions of (7.2.17), we set $b_{n}=0$ for $n \geq 0$. This is equivalent to the recurrence relation

$$
\begin{equation*}
a_{n+2}=-2 \frac{n+1}{n+2} a_{n}, \quad n \geq 0 . \tag{7.2.18}
\end{equation*}
$$

Since the indices on the left and right differ by two, we write (7.2.18) separately for $n=2 m$ and $n=2 m+1$, as in Example 7.2.1. This yields

$$
\begin{equation*}
a_{2 m+2}=-2 \frac{2 m+1}{2 m+2} a_{2 m}=-\frac{2 m+1}{m+1} a_{2 m}, \quad m \geq 0 \tag{7.2.19}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{2 m+3}=-2 \frac{2 m+2}{2 m+3} a_{2 m+1}=-4 \frac{m+1}{2 m+3} a_{2 m+1}, \quad m \geq 0 . \tag{7.2.20}
\end{equation*}
$$

Computing the coefficients of even powers of $x$ from (7.2.19) yields

$$
\begin{aligned}
& a_{2}=-\frac{1}{1} a_{0}, \\
& a_{4}=-\frac{3}{2} a_{2}=\left(-\frac{3}{2}\right)\left(-\frac{1}{1}\right) a_{0}=\frac{1 \cdot 3}{1 \cdot 2} a_{0}, \\
& a_{6}=-\frac{5}{3} a_{4}=-\frac{5}{3}\left(\frac{1 \cdot 3}{1 \cdot 2}\right) a_{0}=-\frac{1 \cdot 3 \cdot 5}{1 \cdot 2 \cdot 3} a_{0}, \\
& a_{8}=-\frac{7}{4} a_{6}=-\frac{7}{4}\left(-\frac{1 \cdot 3 \cdot 5}{1 \cdot 2 \cdot 3}\right) a_{0}=\frac{1 \cdot 3 \cdot 5 \cdot 7}{1 \cdot 2 \cdot 3 \cdot 4} a_{0} .
\end{aligned}
$$

In general,

$$
\begin{equation*}
a_{2 m}=(-1)^{m} \frac{\prod_{j=1}^{m}(2 j-1)}{m!} a_{0}, \quad m \geq 0 . \tag{7.2.21}
\end{equation*}
$$

(Note that (7.2.21) is correct for $m=0$ because we defined $\prod_{j=1}^{0} b_{j}=1$ for any $b_{j}$.)
Computing the coefficients of odd powers of $x$ from (7.2.20) yields

$$
\begin{aligned}
& a_{3}=-4 \frac{1}{3} a_{1}, \\
& a_{5}=-4 \frac{2}{5} a_{3}=-4 \frac{2}{5}\left(-4 \frac{1}{3}\right) a_{1}=4^{2} \frac{1 \cdot 2}{3 \cdot 5} a_{1}, \\
& a_{7}=-4 \frac{3}{7} a_{5}=-4 \frac{3}{7}\left(4^{2} \frac{1 \cdot 2}{3 \cdot 5}\right) a_{1}=-4^{3} \frac{1 \cdot 2 \cdot 3}{3 \cdot 5 \cdot 7} a_{1}, \\
& a_{9}=-4 \frac{4}{9} a_{7}=-4 \frac{4}{9}\left(4^{3} \frac{1 \cdot 2 \cdot 3}{3 \cdot 5 \cdot 7}\right) a_{1}=4^{4} \frac{1 \cdot 2 \cdot 3 \cdot 4}{3 \cdot 5 \cdot 7 \cdot 9} a_{1} .
\end{aligned}
$$

In general,

$$
\begin{equation*}
a_{2 m+1}=\frac{(-1)^{m} 4^{m} m!}{\prod_{j=1}^{m}(2 j+1)} a_{1}, \quad m \geq 0 \tag{7.2.22}
\end{equation*}
$$

From (7.2.21) and (7.2.22),

$$
y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{\prod_{j=1}^{m}(2 j-1)}{m!} x^{2 m}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{4^{m} m!}{\prod_{j=1}^{m}(2 j+1)} x^{2 m+1} .
$$

is the power series in $x$ for the general solution of (7.2.17). Since $P_{0}(x)=1+2 x^{2}$ has no real zeros, Theorem 5.1.1 implies that every solution of (7.2.17) is defined on $(-\infty, \infty)$. However, since $P_{0}( \pm i / \sqrt{2})=0$, Theorem 7.2.1 implies only that the power series converges in $(-1 / \sqrt{2}, 1 / \sqrt{2})$ for any choice of $a_{0}$ and $a_{1}$.

The results in Examples 7.2.1 and 7.2.2 are consequences of the following general theorem.
Theorem 7.2.2 The coefficients $\left\{a_{n}\right\}$ in any solution $y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}$ of

$$
\begin{equation*}
\left(1+\alpha\left(x-x_{0}\right)^{2}\right) y^{\prime \prime}+\beta\left(x-x_{0}\right) y^{\prime}+\gamma y=0 \tag{7.2.23}
\end{equation*}
$$

satisfy the recurrence relation

$$
\begin{equation*}
a_{n+2}=-\frac{p(n)}{(n+2)(n+1)} a_{n}, \quad n \geq 0 \tag{7.2.24}
\end{equation*}
$$

where

$$
\begin{equation*}
p(n)=\alpha n(n-1)+\beta n+\gamma \tag{7.2.25}
\end{equation*}
$$

Moreover, the coefficients of the even and odd powers of $x-x_{0}$ can be computed separately as

$$
\begin{equation*}
a_{2 m+2}=-\frac{p(2 m)}{(2 m+2)(2 m+1)} a_{2 m}, \quad m \geq 0 \tag{7.2.26}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{2 m+3}=-\frac{p(2 m+1)}{(2 m+3)(2 m+2)} a_{2 m+1}, \quad m \geq 0 \tag{7.2.27}
\end{equation*}
$$

where $a_{0}$ and $a_{1}$ are arbitrary.
Proof Here

$$
L y=\left(1+\alpha\left(x-x_{0}\right)^{2}\right) y^{\prime \prime}+\beta\left(x-x_{0}\right) y^{\prime}+\gamma y
$$

If

$$
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

then

$$
y^{\prime}=\sum_{n=1}^{\infty} n a_{n}\left(x-x_{0}\right)^{n-1} \quad \text { and } \quad y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}
$$

Hence,

$$
\begin{aligned}
L y & =\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}+\sum_{n=0}^{\infty}[\alpha n(n-1)+\beta n+\gamma] a_{n}\left(x-x_{0}\right)^{n} \\
& =\sum_{n=2}^{\infty} n(n-1) a_{n}\left(x-x_{0}\right)^{n-2}+\sum_{n=0}^{\infty} p(n) a_{n}\left(x-x_{0}\right)^{n}
\end{aligned}
$$

from (7.2.25). To collect coefficients of powers of $x-x_{0}$, we shift the summation index in the first sum. This yields

$$
L y=\sum_{n=0}^{\infty}\left[(n+2)(n+1) a_{n+2}+p(n) a_{n}\right]\left(x-x_{0}\right)^{n}
$$

Thus, $L y=0$ if and only if

$$
(n+2)(n+1) a_{n+2}+p(n) a_{n}=0, \quad n \geq 0
$$

which is equivalent to (7.2.24). Writing (7.2.24) separately for the cases where $n=2 m$ and $n=2 m+1$ yields (7.2.26) and (7.2.27).

Example 7.2.3 Find the power series in $x-1$ for the general solution of

$$
\begin{equation*}
\left(2+4 x-2 x^{2}\right) y^{\prime \prime}-12(x-1) y^{\prime}-12 y=0 \tag{7.2.28}
\end{equation*}
$$

Solution We must first write the coefficient $P_{0}(x)=2+4 x-x^{2}$ in powers of $x-1$. To do this, we write $x=(x-1)+1$ in $P_{0}(x)$ and then expand the terms, collecting powers of $x-1$; thus,

$$
\begin{aligned}
2+4 x-2 x^{2} & =2+4[(x-1)+1]-2[(x-1)+1]^{2} \\
& =4-2(x-1)^{2}
\end{aligned}
$$

Therefore we can rewrite (7.2.28) as

$$
\left(4-2(x-1)^{2}\right) y^{\prime \prime}-12(x-1) y^{\prime}-12 y=0
$$

or, equivalently,

$$
\left(1-\frac{1}{2}(x-1)^{2}\right) y^{\prime \prime}-3(x-1) y^{\prime}-3 y=0
$$

This is of the form (7.2.23) with $\alpha=-1 / 2, \beta=-3$, and $\gamma=-3$. Therefore, from (7.2.25)

$$
p(n)=-\frac{n(n-1)}{2}-3 n-3=-\frac{(n+2)(n+3)}{2}
$$

Hence, Theorem 7.2.2 implies that

$$
\begin{aligned}
a_{2 m+2} & =-\frac{p(2 m)}{(2 m+2)(2 m+1)} a_{2 m} \\
& =\frac{(2 m+2)(2 m+3)}{2(2 m+2)(2 m+1)} a_{2 m}=\frac{2 m+3}{2(2 m+1)} a_{2 m}, \quad m \geq 0
\end{aligned}
$$

and

$$
\begin{aligned}
a_{2 m+3} & =-\frac{p(2 m+1)}{(2 m+3)(2 m+2)} a_{2 m+1} \\
& =\frac{(2 m+3)(2 m+4)}{2(2 m+3)(2 m+2)} a_{2 m+1}=\frac{m+2}{2(m+1)} a_{2 m+1}, \quad m \geq 0
\end{aligned}
$$

We leave it to you to show that

$$
a_{2 m}=\frac{2 m+1}{2^{m}} a_{0} \quad \text { and } \quad a_{2 m+1}=\frac{m+1}{2^{m}} a_{1}, \quad m \geq 0
$$

which implies that the power series in $x-1$ for the general solution of (7.2.28) is

$$
y=a_{0} \sum_{m=0}^{\infty} \frac{2 m+1}{2^{m}}(x-1)^{2 m}+a_{1} \sum_{m=0}^{\infty} \frac{m+1}{2^{m}}(x-1)^{2 m+1} .
$$

In the examples considered so far we were able to obtain closed formulas for coefficients in the power series solutions. In some cases this is impossible, and we must settle for computing a finite number of terms in the series. The next example illustrates this with an initial value problem.

Example 7.2.4 Compute $a_{0}, a_{1}, \ldots, a_{7}$ in the series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of the initial value problem

$$
\begin{equation*}
\left(1+2 x^{2}\right) y^{\prime \prime}+10 x y^{\prime}+8 y=0, \quad y(0)=2, \quad y^{\prime}(0)=-3 \tag{7.2.29}
\end{equation*}
$$

Solution Since $\alpha=2, \beta=10$, and $\gamma=8$ in (7.2.29),

$$
p(n)=2 n(n-1)+10 n+8=2(n+2)^{2}
$$

Therefore

$$
a_{n+2}=-2 \frac{(n+2)^{2}}{(n+2)(n+1)} a_{n}=-2 \frac{n+2}{n+1} a_{n}, \quad n \geq 0
$$

Writing this equation separately for $n=2 m$ and $n=2 m+1$ yields

$$
\begin{equation*}
a_{2 m+2}=-2 \frac{(2 m+2)}{2 m+1} a_{2 m}=-4 \frac{m+1}{2 m+1} a_{2 m}, \quad m \geq 0 \tag{7.2.30}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{2 m+3}=-2 \frac{2 m+3}{2 m+2} a_{2 m+1}=-\frac{2 m+3}{m+1} a_{2 m+1}, \quad m \geq 0 . \tag{7.2.31}
\end{equation*}
$$

Starting with $a_{0}=y(0)=2$, we compute $a_{2}, a_{4}$, and $a_{6}$ from (7.2.30):

$$
\begin{aligned}
& a_{2}=-4 \frac{1}{1} 2=-8, \\
& a_{4}=-4 \frac{2}{3}(-8)=\frac{64}{3}, \\
& a_{6}=-4 \frac{3}{5}\left(\frac{64}{3}\right)=-\frac{256}{5} .
\end{aligned}
$$

Starting with $a_{1}=y^{\prime}(0)=-3$, we compute $a_{3}, a_{5}$ and $a_{7}$ from (7.2.31):

$$
\begin{aligned}
& a_{3}=-\frac{3}{1}(-3)=9, \\
& a_{5}=-\frac{5}{2} 9=-\frac{45}{2}, \\
& a_{7}=-\frac{7}{3}\left(-\frac{45}{2}\right)=\frac{105}{2} .
\end{aligned}
$$

Therefore the solution of (7.2.29) is

$$
y=2-3 x-8 x^{2}+9 x^{3}+\frac{64}{3} x^{4}-\frac{45}{2} x^{5}-\frac{256}{5} x^{6}+\frac{105}{2} x^{7}+\cdots .
$$

## USING TECHNOLOGY

Computing coefficients recursively as in Example 7.2.4 is tedious. We recommend that you do this kind of computation by writing a short program to implement the appropriate recurrence relation on a calculator or computer. You may wish to do this in verifying examples and doing exercises (identified by the symbol $\boxed{\mathrm{C}}$ ) in this chapter that call for numerical computation of the coefficients in series solutions. We obtained the answers to these exercises by using software that can produce answers in the form of rational numbers. However, it's perfectly acceptable - and more practical - to get your answers in decimal form. You can always check them by converting our fractions to decimals.

If you're interested in actually using series to compute numerical approximations to solutions of a differential equation, then whether or not there's a simple closed form for the coefficents is essentially irrelevant. For computational purposes it's usually more efficient to start with the given coefficients $a_{0}=y\left(x_{0}\right)$ and $a_{1}=y^{\prime}\left(x_{0}\right)$, compute $a_{2}, \ldots, a_{N}$ recursively, and then compute approximate values of the solution from the Taylor polynomial

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n}\left(x-x_{0}\right)^{n} .
$$

The trick is to decide how to choose $N$ so the approximation $y(x) \approx T_{N}(x)$ is sufficiently accurate on the subinterval of the interval of convergence that you're interested in. In the computational exercises in this and the next two sections, you will often be asked to obtain the solution of a given problem by numerical integration with software of your choice (see Section 10.1 for a brief discussion of one such method), and to compare the solution obtained in this way with the approximations obtained with $T_{N}$ for various values of $N$. This is a typical textbook kind of exercise, designed to give you insight into how the accuracy of the approximation $y(x) \approx T_{N}(x)$ behaves as a function of $N$ and the interval that you're working on. In real life, you would choose one or the other of the two methods (numerical integration or series solution). If you choose the method of series solution, then a practical procedure for determining a suitable value of $N$ is to continue increasing $N$ until the maximum of $\left|T_{N}-T_{N-1}\right|$ on the interval of interest is within the margin of error that you're willing to accept.

In doing computational problems that call for numerical solution of differential equations you should choose the most accurate numerical integration procedure your software supports, and experiment with
the step size until you're confident that the numerical results are sufficiently accurate for the problem at hand.

### 7.2 Exercises

In Exercises 1-8 find the power series in $x$ for the general solution.

1. $\left(1+x^{2}\right) y^{\prime \prime}+6 x y^{\prime}+6 y=0$
2. $\left(1+x^{2}\right) y^{\prime \prime}+2 x y^{\prime}-2 y=0$
3. $\left(1+x^{2}\right) y^{\prime \prime}-8 x y^{\prime}+20 y=0$
4. $\left(1+2 x^{2}\right) y^{\prime \prime}+7 x y^{\prime}+2 y=0$
5. $\left(1-x^{2}\right) y^{\prime \prime}-5 x y^{\prime}-4 y=0$
6. $\left(1-x^{2}\right) y^{\prime \prime}-8 x y^{\prime}-12 y=0$
7. $\left(1+x^{2}\right) y^{\prime \prime}+2 x y^{\prime}+\frac{1}{4} y=0$
8. $\left(1+x^{2}\right) y^{\prime \prime}-10 x y^{\prime}+28 y=0$
9. L
(a) Find the power series in $x$ for the general solution of $y^{\prime \prime}+x y^{\prime}+2 y=0$.
(b) For several choices of $a_{0}$ and $a_{1}$, use differential equations software to solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+x y^{\prime}+2 y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}, \tag{A}
\end{equation*}
$$

numerically on $(-5,5)$.
(c) For fixed $r$ in $\{1,2,3,4,5\}$ graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on $(-r, r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs.
10. $L$ Follow the directions of Exercise 9 for the differential equation

$$
y^{\prime \prime}+2 x y^{\prime}+3 y=0 .
$$

In Exercises 11-13 find $a_{0}, \ldots, a_{N}$ for $N$ at least 7 in the power series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of the initial value problem.
11. C $\left(1+x^{2}\right) y^{\prime \prime}+x y^{\prime}+y=0, \quad y(0)=2, \quad y^{\prime}(0)=-1$
12. $\mathrm{C}\left(1+2 x^{2}\right) y^{\prime \prime}-9 x y^{\prime}-6 y=0, \quad y(0)=1, \quad y^{\prime}(0)=-1$
13. $\mathrm{C}\left(1+8 x^{2}\right) y^{\prime \prime}+2 y=0, \quad y(0)=2, \quad y^{\prime}(0)=-1$
14. L Do the next experiment for various choices of real numbers $a_{0}$, $a_{1}$, and $r$, with $0<r<1 / \sqrt{2}$.
(a) Use differential equations software to solve the initial value problem

$$
\begin{equation*}
\left(1-2 x^{2}\right) y^{\prime \prime}-x y^{\prime}+3 y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}, \tag{A}
\end{equation*}
$$

numerically on $(-r, r)$.
(b) For $N=2,3,4, \ldots$, compute $a_{2}, \ldots, a_{N}$ in the power series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of (A), and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on $(-r, r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs.
15. L Do (a) and (b) for several values of $r$ in $(0,1)$ :
(a) Use differential equations software to solve the initial value problem

$$
\begin{equation*}
\left(1+x^{2}\right) y^{\prime \prime}+10 x y^{\prime}+14 y=0, \quad y(0)=5, \quad y^{\prime}(0)=1 \tag{A}
\end{equation*}
$$

numerically on $(-r, r)$.
(b) For $N=2,3,4, \ldots$, compute $a_{2}, \ldots, a_{N}$ in the power series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of (A), and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on $(-r, r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs. What happens to the required $N$ as $r \rightarrow 1$ ?
(c) Try (a) and (b) with $r=1.2$. Explain your results.

In Exercises 16-20 find the power series in $x-x_{0}$ for the general solution.
16. $y^{\prime \prime}-y=0 ; \quad x_{0}=3$ 17. $\quad y^{\prime \prime}-(x-3) y^{\prime}-y=0 ; \quad x_{0}=3$
18. $\left(1-4 x+2 x^{2}\right) y^{\prime \prime}+10(x-1) y^{\prime}+6 y=0 ; \quad x_{0}=1$
19. $\left(11-8 x+2 x^{2}\right) y^{\prime \prime}-16(x-2) y^{\prime}+36 y=0 ; \quad x_{0}=2$
20. $\left(5+6 x+3 x^{2}\right) y^{\prime \prime}+9(x+1) y^{\prime}+3 y=0 ; \quad x_{0}=-1$

In Exercises 21-26 find $a_{0}, \ldots, a_{N}$ for $N$ at least 7 in the power series $y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}$ for the solution of the initial value problem. Take $x_{0}$ to be the point where the initial conditions are imposed.
21. $\mathrm{C}\left(x^{2}-4\right) y^{\prime \prime}-x y^{\prime}-3 y=0, \quad y(0)=-1, \quad y^{\prime}(0)=2$
22. $\mathrm{C} y^{\prime \prime}+(x-3) y^{\prime}+3 y=0, \quad y(3)=-2, \quad y^{\prime}(3)=3$
23. $\mathrm{C}\left(5-6 x+3 x^{2}\right) y^{\prime \prime}+(x-1) y^{\prime}+12 y=0, \quad y(1)=-1, \quad y^{\prime}(1)=1$
24. $\mathrm{C}\left(4 x^{2}-24 x+37\right) y^{\prime \prime}+y=0, \quad y(3)=4, \quad y^{\prime}(3)=-6$
25. C $\left(x^{2}-8 x+14\right) y^{\prime \prime}-8(x-4) y^{\prime}+20 y=0, \quad y(4)=3, \quad y^{\prime}(4)=-4$
26. C $\left(2 x^{2}+4 x+5\right) y^{\prime \prime}-20(x+1) y^{\prime}+60 y=0, \quad y(-1)=3, \quad y^{\prime}(-1)=-3$
27. (a) Find a power series in $x$ for the general solution of

$$
\begin{equation*}
\left(1+x^{2}\right) y^{\prime \prime}+4 x y^{\prime}+2 y=0 \tag{A}
\end{equation*}
$$

(b) Use (a) and the formula

$$
\frac{1}{1-r}=1+r+r^{2}+\cdots+r^{n}+\cdots \quad(-1<r<1)
$$

for the sum of a geometric series to find a closed form expression for the general solution of (A) on $(-1,1)$.
(c) Show that the expression obtained in (b) is actually the general solution of of (A) on $(-\infty, \infty)$.
28. Use Theorem 7.2.2 to show that the power series in $x$ for the general solution of

$$
\left(1+\alpha x^{2}\right) y^{\prime \prime}+\beta x y^{\prime}+\gamma y=0
$$

is

$$
y=a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} p(2 j)\right] \frac{x^{2 m}}{(2 m)!}+a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} p(2 j+1)\right] \frac{x^{2 m+1}}{(2 m+1)!}
$$

29. Use Exercise 28 to show that all solutions of

$$
\left(1+\alpha x^{2}\right) y^{\prime \prime}+\beta x y^{\prime}+\gamma y=0
$$

are polynomials if and only if

$$
\alpha n(n-1)+\beta n+\gamma=\alpha(n-2 r)(n-2 s-1)
$$

where $r$ and $s$ are nonnegative integers.
30. (a) Use Exercise 28 to show that the power series in $x$ for the general solution of

$$
\left(1-x^{2}\right) y^{\prime \prime}-2 b x y^{\prime}+\alpha(\alpha+2 b-1) y=0
$$

is $y=a_{0} y_{1}+a_{1} y_{2}$, where

$$
y_{1}=\sum_{m=0}^{\infty}\left[\prod_{j=0}^{m-1}(2 j-\alpha)(2 j+\alpha+2 b-1)\right] \frac{x^{2 m}}{(2 m)!}
$$

and

$$
y_{2}=\sum_{m=0}^{\infty}\left[\prod_{j=0}^{m-1}(2 j+1-\alpha)(2 j+\alpha+2 b)\right] \frac{x^{2 m+1}}{(2 m+1)!} .
$$

(b) Suppose $2 b$ isn't a negative odd integer and $k$ is a nonnegative integer. Show that $y_{1}$ is a polynomial of degree $2 k$ such that $y_{1}(-x)=y_{1}(x)$ if $\alpha=2 k$, while $y_{2}$ is a polynomial of degree $2 k+1$ such that $y_{2}(-x)=-y_{2}(-x)$ if $\alpha=2 k+1$. Conclude that if $n$ is a nonnegative integer, then there's a polynomial $P_{n}$ of degree $n$ such that $P_{n}(-x)=(-1)^{n} P_{n}(x)$ and

$$
\begin{equation*}
\left(1-x^{2}\right) P_{n}^{\prime \prime}-2 b x P_{n}^{\prime}+n(n+2 b-1) P_{n}=0 \tag{A}
\end{equation*}
$$

(c) Show that (A) implies that

$$
\left[\left(1-x^{2}\right)^{b} P_{n}^{\prime}\right]^{\prime}=-n(n+2 b-1)\left(1-x^{2}\right)^{b-1} P_{n}
$$

and use this to show that if $m$ and $n$ are nonnegative integers, then

$$
\begin{align*}
& {\left[\left(1-x^{2}\right)^{b} P_{n}^{\prime}\right]^{\prime} P_{m}-\left[\left(1-x^{2}\right)^{b} P_{m}^{\prime}\right]^{\prime} P_{n}=} \\
& \quad[m(m+2 b-1)-n(n+2 b-1)]\left(1-x^{2}\right)^{b-1} P_{m} P_{n} \tag{B}
\end{align*}
$$

(d) Now suppose $b>0$. Use (B) and integration by parts to show that if $m \neq n$, then

$$
\int_{-1}^{1}\left(1-x^{2}\right)^{b-1} P_{m}(x) P_{n}(x) d x=0
$$

(We say that $P_{m}$ and $P_{n}$ are orthogonal on $(-1,1)$ with respect to the weighting function $\left(1-x^{2}\right)^{b-1}$.)
31. (a) Use Exercise 28 to show that the power series in $x$ for the general solution of Hermite's equation

$$
y^{\prime \prime}-2 x y^{\prime}+2 \alpha y=0
$$

is $y=a_{0} y_{1}+a_{1} y_{1}$, where

$$
y_{1}=\sum_{m=0}^{\infty}\left[\prod_{j=0}^{m-1}(2 j-\alpha)\right] \frac{2^{m} x^{2 m}}{(2 m)!}
$$

and

$$
y_{2}=\sum_{m=0}^{\infty}\left[\prod_{j=0}^{m-1}(2 j+1-\alpha)\right] \frac{2^{m} x^{2 m+1}}{(2 m+1)!} .
$$

(b) Suppose $k$ is a nonnegative integer. Show that $y_{1}$ is a polynomial of degree $2 k$ such that $y_{1}(-x)=y_{1}(x)$ if $\alpha=2 k$, while $y_{2}$ is a polynomial of degree $2 k+1$ such that $y_{2}(-x)=$ $-y_{2}(-x)$ if $\alpha=2 k+1$. Conclude that if $n$ is a nonnegative integer then there's a polynomial $P_{n}$ of degree $n$ such that $P_{n}(-x)=(-1)^{n} P_{n}(x)$ and

$$
\begin{equation*}
P_{n}^{\prime \prime}-2 x P_{n}^{\prime}+2 n P_{n}=0 \tag{A}
\end{equation*}
$$

(c) Show that (A) implies that

$$
\left[e^{-x^{2}} P_{n}^{\prime}\right]^{\prime}=-2 n e^{-x^{2}} P_{n}
$$

and use this to show that if $m$ and $n$ are nonnegative integers, then

$$
\begin{equation*}
\left[e^{-x^{2}} P_{n}^{\prime}\right]^{\prime} P_{m}-\left[e^{-x^{2}} P_{m}^{\prime}\right]^{\prime} P_{n}=2(m-n) e^{-x^{2}} P_{m} P_{n} \tag{B}
\end{equation*}
$$

(d) Use (B) and integration by parts to show that if $m \neq n$, then

$$
\int_{-\infty}^{\infty} e^{-x^{2}} P_{m}(x) P_{n}(x) d x=0
$$

(We say that $P_{m}$ and $P_{n}$ are orthogonal on $(-\infty, \infty)$ with respect to the weighting function $e^{-x^{2}}$.)
32. Consider the equation

$$
\begin{equation*}
\left(1+\alpha x^{3}\right) y^{\prime \prime}+\beta x^{2} y^{\prime}+\gamma x y=0 \tag{A}
\end{equation*}
$$

and let $p(n)=\alpha n(n-1)+\beta n+\gamma$. (The special case $y^{\prime \prime}-x y=0$ of (A) is Airy's equation.)
(a) Modify the argument used to prove Theorem 7.2.2 to show that

$$
y=\sum_{n=0}^{\infty} a_{n} x^{n}
$$

is a solution of (A) if and only if $a_{2}=0$ and

$$
a_{n+3}=-\frac{p(n)}{(n+3)(n+2)} a_{n}, \quad n \geq 0
$$

(b) Show from (a) that $a_{n}=0$ unless $n=3 m$ or $n=3 m+1$ for some nonnegative integer $m$, and that

$$
a_{3 m+3}=-\frac{p(3 m)}{(3 m+3)(3 m+2)} a_{3 m}, \quad m \geq 0
$$

and

$$
a_{3 m+4}=-\frac{p(3 m+1)}{(3 m+4)(3 m+3)} a_{3 m+1}, \quad m \geq 0
$$

where $a_{0}$ and $a_{1}$ may be specified arbitrarily.
(c) Conclude from (b) that the power series in $x$ for the general solution of (A) is

$$
\begin{aligned}
y= & a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{p(3 j)}{3 j+2}\right] \frac{x^{3 m}}{3^{m} m!} \\
& +a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{p(3 j+1)}{3 j+4}\right] \frac{x^{3 m+1}}{3^{m} m!}
\end{aligned}
$$

In Exercises 33-37 use the method of Exercise 32 to find the power series in $x$ for the general solution.
33. $y^{\prime \prime}-x y=0$
34. $\left(1-2 x^{3}\right) y^{\prime \prime}-10 x^{2} y^{\prime}-8 x y=0$
35. $\left(1+x^{3}\right) y^{\prime \prime}+7 x^{2} y^{\prime}+9 x y=0$
36. $\left(1-2 x^{3}\right) y^{\prime \prime}+6 x^{2} y^{\prime}+24 x y=0$
37. $\left(1-x^{3}\right) y^{\prime \prime}+15 x^{2} y^{\prime}-63 x y=0$
38. Consider the equation

$$
\begin{equation*}
\left(1+\alpha x^{k+2}\right) y^{\prime \prime}+\beta x^{k+1} y^{\prime}+\gamma x^{k} y=0 \tag{A}
\end{equation*}
$$

where $k$ is a positive integer, and let $p(n)=\alpha n(n-1)+\beta n+\gamma$.
(a) Modify the argument used to prove Theorem 7.2.2 to show that

$$
y=\sum_{n=0}^{\infty} a_{n} x^{n}
$$

is a solution of (A) if and only if $a_{n}=0$ for $2 \leq n \leq k+1$ and

$$
a_{n+k+2}=-\frac{p(n)}{(n+k+2)(n+k+1)} a_{n}, \quad n \geq 0
$$

(b) Show from (a) that $a_{n}=0$ unless $n=(k+2) m$ or $n=(k+2) m+1$ for some nonnegative integer $m$, and that

$$
a_{(k+2)(m+1)}=-\frac{p((k+2) m)}{(k+2)(m+1)[(k+2)(m+1)-1]} a_{(k+2) m}, \quad m \geq 0
$$

and

$$
a_{(k+2)(m+1)+1}=-\frac{p((k+2) m+1)}{[(k+2)(m+1)+1](k+2)(m+1)} a_{(k+2) m+1}, \quad m \geq 0
$$

where $a_{0}$ and $a_{1}$ may be specified arbitrarily.
(c) Conclude from (b) that the power series in $x$ for the general solution of (A) is

$$
\begin{aligned}
y= & a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{p((k+2) j)}{(k+2)(j+1)-1}\right] \frac{x^{(k+2) m}}{(k+2)^{m} m!} \\
& +a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{p((k+2) j+1)}{(k+2)(j+1)+1}\right] \frac{x^{(k+2) m+1}}{(k+2)^{m} m!} .
\end{aligned}
$$

In Exercises 39-44 use the method of Exercise 38 to find the power series in $x$ for the general solution.
39. $\left(1+2 x^{5}\right) y^{\prime \prime}+14 x^{4} y^{\prime}+10 x^{3} y=0$
40. $y^{\prime \prime}+x^{2} y=0 \quad$ 41. $y^{\prime \prime}+x^{6} y^{\prime}+7 x^{5} y=0$
42. $\left(1+x^{8}\right) y^{\prime \prime}-16 x^{7} y^{\prime}+72 x^{6} y=0$
43. $\left(1-x^{6}\right) y^{\prime \prime}-12 x^{5} y^{\prime}-30 x^{4} y=0$
44. $y^{\prime \prime}+x^{5} y^{\prime}+6 x^{4} y=0$

### 7.3 SERIES SOLUTIONS NEAR AN ORDINARY POINT II

In this section we continue to find series solutions

$$
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

of initial value problems

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0, \quad y\left(x_{0}\right)=a_{0}, \quad y^{\prime}\left(x_{0}\right)=a_{1}, \tag{7.3.1}
\end{equation*}
$$

where $P_{0}, P_{1}$, and $P_{2}$ are polynomials and $P_{0}\left(x_{0}\right) \neq 0$, so $x_{0}$ is an ordinary point of (7.3.1). However, here we consider cases where the differential equation in (7.3.1) is not of the form

$$
\left(1+\alpha\left(x-x_{0}\right)^{2}\right) y^{\prime \prime}+\beta\left(x-x_{0}\right) y^{\prime}+\gamma y=0
$$

so Theorem 7.2.2 does not apply, and the computation of the coefficients $\left\{a_{n}\right\}$ is more complicated. For the equations considered here it's difficult or impossible to obtain an explicit formula for $a_{n}$ in terms of $n$. Nevertheless, we can calculate as many coefficients as we wish. The next three examples illustrate this.

Example 7.3.1 Find the coefficients $a_{0}, \ldots, a_{7}$ in the series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of the initial value problem

$$
\begin{equation*}
\left(1+x+2 x^{2}\right) y^{\prime \prime}+(1+7 x) y^{\prime}+2 y=0, \quad y(0)=-1, \quad y^{\prime}(0)=-2 . \tag{7.3.2}
\end{equation*}
$$

## Solution Here

$$
L y=\left(1+x+2 x^{2}\right) y^{\prime \prime}+(1+7 x) y^{\prime}+2 y .
$$

The zeros $(-1 \pm i \sqrt{7}) / 4$ of $P_{0}(x)=1+x+2 x^{2}$ have absolute value $1 / \sqrt{2}$, so Theorem 7.2.2 implies that the series solution converges to the solution of (7.3.2) on $(-1 / \sqrt{2}, 1 / \sqrt{2})$. Since

$$
\begin{aligned}
y= & \sum_{n=0}^{\infty} a_{n} x^{n}, \quad y^{\prime}=\sum_{n=1}^{\infty} n a_{n} x^{n-1} \quad \text { and } \quad y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}, \\
L y= & \sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}+\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-1}+2 \sum_{n=2}^{\infty} n(n-1) a_{n} x^{n} \\
& +\sum_{n=1}^{\infty} n a_{n} x^{n-1}+7 \sum_{n=1}^{\infty} n a_{n} x^{n}+2 \sum_{n=0}^{\infty} a_{n} x^{n} .
\end{aligned}
$$

Shifting indices so the general term in each series is a constant multiple of $x^{n}$ yields

$$
\begin{aligned}
L y= & \sum_{n=0}^{\infty}(n+2)(n+1) a_{n+2} x^{n}+\sum_{n=0}^{\infty}(n+1) n a_{n+1} x^{n}+2 \sum_{n=0}^{\infty} n(n-1) a_{n} x^{n} \\
& +\sum_{n=0}^{\infty}(n+1) a_{n+1} x^{n}+7 \sum_{n=0}^{\infty} n a_{n} x^{n}+2 \sum_{n=0}^{\infty} a_{n} x^{n}=\sum_{n=0}^{\infty} b_{n} x^{n},
\end{aligned}
$$

where

$$
b_{n}=(n+2)(n+1) a_{n+2}+(n+1)^{2} a_{n+1}+(n+2)(2 n+1) a_{n} .
$$

Therefore $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ is a solution of $L y=0$ if and only if

$$
\begin{equation*}
a_{n+2}=-\frac{n+1}{n+2} a_{n+1}-\frac{2 n+1}{n+1} a_{n}, n \geq 0 . \tag{7.3.3}
\end{equation*}
$$

From the initial conditions in (7.3.2), $a_{0}=y(0)=-1$ and $a_{1}=y^{\prime}(0)=-2$. Setting $n=0$ in (7.3.3) yields

$$
a_{2}=-\frac{1}{2} a_{1}-a_{0}=-\frac{1}{2}(-2)-(-1)=2 .
$$

Setting $n=1$ in (7.3.3) yields

$$
a_{3}=-\frac{2}{3} a_{2}-\frac{3}{2} a_{1}=-\frac{2}{3}(2)-\frac{3}{2}(-2)=\frac{5}{3} .
$$

We leave it to you to compute $a_{4}, a_{5}, a_{6}, a_{7}$ from (7.3.3) and show that

$$
y=-1-2 x+2 x^{2}+\frac{5}{3} x^{3}-\frac{55}{12} x^{4}+\frac{3}{4} x^{5}+\frac{61}{8} x^{6}-\frac{443}{56} x^{7}+\cdots .
$$

We also leave it to you (Exercise 13) to verify numerically that the Taylor polynomials $T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}$ converge to the solution of $(7.3 .2)$ on $(-1 / \sqrt{2}, 1 / \sqrt{2})$.

Example 7.3.2 Find the coefficients $a_{0}, \ldots, a_{5}$ in the series solution

$$
y=\sum_{n=0}^{\infty} a_{n}(x+1)^{n}
$$

of the initial value problem

$$
\begin{equation*}
(3+x) y^{\prime \prime}+(1+2 x) y^{\prime}-(2-x) y=0, \quad y(-1)=2, \quad y^{\prime}(-1)=-3 \tag{7.3.4}
\end{equation*}
$$

Solution Since the desired series is in powers of $x+1$ we rewrite the differential equation in (7.3.4) as $L y=0$, with

$$
L y=(2+(x+1)) y^{\prime \prime}-(1-2(x+1)) y^{\prime}-(3-(x+1)) y
$$

Since

$$
\begin{gathered}
y=\sum_{n=0}^{\infty} a_{n}(x+1)^{n}, \quad y^{\prime}=\sum_{n=1}^{\infty} n a_{n}(x+1)^{n-1} \quad \text { and } \quad y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n}(x+1)^{n-2}, \\
L y=2 \sum_{n=2}^{\infty} n(n-1) a_{n}(x+1)^{n-2}+\sum_{n=2}^{\infty} n(n-1) a_{n}(x+1)^{n-1} \\
-\sum_{n=1}^{\infty} n a_{n}(x+1)^{n-1}+2 \sum_{n=1}^{\infty} n a_{n}(x+1)^{n} \\
-3 \sum_{n=0}^{\infty} a_{n}(x+1)^{n}+\sum_{n=0}^{\infty} a_{n}(x+1)^{n+1}
\end{gathered}
$$

Shifting indices so that the general term in each series is a constant multiple of $(x+1)^{n}$ yields

$$
\begin{aligned}
L y= & 2 \sum_{n=0}^{\infty}(n+2)(n+1) a_{n+2}(x+1)^{n}+\sum_{n=0}^{\infty}(n+1) n a_{n+1}(x+1)^{n} \\
& -\sum_{n=0}^{\infty}(n+1) a_{n+1}(x+1)^{n}+\sum_{n=0}^{\infty}(2 n-3) a_{n}(x+1)^{n}+\sum_{n=1}^{\infty} a_{n-1}(x+1)^{n} \\
= & \sum_{n=0}^{\infty} b_{n}(x+1)^{n},
\end{aligned}
$$

where

$$
b_{0}=4 a_{2}-a_{1}-3 a_{0}
$$

and

$$
b_{n}=2(n+2)(n+1) a_{n+2}+\left(n^{2}-1\right) a_{n+1}+(2 n-3) a_{n}+a_{n-1}, \quad n \geq 1
$$

Therefore $y=\sum_{n=0}^{\infty} a_{n}(x+1)^{n}$ is a solution of $L y=0$ if and only if

$$
\begin{equation*}
a_{2}=\frac{1}{4}\left(a_{1}+3 a_{0}\right) \tag{7.3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{n+2}=-\frac{1}{2(n+2)(n+1)}\left[\left(n^{2}-1\right) a_{n+1}+(2 n-3) a_{n}+a_{n-1}\right], \quad n \geq 1 . \tag{7.3.6}
\end{equation*}
$$

From the initial conditions in (7.3.4), $a_{0}=y(-1)=2$ and $a_{1}=y^{\prime}(-1)=-3$. We leave it to you to compute $a_{2}, \ldots, a_{5}$ with (7.3.5) and (7.3.6) and show that the solution of (7.3.4) is

$$
y=-2-3(x+1)+\frac{3}{4}(x+1)^{2}-\frac{5}{12}(x+1)^{3}+\frac{7}{48}(x+1)^{4}-\frac{1}{60}(x+1)^{5}+\cdots .
$$

We also leave it to you (Exercise 14) to verify numerically that the Taylor polynomials $T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}$ converge to the solution of (7.3.4) on the interval of convergence of the power series solution.

Example 7.3.3 Find the coefficients $a_{0}, \ldots, a_{5}$ in the series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+3 x y^{\prime}+\left(4+2 x^{2}\right) y=0, \quad y(0)=2, \quad y^{\prime}(0)=-3 \tag{7.3.7}
\end{equation*}
$$

## Solution Here

$$
L y=y^{\prime \prime}+3 x y^{\prime}+\left(4+2 x^{2}\right) y
$$

Since

$$
\begin{aligned}
y & =\sum_{n=0}^{\infty} a_{n} x^{n}, \quad y^{\prime}=\sum_{n=1}^{\infty} n a_{n} x^{n-1}, \quad \text { and } \quad y^{\prime \prime}=\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}, \\
L y & =\sum_{n=2}^{\infty} n(n-1) a_{n} x^{n-2}+3 \sum_{n=1}^{\infty} n a_{n} x^{n}+4 \sum_{n=0}^{\infty} a_{n} x^{n}+2 \sum_{n=0}^{\infty} a_{n} x^{n+2} .
\end{aligned}
$$

Shifting indices so that the general term in each series is a constant multiple of $x^{n}$ yields

$$
L y=\sum_{n=0}^{\infty}(n+2)(n+1) a_{n+2} x^{n}+\sum_{n=0}^{\infty}(3 n+4) a_{n} x^{n}+2 \sum_{n=2}^{\infty} a_{n-2} x^{n}=\sum_{n=0}^{\infty} b_{n} x^{n}
$$

where

$$
b_{0}=2 a_{2}+4 a_{0}, \quad b_{1}=6 a_{3}+7 a_{1},
$$

and

$$
b_{n}=(n+2)(n+1) a_{n+2}+(3 n+4) a_{n}+2 a_{n-2}, \quad n \geq 2
$$

Therefore $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ is a solution of $L y=0$ if and only if

$$
\begin{equation*}
a_{2}=-2 a_{0}, \quad a_{3}=-\frac{7}{6} a_{1}, \tag{7.3.8}
\end{equation*}
$$

and

$$
\begin{equation*}
a_{n+2}=-\frac{1}{(n+2)(n+1)}\left[(3 n+4) a_{n}+2 a_{n-2}\right], \quad n \geq 2 . \tag{7.3.9}
\end{equation*}
$$

From the initial conditions in (7.3.7), $a_{0}=y(0)=2$ and $a_{1}=y^{\prime}(0)=-3$. We leave it to you to compute $a_{2}, \ldots, a_{5}$ with (7.3.8) and (7.3.9) and show that the solution of (7.3.7) is

$$
y=2-3 x-4 x^{2}+\frac{7}{2} x^{3}+3 x^{4}-\frac{79}{40} x^{5}+\cdots
$$

We also leave it to you (Exercise 15) to verify numerically that the Taylor polynomials $T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}$ converge to the solution of (7.3.9) on the interval of convergence of the power series solution.

### 7.3 Exercises

In Exercises 1-12 find the coefficients $a_{0}, \ldots, a_{N}$ for $N$ at least 7 in the series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of the initial value problem.

1. C $(1+3 x) y^{\prime \prime}+x y^{\prime}+2 y=0, \quad y(0)=2, \quad y^{\prime}(0)=-3$
2. $\mathrm{C}\left(1+x+2 x^{2}\right) y^{\prime \prime}+(2+8 x) y^{\prime}+4 y=0, \quad y(0)=-1, \quad y^{\prime}(0)=2$
3. $\mathrm{C}\left(1-2 x^{2}\right) y^{\prime \prime}+(2-6 x) y^{\prime}-2 y=0, \quad y(0)=1, \quad y^{\prime}(0)=0$
4. $\mathrm{C}\left(1+x+3 x^{2}\right) y^{\prime \prime}+(2+15 x) y^{\prime}+12 y=0, \quad y(0)=0, \quad y^{\prime}(0)=1$
5. $\mathrm{C}(2+x) y^{\prime \prime}+(1+x) y^{\prime}+3 y=0, \quad y(0)=4, \quad y^{\prime}(0)=3$
6. $\mathrm{C}\left(3+3 x+x^{2}\right) y^{\prime \prime}+(6+4 x) y^{\prime}+2 y=0, \quad y(0)=7, \quad y^{\prime}(0)=3$
7. $\mathrm{C}(4+x) y^{\prime \prime}+(2+x) y^{\prime}+2 y=0, \quad y(0)=2, \quad y^{\prime}(0)=5$
8. C $\left(2-3 x+2 x^{2}\right) y^{\prime \prime}-(4-6 x) y^{\prime}+2 y=0, \quad y(1)=1, \quad y^{\prime}(1)=-1$
9. $\mathrm{C}\left(3 x+2 x^{2}\right) y^{\prime \prime}+10(1+x) y^{\prime}+8 y=0, \quad y(-1)=1, \quad y^{\prime}(-1)=-1$
10. $\mathrm{C}\left(1-x+x^{2}\right) y^{\prime \prime}-(1-4 x) y^{\prime}+2 y=0, \quad y(1)=2, \quad y^{\prime}(1)=-1$
11. $\mathrm{C}(2+x) y^{\prime \prime}+(2+x) y^{\prime}+y=0, \quad y(-1)=-2, \quad y^{\prime}(-1)=3$
12. C $x^{2} y^{\prime \prime}-(6-7 x) y^{\prime}+8 y=0, \quad y(1)=1, \quad y^{\prime}(1)=-2$
13. $L$ Do the following experiment for various choices of real numbers $a_{0}$, $a_{1}$, and $r$, with $0<r<$ $1 / \sqrt{2}$.
(a) Use differential equations software to solve the initial value problem

$$
\begin{equation*}
\left(1+x+2 x^{2}\right) y^{\prime \prime}+(1+7 x) y^{\prime}+2 y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}, \tag{A}
\end{equation*}
$$

numerically on $(-r, r)$. (See Example 7.3.1.)
(b) For $N=2,3,4, \ldots$, compute $a_{2}, \ldots, a_{N}$ in the power series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of (A), and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on $(-r, r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs.
14. L Do the following experiment for various choices of real numbers $a_{0}$, $a_{1}$, and $r$, with $0<r<2$.
(a) Use differential equations software to solve the initial value problem

$$
\begin{equation*}
(3+x) y^{\prime \prime}+(1+2 x) y^{\prime}-(2-x) y=0, \quad y(-1)=a_{0}, \quad y^{\prime}(-1)=a_{1} \tag{A}
\end{equation*}
$$

numerically on $(-1-r,-1+r)$. (See Example 7.3.2. Why this interval?)
(b) For $N=2,3,4, \ldots$, compute $a_{2}, \ldots, a_{N}$ in the power series solution

$$
y=\sum_{n=0}^{\infty} a_{n}(x+1)^{n}
$$

of (A), and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n}(x+1)^{n}
$$

and the solution obtained in (a) on $(-1-r,-1+r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs.
15. L Do the following experiment for several choices of $a_{0}, a_{1}$, and $r$, with $r>0$.
(a) Use differential equations software to solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+3 x y^{\prime}+\left(4+2 x^{2}\right) y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}, \tag{A}
\end{equation*}
$$

numerically on $(-r, r)$. (See Example 7.3.3.)
(b) Find the coefficients $a_{0}, a_{1}, \ldots, a_{N}$ in the power series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of (A), and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on $(-r, r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs.
16. L Do the following experiment for several choices of $a_{0}$ and $a_{1}$.
(a) Use differential equations software to solve the initial value problem

$$
\begin{equation*}
(1-x) y^{\prime \prime}-(2-x) y^{\prime}+y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}, \tag{A}
\end{equation*}
$$

numerically on $(-r, r)$.
(b) Find the coefficients $a_{0}, a_{1}, \ldots, a_{N}$ in the power series solution $y=\sum_{n=0}^{N} a_{n} x^{n}$ of (A), and graph

$$
T_{N}(x)=\sum_{n=0}^{N} a_{n} x^{n}
$$

and the solution obtained in (a) on $(-r, r)$. Continue increasing $N$ until there's no perceptible difference between the two graphs. What happens as you let $r \rightarrow 1$ ?
17. $L$ Follow the directions of Exercise 16 for the initial value problem

$$
(1+x) y^{\prime \prime}+3 y^{\prime}+32 y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1} .
$$

18. L Follow the directions of Exercise 16 for the initial value problem

$$
\left(1+x^{2}\right) y^{\prime \prime}+y^{\prime}+2 y=0, \quad y(0)=a_{0}, \quad y^{\prime}(0)=a_{1}
$$

In Exercises 19-28 find the coefficients $a_{0}, \ldots, a_{N}$ for $N$ at least 7 in the series solution

$$
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

of the initial value problem. Take $x_{0}$ to be the point where the initial conditions are imposed.
19. $\mathrm{C}(2+4 x) y^{\prime \prime}-4 y^{\prime}-(6+4 x) y=0, \quad y(0)=2, \quad y^{\prime}(0)=-7$
20. $\mathrm{C}(1+2 x) y^{\prime \prime}-(1-2 x) y^{\prime}-(3-2 x) y=0, \quad y(1)=1, \quad y^{\prime}(1)=-2$
21. $\mathrm{C}(5+2 x) y^{\prime \prime}-y^{\prime}+(5+x) y=0, \quad y(-2)=2, \quad y^{\prime}(-2)=-1$
22. $\mathrm{C}(4+x) y^{\prime \prime}-(4+2 x) y^{\prime}+(6+x) y=0, \quad y(-3)=2, \quad y^{\prime}(-3)=-2$
23. C $(2+3 x) y^{\prime \prime}-x y^{\prime}+2 x y=0, \quad y(0)=-1, \quad y^{\prime}(0)=2$
24. $\mathrm{C}(3+2 x) y^{\prime \prime}+3 y^{\prime}-x y=0, \quad y(-1)=2, \quad y^{\prime}(-1)=-3$
25. $\mathrm{C}(3+2 x) y^{\prime \prime}-3 y^{\prime}-(2+x) y=0, \quad y(-2)=-2, \quad y^{\prime}(-2)=3$
26. $\mathrm{C}(10-2 x) y^{\prime \prime}+(1+x) y=0, \quad y(2)=2, \quad y^{\prime}(2)=-4$
27. $\mathrm{C}(7+x) y^{\prime \prime}+(8+2 x) y^{\prime}+(5+x) y=0, \quad y(-4)=1, \quad y^{\prime}(-4)=2$
28. $\mathrm{C}(6+4 x) y^{\prime \prime}+(1+2 x) y=0, \quad y(-1)=-1, \quad y^{\prime}(-1)=2$
29. Show that the coefficients in the power series in $x$ for the general solution of

$$
\left(1+\alpha x+\beta x^{2}\right) y^{\prime \prime}+(\gamma+\delta x) y^{\prime}+\epsilon y=0
$$

satisfy the recurrrence relation

$$
a_{n+2}=-\frac{\gamma+\alpha n}{n+2} a_{n+1}-\frac{\beta n(n-1)+\delta n+\epsilon}{(n+2)(n+1)} a_{n}
$$

30. (a) Let $\alpha$ and $\beta$ be constants, with $\beta \neq 0$. Show that $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ is a solution of

$$
\begin{equation*}
\left(1+\alpha x+\beta x^{2}\right) y^{\prime \prime}+(2 \alpha+4 \beta x) y^{\prime}+2 \beta y=0 \tag{A}
\end{equation*}
$$

if and only if

$$
\begin{equation*}
a_{n+2}+\alpha a_{n+1}+\beta a_{n}=0, \quad n \geq 0 \tag{B}
\end{equation*}
$$

An equation of this form is called a second order homogeneous linear difference equation. The polynomial $p(r)=r^{2}+\alpha r+\beta$ is called the characteristic polynomial of (B). If $r_{1}$ and $r_{2}$ are the zeros of $p$, then $1 / r_{1}$ and $1 / r_{2}$ are the zeros of

$$
P_{0}(x)=1+\alpha x+\beta x^{2} .
$$

(b) Suppose $p(r)=\left(r-r_{1}\right)\left(r-r_{2}\right)$ where $r_{1}$ and $r_{2}$ are real and distinct, and let $\rho$ be the smaller of the two numbers $\left\{1 /\left|r_{1}\right|, 1 /\left|r_{2}\right|\right\}$. Show that if $c_{1}$ and $c_{2}$ are constants then the sequence

$$
a_{n}=c_{1} r_{1}^{n}+c_{2} r_{2}^{n}, \quad n \geq 0
$$

satisfies (B). Conclude from this that any function of the form

$$
y=\sum_{n=0}^{\infty}\left(c_{1} r_{1}^{n}+c_{2} r_{2}^{n}\right) x^{n}
$$

is a solution of $(\mathrm{A})$ on $(-\rho, \rho)$.
(c) Use (b) and the formula for the sum of a geometric series to show that the functions

$$
y_{1}=\frac{1}{1-r_{1} x} \quad \text { and } \quad y_{2}=\frac{1}{1-r_{2} x}
$$

form a fundamental set of solutions of (A) on $(-\rho, \rho)$.
(d) Show that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (A) on any interval that does'nt contain either $1 / r_{1}$ or $1 / r_{2}$.
(e) Suppose $p(r)=\left(r-r_{1}\right)^{2}$, and let $\rho=1 /\left|r_{1}\right|$. Show that if $c_{1}$ and $c_{2}$ are constants then the sequence

$$
a_{n}=\left(c_{1}+c_{2} n\right) r_{1}^{n}, \quad n \geq 0
$$

satisfies (B). Conclude from this that any function of the form

$$
y=\sum_{n=0}^{\infty}\left(c_{1}+c_{2} n\right) r_{1}^{n} x^{n}
$$

is a solution of $(\mathrm{A})$ on $(-\rho, \rho)$.
(f) Use (e) and the formula for the sum of a geometric series to show that the functions

$$
y_{1}=\frac{1}{1-r_{1} x} \quad \text { and } \quad y_{2}=\frac{x}{\left(1-r_{1} x\right)^{2}}
$$

form a fundamental set of solutions of (A) on $(-\rho, \rho)$.
(g) Show that $\left\{y_{1}, y_{2}\right\}$ is a fundamental set of solutions of (A) on any interval that does not contain $1 / r_{1}$.
31. Use the results of Exercise 30 to find the general solution of the given equation on any interval on which polynomial multiplying $y^{\prime \prime}$ has no zeros.
(a) $\left(1+3 x+2 x^{2}\right) y^{\prime \prime}+(6+8 x) y^{\prime}+4 y=0$
(b) $\left(1-5 x+6 x^{2}\right) y^{\prime \prime}-(10-24 x) y^{\prime}+12 y=0$
(c) $\left(1-4 x+4 x^{2}\right) y^{\prime \prime}-(8-16 x) y^{\prime}+8 y=0$
(d) $\left(4+4 x+x^{2}\right) y^{\prime \prime}+(8+4 x) y^{\prime}+2 y=0$
(e) $\left(4+8 x+3 x^{2}\right) y^{\prime \prime}+(16+12 x) y^{\prime}+6 y=0$

In Exercises 32-38 find the coefficients $a_{0}, \ldots, a_{N}$ for $N$ at least 7 in the series solution $y=\sum_{n=0}^{\infty} a_{n} x^{n}$ of the initial value problem.
32. C $y^{\prime \prime}+2 x y^{\prime}+\left(3+2 x^{2}\right) y=0, \quad y(0)=1, \quad y^{\prime}(0)=-2$
33. C $y^{\prime \prime}-3 x y^{\prime}+\left(5+2 x^{2}\right) y=0, \quad y(0)=1, \quad y^{\prime}(0)=-2$
34. C $y^{\prime \prime}+5 x y^{\prime}-\left(3-x^{2}\right) y=0, \quad y(0)=6, \quad y^{\prime}(0)=-2$
35. C $y^{\prime \prime}-2 x y^{\prime}-\left(2+3 x^{2}\right) y=0, \quad y(0)=2, \quad y^{\prime}(0)=-5$
36. C $y^{\prime \prime}-3 x y^{\prime}+\left(2+4 x^{2}\right) y=0, \quad y(0)=3, \quad y^{\prime}(0)=6$
37. C $2 y^{\prime \prime}+5 x y^{\prime}+\left(4+2 x^{2}\right) y=0, \quad y(0)=3, \quad y^{\prime}(0)=-2$
38. C $3 y^{\prime \prime}+2 x y^{\prime}+\left(4-x^{2}\right) y=0, \quad y(0)=-2, \quad y^{\prime}(0)=3$
39. Find power series in $x$ for the solutions $y_{1}$ and $y_{2}$ of

$$
y^{\prime \prime}+4 x y^{\prime}+\left(2+4 x^{2}\right) y=0
$$

such that $y_{1}(0)=1, y_{1}^{\prime}(0)=0, y_{2}(0)=0, y_{2}^{\prime}(0)=1$, and identify $y_{1}$ and $y_{2}$ in terms of familiar elementary functions.

In Exercises 40-49 find the coefficients $a_{0}, \ldots, a_{N}$ for $N$ at least 7 in the series solution

$$
y=\sum_{n=0}^{\infty} a_{n}\left(x-x_{0}\right)^{n}
$$

of the initial value problem. Take $x_{0}$ to be the point where the initial conditions are imposed.
40. $\mathrm{C}(1+x) y^{\prime \prime}+x^{2} y^{\prime}+(1+2 x) y=0, \quad y(0)-2, \quad y^{\prime}(0)=3$
41. C $y^{\prime \prime}+\left(1+2 x+x^{2}\right) y^{\prime}+2 y=0, \quad y(0)=2, \quad y^{\prime}(0)=3$
42. C $\left(1+x^{2}\right) y^{\prime \prime}+\left(2+x^{2}\right) y^{\prime}+x y=0, \quad y(0)=-3, \quad y^{\prime}(0)=5$
43. $\quad$ C $(1+x) y^{\prime \prime}+\left(1-3 x+2 x^{2}\right) y^{\prime}-(x-4) y=0, \quad y(1)=-2, \quad y^{\prime}(1)=3$
44. C $y^{\prime \prime}+\left(13+12 x+3 x^{2}\right) y^{\prime}+(5+2 x), \quad y(-2)=2, \quad y^{\prime}(-2)=-3$
45. $\mathrm{C}\left(1+2 x+3 x^{2}\right) y^{\prime \prime}+\left(2-x^{2}\right) y^{\prime}+(1+x) y=0, \quad y(0)=1, \quad y^{\prime}(0)=-2$
46. $\mathrm{C}\left(3+4 x+x^{2}\right) y^{\prime \prime}-\left(5+4 x-x^{2}\right) y^{\prime}-(2+x) y=0, \quad y(-2)=2, \quad y^{\prime}(-2)=-1$
47. $\mathrm{C}\left(1+2 x+x^{2}\right) y^{\prime \prime}+(1-x) y=0, \quad y(0)=2, \quad y^{\prime}(0)=-1$
48. $\mathrm{C}\left(x-2 x^{2}\right) y^{\prime \prime}+\left(1+3 x-x^{2}\right) y^{\prime}+(2+x) y=0, \quad y(1)=1, \quad y^{\prime}(1)=0$
49. $\mathrm{C}\left(16-11 x+2 x^{2}\right) y^{\prime \prime}+\left(10-6 x+x^{2}\right) y^{\prime}-(2-x) y, \quad y(3)=1, \quad y^{\prime}(3)=-2$

## CHAPTER 8 Laplace Transforms

IN THIS CHAPTER we study the method of Laplace transforms, which illustrates one of the basic problem solving techniques in mathematics: transform a difficult problem into an easier one, solve the latter, and then use its solution to obtain a solution of the original problem. The method discussed here transforms an initial value problem for a constant coefficient equation into an algebraic equation whose solution can then be used to solve the initial value problem. In some cases this method is merely an alternative procedure for solving problems that can be solved equally well by methods that we considered previously; however, in other cases the method of Laplace transforms is more efficient than the methods previously discussed. This is especially true in physical problems dealing with discontinuous forcing functions.

SECTION 8.1 defines the Laplace transform and developes its properties.
SECTION 8.2 deals with the problem of finding a function that has a given Laplace transform.
SECTION 8.3 applies the Laplace transform to solve initial value problems for constant coefficient second order differential equations on $(0, \infty)$.

SECTION 8.4 introduces the unit step function.
SECTION 8.5 uses the unit step function to solve constant coefficient equations with piecewise continuous forcing functions.

SECTION 8.6 deals with the convolution theorem, an important theoretical property of the Laplace transform.
SECTION 8.7 introduces the idea of impulsive force, and treats constant coefficient equations with impulsive forcing functions.

SECTION 8.8 is a brief table of Laplace transforms.

### 8.1 INTRODUCTION TO THE LAPLACE TRANSFORM

Definition of the Laplace Transform
To define the Laplace transform, we first recall the definition of an improper integral. If $g$ is integrable over the interval $[a, T]$ for every $T>a$, then the improper integral of $g$ over $[a, \infty)$ is defined as

$$
\begin{equation*}
\int_{a}^{\infty} g(t) d t=\lim _{T \rightarrow \infty} \int_{a}^{T} g(t) d t \tag{8.1.1}
\end{equation*}
$$

We say that the improper integral converges if the limit in (8.1.1) exists; otherwise, we say that the improper integral diverges or does not exist. Here's the definition of the Laplace transform of a function $f$.

Definition 8.1.1 Let $f$ be defined for $t \geq 0$ and let $s$ be a real number. Then the Laplace transform of $f$ is the function $F$ defined by

$$
\begin{equation*}
F(s)=\int_{0}^{\infty} e^{-s t} f(t) d t \tag{8.1.2}
\end{equation*}
$$

for those values of $s$ for which the improper integral converges.
It is important to keep in mind that the variable of integration in (8.1.2) is $t$, while $s$ is a parameter independent of $t$. We use $t$ as the independent variable for $f$ because in applications the Laplace transform is usually applied to functions of time.

The Laplace transform can be viewed as an operator $\mathcal{L}$ that transforms the function $f=f(t)$ into the function $F=F(s)$. Thus, (8.1.2) can be expressed as

$$
F=\mathcal{L}(f)
$$

The functions $f$ and $F$ form a transform pair, which we'll sometimes denote by

$$
f(t) \leftrightarrow F(s) .
$$

It can be shown that if $F(s)$ is defined for $s=s_{0}$ then it's defined for all $s>s_{0}$ (Exercise 14(b)).
Computation of Some Simple Laplace Transforms

Example 8.1.1 Find the Laplace transform of $f(t)=1$.

Solution From (8.1.2) with $f(t)=1$,

$$
F(s)=\int_{0}^{\infty} e^{-s t} d t=\lim _{T \rightarrow \infty} \int_{0}^{T} e^{-s t} d t
$$

If $s \neq 0$ then

$$
\begin{equation*}
\int_{0}^{T} e^{-s t} d t=-\left.\frac{1}{s} e^{-s t}\right|_{0} ^{T}=\frac{1-e^{-s T}}{s} . \tag{8.1.3}
\end{equation*}
$$

Therefore

$$
\lim _{T \rightarrow \infty} \int_{0}^{T} e^{-s t} d t=\left\{\begin{array}{cc}
\frac{1}{s}, & s>0  \tag{8.1.4}\\
\infty, & s<0
\end{array}\right.
$$

If $s=0$ the integrand reduces to the constant 1 , and

$$
\lim _{T \rightarrow \infty} \int_{0}^{T} 1 d t=\lim _{T \rightarrow \infty} \int_{0}^{T} 1 d t=\lim _{T \rightarrow \infty} T=\infty
$$

Therefore $F(0)$ is undefined, and

$$
F(s)=\int_{0}^{\infty} e^{-s t} d t=\frac{1}{s}, \quad s>0
$$

This result can be written in operator notation as

$$
\mathcal{L}(1)=\frac{1}{s}, \quad s>0,
$$

or as the transform pair

$$
1 \leftrightarrow \frac{1}{s}, \quad s>0 .
$$

REMARK: It is convenient to combine the steps of integrating from 0 to $T$ and letting $T \rightarrow \infty$. Therefore, instead of writing (8.1.3) and (8.1.4) as separate steps we write

$$
\int_{0}^{\infty} e^{-s t} d t=-\left.\frac{1}{s} e^{-s t}\right|_{0} ^{\infty}=\left\{\begin{array}{cc}
\frac{1}{s}, & s>0 \\
\infty, & s<0
\end{array}\right.
$$

We'll follow this practice throughout this chapter.
Example 8.1.2 Find the Laplace transform of $f(t)=t$.

Solution From (8.1.2) with $f(t)=t$,

$$
\begin{equation*}
F(s)=\int_{0}^{\infty} e^{-s t} t d t \tag{8.1.5}
\end{equation*}
$$

If $s \neq 0$, integrating by parts yields

$$
\begin{aligned}
\int_{0}^{\infty} e^{-s t} t d t & =-\left.\frac{t e^{-s t}}{s}\right|_{0} ^{\infty}+\frac{1}{s} \int_{0}^{\infty} e^{-s t} d t=-\left.\left[\frac{t}{s}+\frac{1}{s^{2}}\right] e^{-s t}\right|_{0} ^{\infty} \\
& =\left\{\begin{aligned}
\frac{1}{s^{2}}, & s>0 \\
\infty, & s<0
\end{aligned}\right.
\end{aligned}
$$

If $s=0$, the integral in (8.1.5) becomes

$$
\int_{0}^{\infty} t d t=\left.\frac{t^{2}}{2}\right|_{0} ^{\infty}=\infty
$$

Therefore $F(0)$ is undefined and

$$
F(s)=\frac{1}{s^{2}}, \quad s>0
$$

This result can also be written as

$$
\mathcal{L}(t)=\frac{1}{s^{2}}, \quad s>0
$$

or as the transform pair

$$
t \leftrightarrow \frac{1}{s^{2}}, \quad s>0 .
$$

Example 8.1.3 Find the Laplace transform of $f(t)=e^{a t}$, where $a$ is a constant.

Solution From (8.1.2) with $f(t)=e^{a t}$,

$$
F(s)=\int_{0}^{\infty} e^{-s t} e^{a t} d t
$$

Combining the exponentials yields

$$
F(s)=\int_{0}^{\infty} e^{-(s-a) t} d t
$$

However, we know from Example 8.1.1 that

$$
\int_{0}^{\infty} e^{-s t} d t=\frac{1}{s}, \quad s>0
$$

Replacing $s$ by $s-a$ here shows that

$$
F(s)=\frac{1}{s-a}, \quad s>a
$$

This can also be written as

$$
\mathcal{L}\left(e^{a t}\right)=\frac{1}{s-a}, \quad s>a, \quad \text { or } \quad e^{a t} \leftrightarrow \frac{1}{s-a}, \quad s>a
$$

Example 8.1.4 Find the Laplace transforms of $f(t)=\sin \omega t$ and $g(t)=\cos \omega t$, where $\omega$ is a constant.

## Solution Define

$$
\begin{equation*}
F(s)=\int_{0}^{\infty} e^{-s t} \sin \omega t d t \tag{8.1.6}
\end{equation*}
$$

and

$$
\begin{equation*}
G(s)=\int_{0}^{\infty} e^{-s t} \cos \omega t d t \tag{8.1.7}
\end{equation*}
$$

If $s>0$, integrating (8.1.6) by parts yields

$$
F(s)=-\left.\frac{e^{-s t}}{s} \sin \omega t\right|_{0} ^{\infty}+\frac{\omega}{s} \int_{0}^{\infty} e^{-s t} \cos \omega t d t
$$

so

$$
\begin{equation*}
F(s)=\frac{\omega}{s} G(s) \tag{8.1.8}
\end{equation*}
$$

If $s>0$, integrating (8.1.7) by parts yields

$$
G(s)=-\left.\frac{e^{-s t} \cos \omega t}{s}\right|_{0} ^{\infty}-\frac{\omega}{s} \int_{0}^{\infty} e^{-s t} \sin \omega t d t
$$

so

$$
G(s)=\frac{1}{s}-\frac{\omega}{s} F(s) .
$$

Now substitute from (8.1.8) into this to obtain

$$
G(s)=\frac{1}{s}-\frac{\omega^{2}}{s^{2}} G(s)
$$

Solving this for $G(s)$ yields

$$
G(s)=\frac{s}{s^{2}+\omega^{2}}, \quad s>0
$$

This and (8.1.8) imply that

$$
F(s)=\frac{\omega}{s^{2}+\omega^{2}}, \quad s>0
$$

Tables of Laplace transforms
Extensive tables of Laplace transforms have been compiled and are commonly used in applications. The brief table of Laplace transforms in the Appendix will be adequate for our purposes.

Example 8.1.5 Use the table of Laplace transforms to find $\mathcal{L}\left(t^{3} e^{4 t}\right)$.

Solution The table includes the transform pair

$$
t^{n} e^{a t} \leftrightarrow \frac{n!}{(s-a)^{n+1}}
$$

Setting $n=3$ and $a=4$ here yields

$$
\mathcal{L}\left(t^{3} e^{4 t}\right)=\frac{3!}{(s-4)^{4}}=\frac{6}{(s-4)^{4}}
$$

We'll sometimes write Laplace transforms of specific functions without explicitly stating how they are obtained. In such cases you should refer to the table of Laplace transforms.
Linearity of the Laplace Transform
The next theorem presents an important property of the Laplace transform.

Theorem 8.1.2 [Linearity Property] Suppose $\mathcal{L}\left(f_{i}\right)$ is defined for $s>s_{i}, 1 \leq i \leq n$. Let $s_{0}$ be the largest of the numbers $s_{1}, s_{2}, \ldots, s_{n}$, and let $c_{1}, c_{2}, \ldots, c_{n}$ be constants. Then

$$
\mathcal{L}\left(c_{1} f_{1}+c_{2} f_{2}+\cdots+c_{n} f_{n}\right)=c_{1} \mathcal{L}\left(f_{1}\right)+c_{2} \mathcal{L}\left(f_{2}\right)+\cdots+c_{n} \mathcal{L}\left(f_{n}\right) \text { for } s>s_{0} .
$$

Proof We give the proof for the case where $n=2$. If $s>s_{0}$ then

$$
\begin{aligned}
\mathcal{L}\left(c_{1} f_{1}+c_{2} f_{2}\right) & =\int_{0}^{\infty} e^{-s t}\left(c_{1} f_{1}(t)+c_{2} f_{2}(t)\right) d t \\
& =c_{1} \int_{0}^{\infty} e^{-s t} f_{1}(t) d t+c_{2} \int_{0}^{\infty} e^{-s t} f_{2}(t) d t \\
& =c_{1} \mathcal{L}\left(f_{1}\right)+c_{2} \mathcal{L}\left(f_{2}\right)
\end{aligned}
$$

Example 8.1.6 Use Theorem 8.1.2 and the known Laplace transform

$$
\mathcal{L}\left(e^{a t}\right)=\frac{1}{s-a}
$$

to find $\mathcal{L}(\cosh b t)(b \neq 0)$.

Solution By definition,

$$
\cosh b t=\frac{e^{b t}+e^{-b t}}{2}
$$

Therefore

$$
\begin{align*}
\mathcal{L}(\cosh b t) & =\mathcal{L}\left(\frac{1}{2} e^{b t}+\frac{1}{2} e^{-b t}\right) \\
& =\frac{1}{2} \mathcal{L}\left(e^{b t}\right)+\frac{1}{2} \mathcal{L}\left(e^{-b t}\right) \quad \text { (linearity property) }  \tag{8.1.9}\\
& =\frac{1}{2} \frac{1}{s-b}+\frac{1}{2} \frac{1}{s+b}
\end{align*}
$$

where the first transform on the right is defined for $s>b$ and the second for $s>-b$; hence, both are defined for $s>|b|$. Simplifying the last expression in (8.1.9) yields

$$
\mathcal{L}(\cosh b t)=\frac{s}{s^{2}-b^{2}}, \quad s>|b| .
$$

The First Shifting Theorem
The next theorem enables us to start with known transform pairs and derive others. (For other results of this kind, see Exercises 6 and 13.)

Theorem 8.1.3 [First Shifting Theorem] If

$$
\begin{equation*}
F(s)=\int_{0}^{\infty} e^{-s t} f(t) d t \tag{8.1.10}
\end{equation*}
$$

is the Laplace transform of $f(t)$ for $s>s_{0}$, then $F(s-a)$ is the Laplace transform of $e^{a t} f(t)$ for $s>s_{0}+a$.

PROOF. Replacing $s$ by $s-a$ in (8.1.10) yields

$$
\begin{equation*}
F(s-a)=\int_{0}^{\infty} e^{-(s-a) t} f(t) d t \tag{8.1.11}
\end{equation*}
$$

if $s-a>s_{0}$; that is, if $s>s_{0}+a$. However, (8.1.11) can be rewritten as

$$
F(s-a)=\int_{0}^{\infty} e^{-s t}\left(e^{a t} f(t)\right) d t
$$

which implies the conclusion.

Example 8.1.7 Use Theorem 8.1.3 and the known Laplace transforms of $1, t, \cos \omega t$, and $\sin \omega t$ to find

$$
\mathcal{L}\left(e^{a t}\right), \quad \mathcal{L}\left(t e^{a t}\right), \quad \mathcal{L}\left(e^{\lambda t} \sin \omega t\right), \text { and } \mathcal{L}\left(e^{\lambda t} \cos \omega t\right)
$$

Solution In the following table the known transform pairs are listed on the left and the required transform pairs listed on the right are obtained by applying Theorem 8.1.3.

| $f(t) \leftrightarrow F(s)$ | $e^{a t} f(t) \leftrightarrow F(s-a)$ |
| :---: | :---: |
| $1 \leftrightarrow \frac{1}{s}, \quad s>0$ | $e^{a t} \leftrightarrow \frac{1}{(s-a)}, \quad s>a$ |
| $t \leftrightarrow \frac{1}{s^{2}}, \quad s>0$ | $t e^{a t} \leftrightarrow \frac{1}{(s-a)^{2}}, \quad s>a$ |
| $\sin \omega t \leftrightarrow \frac{\omega}{s^{2}+\omega^{2}}, \quad s>0$ | $e^{\lambda t} \sin \omega t \leftrightarrow \frac{\omega}{(s-\lambda)^{2}+\omega^{2}}, s>\lambda$ |
| $\cos \omega t \leftrightarrow \frac{s}{s^{2}+\omega^{2}}, \quad s>0$ | $e^{\lambda t} \cos \omega t \leftrightarrow \frac{s-\lambda}{(s-\lambda)^{2}+\omega^{2}}, s>\lambda$ |

Existence of Laplace Transforms
Not every function has a Laplace transform. For example, it can be shown (Exercise 3) that

$$
\int_{0}^{\infty} e^{-s t} e^{t^{2}} d t=\infty
$$

for every real number $s$. Hence, the function $f(t)=e^{t^{2}}$ does not have a Laplace transform.
Our next objective is to establish conditions that ensure the existence of the Laplace transform of a function. We first review some relevant definitions from calculus.
Recall that a limit

$$
\lim _{t \rightarrow t_{0}} f(t)
$$

exists if and only if the one-sided limits

$$
\lim _{t \rightarrow t_{0}-} f(t) \text { and } \lim _{t \rightarrow t_{0}+} f(t)
$$

both exist and are equal; in this case,

$$
\lim _{t \rightarrow t_{0}} f(t)=\lim _{t \rightarrow t_{0}-} f(t)=\lim _{t \rightarrow t_{0}+} f(t) .
$$

Recall also that $f$ is continuous at a point $t_{0}$ in an open interval $(a, b)$ if and only if

$$
\lim _{t \rightarrow t_{0}} f(t)=f\left(t_{0}\right)
$$

which is equivalent to

$$
\begin{equation*}
\lim _{t \rightarrow t_{0}+} f(t)=\lim _{t \rightarrow t_{0}-} f(t)=f\left(t_{0}\right) \tag{8.1.12}
\end{equation*}
$$

For simplicity, we define

$$
f\left(t_{0}+\right)=\lim _{t \rightarrow t_{0}+} f(t) \quad \text { and } \quad f\left(t_{0}-\right)=\lim _{t \rightarrow t_{0}-} f(t)
$$

so (8.1.12) can be expressed as

$$
f\left(t_{0}+\right)=f\left(t_{0}-\right)=f\left(t_{0}\right)
$$

If $f\left(t_{0}+\right)$ and $f\left(t_{0}-\right)$ have finite but distinct values, we say that $f$ has a jump discontinuity at $t_{0}$, and

$$
f\left(t_{0}+\right)-f\left(t_{0}-\right)
$$

is called the jump in $f$ at $t_{0}$ (Figure 8.1.1).


Figure 8.1.1 A jump discontinuity


Figure 8.1.3 A piecewise continuous function on
Figure 8.1.2
$[a, b]$

If $f\left(t_{0}+\right)$ and $f\left(t_{0}-\right)$ are finite and equal, but either $f$ isn't defined at $t_{0}$ or it's defined but

$$
f\left(t_{0}\right) \neq f\left(t_{0}+\right)=f\left(t_{0}-\right)
$$

we say that $f$ has a removable discontinuity at $t_{0}$ (Figure 8.1.2). This terminolgy is appropriate since a function $f$ with a removable discontinuity at $t_{0}$ can be made continuous at $t_{0}$ by defining (or redefining)

$$
f\left(t_{0}\right)=f\left(t_{0}+\right)=f\left(t_{0}-\right)
$$

REMARK: We know from calculus that a definite integral isn't affected by changing the values of its integrand at isolated points. Therefore, redefining a function $f$ to make it continuous at removable discontinuities does not change $\mathcal{L}(f)$.

## Definition 8.1.4

(i) A function $f$ is said to be piecewise continuous on a finite closed interval $[0, T]$ if $f(0+)$ and $f(T-)$ are finite and $f$ is continuous on the open interval $(0, T)$ except possibly at finitely many points, where $f$ may have jump discontinuities or removable discontinuities.
(ii) A function $f$ is said to be piecewise continuous on the infinite interval $[0, \infty)$ if it's piecewise continuous on $[0, T]$ for every $T>0$.

Figure 8.1.3 shows the graph of a typical piecewise continuous function.
It is shown in calculus that if a function is piecewise continuous on a finite closed interval then it's integrable on that interval. But if $f$ is piecewise continuous on $[0, \infty)$, then so is $e^{-s t} f(t)$, and therefore

$$
\int_{0}^{T} e^{-s t} f(t) d t
$$

exists for every $T>0$. However, piecewise continuity alone does not guarantee that the improper integral

$$
\begin{equation*}
\int_{0}^{\infty} e^{-s t} f(t) d t=\lim _{T \rightarrow \infty} \int_{0}^{T} e^{-s t} f(t) d t \tag{8.1.13}
\end{equation*}
$$

converges for $s$ in some interval $\left(s_{0}, \infty\right)$. For example, we noted earlier that (8.1.13) diverges for all $s$ if $f(t)=e^{t^{2}}$. Stated informally, this occurs because $e^{t^{2}}$ increases too rapidly as $t \rightarrow \infty$. The next definition provides a constraint on the growth of a function that guarantees convergence of its Laplace transform for $s$ in some interval $\left(s_{0}, \infty\right)$.

Definition 8.1.5 A function $f$ is said to be of exponential order $s_{0}$ if there are constants $M$ and $t_{0}$ such that

$$
\begin{equation*}
|f(t)| \leq M e^{s_{0} t}, \quad t \geq t_{0} \tag{8.1.14}
\end{equation*}
$$

In situations where the specific value of $s_{0}$ is irrelevant we say simply that $f$ is of exponential order.
The next theorem gives useful sufficient conditions for a function $f$ to have a Laplace transform. The proof is sketched in Exercise 10.

Theorem 8.1.6 If $f$ is piecewise continuous on $[0, \infty)$ and of exponential order $s_{0}$, then $\mathcal{L}(f)$ is defined for $s>s_{0}$.

REMARK: We emphasize that the conditions of Theorem 8.1.6 are sufficient, but not necessary, for $f$ to have a Laplace transform. For example, Exercise 14(c) shows that $f$ may have a Laplace transform even though $f$ isn't of exponential order.

Example 8.1.8 If $f$ is bounded on some interval $\left[t_{0}, \infty\right)$, say

$$
|f(t)| \leq M, \quad t \geq t_{0}
$$

then (8.1.14) holds with $s_{0}=0$, so $f$ is of exponential order zero. Thus, for example, $\sin \omega t$ and $\cos \omega t$ are of exponential order zero, and Theorem 8.1.6 implies that $\mathcal{L}(\sin \omega t)$ and $\mathcal{L}(\cos \omega t)$ exist for $s>0$. This is consistent with the conclusion of Example 8.1.4.

Example 8.1.9 It can be shown that if $\lim _{t \rightarrow \infty} e^{-s_{0} t} f(t)$ exists and is finite then $f$ is of exponential order $s_{0}$ (Exercise 9). If $\alpha$ is any real number and $s_{0}>0$ then $f(t)=t^{\alpha}$ is of exponential order $s_{0}$, since

$$
\lim _{t \rightarrow \infty} e^{-s_{0} t} t^{\alpha}=0
$$

by L'Hôpital's rule. If $\alpha \geq 0, f$ is also continuous on $[0, \infty)$. Therefore Exercise 9 and Theorem 8.1.6 imply that $\mathcal{L}\left(t^{\alpha}\right)$ exists for $s \geq s_{0}$. However, since $s_{0}$ is an arbitrary positive number, this really implies that $\mathcal{L}\left(t^{\alpha}\right)$ exists for all $s>0$. This is consistent with the results of Example 8.1.2 and Exercises 6 and 8.

Example 8.1.10 Find the Laplace transform of the piecewise continuous function

$$
f(t)=\left\{\begin{array}{cl}
1, & 0 \leq t<1 \\
-3 e^{-t}, & t \geq 1
\end{array}\right.
$$

Solution Since $f$ is defined by different formulas on $[0,1)$ and $[1, \infty)$, we write

$$
F(s)=\int_{0}^{\infty} e^{-s t} f(t) d t=\int_{0}^{1} e^{-s t}(1) d t+\int_{1}^{\infty} e^{-s t}\left(-3 e^{-t}\right) d t
$$

Since

$$
\int_{0}^{1} e^{-s t} d t=\left\{\begin{array}{cl}
\frac{1-e^{-s}}{s}, & s \neq 0 \\
1, & s=0
\end{array}\right.
$$

and

$$
\int_{1}^{\infty} e^{-s t}\left(-3 e^{-t}\right) d t=-3 \int_{1}^{\infty} e^{-(s+1) t} d t=-\frac{3 e^{-(s+1)}}{s+1}, \quad s>-1
$$

it follows that

$$
F(s)=\left\{\begin{array}{cl}
\frac{1-e^{-s}}{s}-3 \frac{e^{-(s+1)}}{s+1}, & s>-1, s \neq 0 \\
1-\frac{3}{e}, & s=0
\end{array}\right.
$$

This is consistent with Theorem 8.1.6, since

$$
|f(t)| \leq 3 e^{-t}, \quad t \geq 1
$$

and therefore $f$ is of exponential order $s_{0}=-1$.
REMARK: In Section 8.4 we'll develop a more efficient method for finding Laplace transforms of piecewise continuous functions.

Example 8.1.11 We stated earlier that

$$
\int_{0}^{\infty} e^{-s t} e^{t^{2}} d t=\infty
$$

for all $s$, so Theorem 8.1.6 implies that $f(t)=e^{t^{2}}$ is not of exponential order, since

$$
\lim _{t \rightarrow \infty} \frac{e^{t^{2}}}{M e^{s_{0} t}}=\lim _{t \rightarrow \infty} \frac{1}{M} e^{t^{2}-s_{0} t}=\infty
$$

so

$$
e^{t^{2}}>M e^{s_{0} t}
$$

for sufficiently large values of $t$, for any choice of $M$ and $s_{0}$ (Exercise 3).

### 8.1 Exercises

1. Find the Laplace transforms of the following functions by evaluating the integral $F(s)=\int_{0}^{\infty} e^{-s t} f(t) d t$.
(a) $t$
(b) $t e^{-t}$
(d) $e^{2 t}-3 e^{t}$
(e) $t^{2}$
(c) $\sinh b t$
2. Use the table of Laplace transforms to find the Laplace transforms of the following functions.
(a) $\cosh t \sin t$
(b) $\sin ^{2} t$
(c) $\cos ^{2} 2 t$
(d) $\cosh ^{2} t$
(e) $t \sinh 2 t$
(f) $\sin t \cos t$
(g) $\sin \left(t+\frac{\pi}{4}\right)$
(h) $\cos 2 t-\cos 3 t$
(i) $\sin 2 t+\cos 4 t$
3. Show that

$$
\int_{0}^{\infty} e^{-s t} e^{t^{2}} d t=\infty
$$

for every real number $s$.
4. Graph the following piecewise continuous functions and evaluate $f(t+), f(t-)$, and $f(t)$ at each point of discontinuity.
(a) $f(t)=\left\{\begin{array}{cl}-t, & 0 \leq t<2, \\ t-4, & 2 \leq t<3, \\ 1, & t \geq 3 .\end{array}\right.$
(b) $f(t)=\left\{\begin{array}{cl}t^{2}+2, & 0 \leq t<1, \\ 4, & t=1, \\ t, & t>1 .\end{array}\right.$
(c) $f(t)=\left\{\begin{aligned} \sin t, & 0 \leq t<\pi / 2, \\ 2 \sin t, & \pi / 2 \leq t<\pi, \\ \cos t, & t \geq \pi .\end{aligned}\right.$
(d) $f(t)=\left\{\begin{array}{cl}t, & 0 \leq t<1, \\ 2, & t=1, \\ 2-t, & 1 \leq t<2, \\ 3, & t=2, \\ 6, & t>2 .\end{array}\right.$
5. Find the Laplace transform:
(a) $f(t)=\left\{\begin{array}{cl}e^{-t}, & 0 \leq t<1, \\ e^{-2 t}, & t \geq 1 .\end{array}\right.$
(b) $f(t)= \begin{cases}1, & 0 \leq t<4, \\ t, & t \geq 4 .\end{cases}$
(c) $f(t)= \begin{cases}t, & 0 \leq t<1, \\ 1, & t \geq 1 .\end{cases}$
(d) $f(t)=\left\{\begin{aligned} t e^{t}, & 0 \leq t<1, \\ e^{t}, & t \geq 1 .\end{aligned}\right.$
6. Prove that if $f(t) \leftrightarrow F(s)$ then $t^{k} f(t) \leftrightarrow(-1)^{k} F^{(k)}(s)$. Hint: Assume that it's permissible to differentiate the integral $\int_{0}^{\infty} e^{-s t} f(t) d t$ with respect to $s$ under the integral sign.
7. Use the known Laplace transforms

$$
\mathcal{L}\left(e^{\lambda t} \sin \omega t\right)=\frac{\omega}{(s-\lambda)^{2}+\omega^{2}} \quad \text { and } \quad \mathcal{L}\left(e^{\lambda t} \cos \omega t\right)=\frac{s-\lambda}{(s-\lambda)^{2}+\omega^{2}}
$$

and the result of Exercise 6 to find $\mathcal{L}\left(t e^{\lambda t} \cos \omega t\right)$ and $\mathcal{L}\left(t e^{\lambda t} \sin \omega t\right)$.
8. Use the known Laplace transform $\mathcal{L}(1)=1 / s$ and the result of Exercise 6 to show that

$$
\mathcal{L}\left(t^{n}\right)=\frac{n!}{s^{n+1}}, \quad n=\text { integer }
$$

9. (a) Show that if $\lim _{t \rightarrow \infty} e^{-s_{0} t} f(t)$ exists and is finite then $f$ is of exponential order $s_{0}$.
(b) Show that if $f$ is of exponential order $s_{0}$ then $\lim _{t \rightarrow \infty} e^{-s t} f(t)=0$ for all $s>s_{0}$.
(c) Show that if $f$ is of exponential order $s_{0}$ and $g(t)=f(t+\tau)$ where $\tau>0$, then $g$ is also of exponential order $s_{0}$.
10. Recall the next theorem from calculus.

Theorem A. Let $g$ be integrable on $[0, T]$ for every $T>0$. Suppose there's a function $w$ defined on some interval $[\tau, \infty)($ with $\tau \geq 0)$ such that $|g(t)| \leq w(t)$ for $t \geq \tau$ and $\int_{\tau}^{\infty} w(t) d t$ converges. Then $\int_{0}^{\infty} g(t) d t$ converges.
Use Theorem A to show that if $f$ is piecewise continuous on $[0, \infty)$ and of exponential order $s_{0}$, then $f$ has a Laplace transform $F(s)$ defined for $s>s_{0}$.
11. Prove: If $f$ is piecewise continuous and of exponential order then $\lim _{s \rightarrow \infty} F(s)=0$.
12. Prove: If $f$ is continuous on $[0, \infty)$ and of exponential order $s_{0}>0$, then

$$
\mathcal{L}\left(\int_{0}^{t} f(\tau) d \tau\right)=\frac{1}{s} \mathcal{L}(f), \quad s>s_{0}
$$

Hint: Use integration by parts to evaluate the transform on the left.
13. Suppose $f$ is piecewise continuous and of exponential order, and that $\lim _{t \rightarrow 0+} f(t) / t$ exists. Show that

$$
\mathcal{L}\left(\frac{f(t)}{t}\right)=\int_{s}^{\infty} F(r) d r .
$$

Hint: Use the results of Exercises 6 and 11.
14. Suppose $f$ is piecewise continuous on $[0, \infty)$.
(a) Prove: If the integral $g(t)=\int_{0}^{t} e^{-s_{0} \tau} f(\tau) d \tau$ satisfies the inequality $|g(t)| \leq M(t \geq 0)$, then $f$ has a Laplace transform $F(s)$ defined for $s>s_{0}$. Hint: Use integration by parts to show that

$$
\int_{0}^{T} e^{-s t} f(t) d t=e^{-\left(s-s_{0}\right) T} g(T)+\left(s-s_{0}\right) \int_{0}^{T} e^{-\left(s-s_{0}\right) t} g(t) d t
$$

(b) Show that if $\mathcal{L}(f)$ exists for $s=s_{0}$ then it exists for $s>s_{0}$. Show that the function

$$
f(t)=t e^{t^{2}} \cos \left(e^{t^{2}}\right)
$$

has a Laplace transform defined for $s>0$, even though $f$ isn't of exponential order.
(c) Show that the function

$$
f(t)=t e^{t^{2}} \cos \left(e^{t^{2}}\right)
$$

has a Laplace transform defined for $s>0$, even though $f$ isn't of exponential order.
15. Use the table of Laplace transforms and the result of Exercise 13 to find the Laplace transforms of the following functions.
(a) $\frac{\sin \omega t}{t} \quad(\omega>0)$
(b) $\frac{\cos \omega t-1}{t} \quad(\omega>0)$
(c) $\frac{e^{a t}-e^{b t}}{t}$
(d) $\frac{\cosh t-1}{t}$
(e) $\frac{\sinh ^{2} t}{t}$
16. The gamma function is defined by

$$
\Gamma(\alpha)=\int_{0}^{\infty} x^{\alpha-1} e^{-x} d x
$$

which can be shown to converge if $\alpha>0$.
(a) Use integration by parts to show that

$$
\Gamma(\alpha+1)=\alpha \Gamma(\alpha), \quad \alpha>0
$$

(b) Show that $\Gamma(n+1)=n$ ! if $n=1,2,3, \ldots$.
(c) From (b) and the table of Laplace transforms,

$$
\mathcal{L}\left(t^{\alpha}\right)=\frac{\Gamma(\alpha+1)}{s^{\alpha+1}}, \quad s>0
$$

if $\alpha$ is a nonnegative integer. Show that this formula is valid for any $\alpha>-1$. Hint: Change the variable of integration in the integral for $\Gamma(\alpha+1)$.
17. Suppose $f$ is continuous on $[0, T]$ and $f(t+T)=f(t)$ for all $t \geq 0$. (We say in this case that $f$ is periodic with period T.)
(a) Conclude from Theorem 8.1.6 that the Laplace transform of $f$ is defined for $s>0$. HINT: Since $f$ is continuous on $[0, T]$ and periodic with period $T$, it's bounded on $[0, \infty)$.
(b) (b) Show that

$$
F(s)=\frac{1}{1-e^{-s T}} \int_{0}^{T} e^{-s t} f(t) d t, \quad s>0
$$

Hint: Write

$$
F(s)=\sum_{n=0}^{\infty} \int_{n T}^{(n+1) T} e^{-s t} f(t) d t
$$

Then show that

$$
\int_{n T}^{(n+1) T} e^{-s t} f(t) d t=e^{-n s T} \int_{0}^{T} e^{-s t} f(t) d t
$$

and recall the formula for the sum of a geometric series.
18. Use the formula given in Exercise 17 (b) to find the Laplace transforms of the given periodic functions:
(a) $f(t)=\left\{\begin{array}{cc}t, & 0 \leq t<1, \\ 2-t, & 1 \leq t<2,\end{array} \quad f(t+2)=f(t), \quad t \geq 0\right.$
(b) $f(t)=\left\{\begin{array}{rl}1, & 0 \leq t<\frac{1}{2}, \\ -1, & \frac{1}{2} \leq t<1,\end{array} \quad f(t+1)=f(t), \quad t \geq 0\right.$
(c) $f(t)=|\sin t|$
(d) $f(t)=\left\{\begin{array}{cl}\sin t, & 0 \leq t<\pi, \\ 0, & \pi \leq t<2 \pi,\end{array} \quad f(t+2 \pi)=f(t)\right.$

### 8.2 THE INVERSE LAPLACE TRANSFORM

Definition of the Inverse Laplace Transform
In Section 8.1 we defined the Laplace transform of $f$ by

$$
F(s)=\mathcal{L}(f)=\int_{0}^{\infty} e^{-s t} f(t) d t
$$

We'll also say that $f$ is an inverse Laplace Transform of $F$, and write

$$
f=\mathcal{L}^{-1}(F)
$$

To solve differential equations with the Laplace transform, we must be able to obtain $f$ from its transform $F$. There's a formula for doing this, but we can't use it because it requires the theory of functions of a complex variable. Fortunately, we can use the table of Laplace transforms to find inverse transforms that we'll need.

Example 8.2.1 Use the table of Laplace transforms to find

$$
\text { (a) } \mathcal{L}^{-1}\left(\frac{1}{s^{2}-1}\right) \quad \text { and } \quad \text { (b) } \mathcal{L}^{-1}\left(\frac{s}{s^{2}+9}\right)
$$

$\underline{\text { SOLUTION(a) }}$ Setting $b=1$ in the transform pair

$$
\sinh b t \leftrightarrow \frac{b}{s^{2}-b^{2}}
$$

shows that

$$
\mathcal{L}^{-1}\left(\frac{1}{s^{2}-1}\right)=\sinh t
$$

SOLUTION(b) Setting $\omega=3$ in the transform pair

$$
\cos \omega t \leftrightarrow \frac{s}{s^{2}+\omega^{2}}
$$

shows that

$$
\mathcal{L}^{-1}\left(\frac{s}{s^{2}+9}\right)=\cos 3 t
$$

The next theorem enables us to find inverse transforms of linear combinations of transforms in the table. We omit the proof.

Theorem 8.2.1 [Linearity Property] If $F_{1}, F_{2}, \ldots, F_{n}$ are Laplace transforms and $c_{1}, c_{2}, \ldots, c_{n}$ are constants, then

$$
\mathcal{L}^{-1}\left(c_{1} F_{1}+c_{2} F_{2}+\cdots+c_{n} F_{n}\right)=c_{1} \mathcal{L}^{-1}\left(F_{1}\right)+c_{2} \mathcal{L}^{-1}\left(F_{2}\right)+\cdots+c_{n} \mathcal{L}^{-1} F_{n}
$$

## Example 8.2.2 Find

$$
\mathcal{L}^{-1}\left(\frac{8}{s+5}+\frac{7}{s^{2}+3}\right)
$$

Solution From the table of Laplace transforms in Section 8.8,,

$$
e^{a t} \leftrightarrow \frac{1}{s-a} \quad \text { and } \quad \sin \omega t \leftrightarrow \frac{\omega}{s^{2}+\omega^{2}}
$$

Theorem 8.2.1 with $a=-5$ and $\omega=\sqrt{3}$ yields

$$
\begin{aligned}
\mathcal{L}^{-1}\left(\frac{8}{s+5}+\frac{7}{s^{2}+3}\right) & =8 \mathcal{L}^{-1}\left(\frac{1}{s+5}\right)+7 \mathcal{L}^{-1}\left(\frac{1}{s^{2}+3}\right) \\
& =8 \mathcal{L}^{-1}\left(\frac{1}{s+5}\right)+\frac{7}{\sqrt{3}} \mathcal{L}^{-1}\left(\frac{\sqrt{3}}{s^{2}+3}\right) \\
& =8 e^{-5 t}+\frac{7}{\sqrt{3}} \sin \sqrt{3} t
\end{aligned}
$$

## Example 8.2.3 Find

$$
\mathcal{L}^{-1}\left(\frac{3 s+8}{s^{2}+2 s+5}\right)
$$

Solution Completing the square in the denominator yields

$$
\frac{3 s+8}{s^{2}+2 s+5}=\frac{3 s+8}{(s+1)^{2}+4}
$$

Because of the form of the denominator, we consider the transform pairs

$$
e^{-t} \cos 2 t \leftrightarrow \frac{s+1}{(s+1)^{2}+4} \quad \text { and } \quad e^{-t} \sin 2 t \leftrightarrow \frac{2}{(s+1)^{2}+4},
$$

and write

$$
\begin{aligned}
\mathcal{L}^{-1}\left(\frac{3 s+8}{(s+1)^{2}+4}\right) & =\mathcal{L}^{-1}\left(\frac{3 s+3}{(s+1)^{2}+4}\right)+\mathcal{L}^{-1}\left(\frac{5}{(s+1)^{2}+4}\right) \\
& =3 \mathcal{L}^{-1}\left(\frac{s+1}{(s+1)^{2}+4}\right)+\frac{5}{2} \mathcal{L}^{-1}\left(\frac{2}{(s+1)^{2}+4}\right) \\
& =e^{-t}\left(3 \cos 2 t+\frac{5}{2} \sin 2 t\right)
\end{aligned}
$$

REMARK: We'll often write inverse Laplace transforms of specific functions without explicitly stating how they are obtained. In such cases you should refer to the table of Laplace transforms in Section 8.8.

Inverse Laplace Transforms of Rational Functions
Using the Laplace transform to solve differential equations often requires finding the inverse transform of a rational function

$$
F(s)=\frac{P(s)}{Q(s)}
$$

where $P$ and $Q$ are polynomials in $s$ with no common factors. Since it can be shown that $\lim _{s \rightarrow \infty} F(s)=$ 0 if $F$ is a Laplace transform, we need only consider the case where degree $(P)<\operatorname{degree}(Q)$. To obtain $\mathcal{L}^{-1}(F)$, we find the partial fraction expansion of $F$, obtain inverse transforms of the individual terms in the expansion from the table of Laplace transforms, and use the linearity property of the inverse transform. The next two examples illustrate this.

Example 8.2.4 Find the inverse Laplace transform of

$$
\begin{equation*}
F(s)=\frac{3 s+2}{s^{2}-3 s+2} \tag{8.2.1}
\end{equation*}
$$

Solution (METHOD 1) Factoring the denominator in (8.2.1) yields

$$
\begin{equation*}
F(s)=\frac{3 s+2}{(s-1)(s-2)} \tag{8.2.2}
\end{equation*}
$$

The form for the partial fraction expansion is

$$
\begin{equation*}
\frac{3 s+2}{(s-1)(s-2)}=\frac{A}{s-1}+\frac{B}{s-2} . \tag{8.2.3}
\end{equation*}
$$

Multiplying this by $(s-1)(s-2)$ yields

$$
3 s+2=(s-2) A+(s-1) B
$$

Setting $s=2$ yields $B=8$ and setting $s=1$ yields $A=-5$. Therefore

$$
F(s)=-\frac{5}{s-1}+\frac{8}{s-2}
$$

and

$$
\mathcal{L}^{-1}(F)=-5 \mathcal{L}^{-1}\left(\frac{1}{s-1}\right)+8 \mathcal{L}^{-1}\left(\frac{1}{s-2}\right)=-5 e^{t}+8 e^{2 t}
$$

Solution (METHOD 2) We don't really have to multiply (8.2.3) by $(s-1)(s-2)$ to compute $A$ and $B$. We can obtain $A$ by simply ignoring the factor $s-1$ in the denominator of (8.2.2) and setting $s=1$ elsewhere; thus,

$$
\begin{equation*}
A=\left.\frac{3 s+2}{s-2}\right|_{s=1}=\frac{3 \cdot 1+2}{1-2}=-5 \tag{8.2.4}
\end{equation*}
$$

Similarly, we can obtain $B$ by ignoring the factor $s-2$ in the denominator of (8.2.2) and setting $s=2$ elsewhere; thus,

$$
\begin{equation*}
B=\left.\frac{3 s+2}{s-1}\right|_{s=2}=\frac{3 \cdot 2+2}{2-1}=8 . \tag{8.2.5}
\end{equation*}
$$

To justify this, we observe that multiplying (8.2.3) by $s-1$ yields

$$
\frac{3 s+2}{s-2}=A+(s-1) \frac{B}{s-2}
$$

and setting $s=1$ leads to (8.2.4). Similarly, multiplying (8.2.3) by $s-2$ yields

$$
\frac{3 s+2}{s-1}=(s-2) \frac{A}{s-2}+B
$$

and setting $s=2$ leads to (8.2.5). (It isn't necesary to write the last two equations. We wrote them only to justify the shortcut procedure indicated in (8.2.4) and (8.2.5).)
The shortcut employed in the second solution of Example 8.2.4 is Heaviside's method. The next theorem states this method formally. For a proof and an extension of this theorem, see Exercise 10.

Theorem 8.2.2 Suppose

$$
\begin{equation*}
F(s)=\frac{P(s)}{\left(s-s_{1}\right)\left(s-s_{2}\right) \cdots\left(s-s_{n}\right)} \tag{8.2.6}
\end{equation*}
$$

where $s_{1}, s_{2}, \ldots, s_{n}$ are distinct and $P$ is a polynomial of degree less than $n$. Then

$$
F(s)=\frac{A_{1}}{s-s_{1}}+\frac{A_{2}}{s-s_{2}}+\cdots+\frac{A_{n}}{s-s_{n}}
$$

where $A_{i}$ can be computed from (8.2.6) by ignoring the factor $s-s_{i}$ and setting $s=s_{i}$ elsewhere.

Example 8.2.5 Find the inverse Laplace transform of

$$
\begin{equation*}
F(s)=\frac{6+(s+1)\left(s^{2}-5 s+11\right)}{s(s-1)(s-2)(s+1)} \tag{8.2.7}
\end{equation*}
$$

Solution The partial fraction expansion of (8.2.7) is of the form

$$
\begin{equation*}
F(s)=\frac{A}{s}+\frac{B}{s-1}+\frac{C}{s-2}+\frac{D}{s+1} \tag{8.2.8}
\end{equation*}
$$

To find $A$, we ignore the factor $s$ in the denominator of (8.2.7) and set $s=0$ elsewhere. This yields

$$
A=\frac{6+(1)(11)}{(-1)(-2)(1)}=\frac{17}{2}
$$

Similarly, the other coefficients are given by

$$
\begin{gathered}
B=\frac{6+(2)(7)}{(1)(-1)(2)}=-10 \\
C=\frac{6+3(5)}{2(1)(3)}=\frac{7}{2}
\end{gathered}
$$

and

$$
D=\frac{6}{(-1)(-2)(-3)}=-1
$$

Therefore

$$
F(s)=\frac{17}{2} \frac{1}{s}-\frac{10}{s-1}+\frac{7}{2} \frac{1}{s-2}-\frac{1}{s+1}
$$

and

$$
\begin{aligned}
\mathcal{L}^{-1}(F) & =\frac{17}{2} \mathcal{L}^{-1}\left(\frac{1}{s}\right)-10 \mathcal{L}^{-1}\left(\frac{1}{s-1}\right)+\frac{7}{2} \mathcal{L}^{-1}\left(\frac{1}{s-2}\right)-\mathcal{L}^{-1}\left(\frac{1}{s+1}\right) \\
& =\frac{17}{2}-10 e^{t}+\frac{7}{2} e^{2 t}-e^{-t}
\end{aligned}
$$

REMARK: We didn't "multiply out" the numerator in (8.2.7) before computing the coefficients in (8.2.8), since it wouldn't simplify the computations.
Example 8.2.6 Find the inverse Laplace transform of

$$
\begin{equation*}
F(s)=\frac{8-(s+2)(4 s+10)}{(s+1)(s+2)^{2}} \tag{8.2.9}
\end{equation*}
$$

Solution The form for the partial fraction expansion is

$$
\begin{equation*}
F(s)=\frac{A}{s+1}+\frac{B}{s+2}+\frac{C}{(s+2)^{2}} \tag{8.2.10}
\end{equation*}
$$

Because of the repeated factor $(s+2)^{2}$ in (8.2.9), Heaviside's method doesn't work. Instead, we find a common denominator in (8.2.10). This yields

$$
\begin{equation*}
F(s)=\frac{A(s+2)^{2}+B(s+1)(s+2)+C(s+1)}{(s+1)(s+2)^{2}} \tag{8.2.11}
\end{equation*}
$$

If (8.2.9) and (8.2.11) are to be equivalent, then

$$
\begin{equation*}
A(s+2)^{2}+B(s+1)(s+2)+C(s+1)=8-(s+2)(4 s+10) \tag{8.2.12}
\end{equation*}
$$

The two sides of this equation are polynomials of degree two. From a theorem of algebra, they will be equal for all $s$ if they are equal for any three distinct values of $s$. We may determine $A, B$ and $C$ by choosing convenient values of $s$.

The left side of (8.2.12) suggests that we take $s=-2$ to obtain $C=-8$, and $s=-1$ to obtain $A=2$. We can now choose any third value of $s$ to determine $B$. Taking $s=0$ yields $4 A+2 B+C=-12$. Since $A=2$ and $C=-8$ this implies that $B=-6$. Therefore

$$
F(s)=\frac{2}{s+1}-\frac{6}{s+2}-\frac{8}{(s+2)^{2}}
$$

and

$$
\begin{aligned}
\mathcal{L}^{-1}(F) & =2 \mathcal{L}^{-1}\left(\frac{1}{s+1}\right)-6 \mathcal{L}^{-1}\left(\frac{1}{s+2}\right)-8 \mathcal{L}^{-1}\left(\frac{1}{(s+2)^{2}}\right) \\
& =2 e^{-t}-6 e^{-2 t}-8 t e^{-2 t}
\end{aligned}
$$

Example 8.2.7 Find the inverse Laplace transform of

$$
F(s)=\frac{s^{2}-5 s+7}{(s+2)^{3}}
$$

Solution The form for the partial fraction expansion is

$$
F(s)=\frac{A}{s+2}+\frac{B}{(s+2)^{2}}+\frac{C}{(s+2)^{3}} .
$$

The easiest way to obtain $A, B$, and $C$ is to expand the numerator in powers of $s+2$. This yields

$$
s^{2}-5 s+7=[(s+2)-2]^{2}-5[(s+2)-2]+7=(s+2)^{2}-9(s+2)+21 .
$$

Therefore

$$
\begin{aligned}
F(s) & =\frac{(s+2)^{2}-9(s+2)+21}{(s+2)^{3}} \\
& =\frac{1}{s+2}-\frac{9}{(s+2)^{2}}+\frac{21}{(s+2)^{3}}
\end{aligned}
$$

and

$$
\begin{aligned}
\mathcal{L}^{-1}(F) & =\mathcal{L}^{-1}\left(\frac{1}{s+2}\right)-9 \mathcal{L}^{-1}\left(\frac{1}{(s+2)^{2}}\right)+\frac{21}{2} \mathcal{L}^{-1}\left(\frac{2}{(s+2)^{3}}\right) \\
& =e^{-2 t}\left(1-9 t+\frac{21}{2} t^{2}\right)
\end{aligned}
$$

Example 8.2.8 Find the inverse Laplace transform of

$$
\begin{equation*}
F(s)=\frac{1-s(5+3 s)}{s\left[(s+1)^{2}+1\right]} \tag{8.2.13}
\end{equation*}
$$

Solution One form for the partial fraction expansion of $F$ is

$$
\begin{equation*}
F(s)=\frac{A}{s}+\frac{B s+C}{(s+1)^{2}+1} \tag{8.2.14}
\end{equation*}
$$

However, we see from the table of Laplace transforms that the inverse transform of the second fraction on the right of (8.2.14) will be a linear combination of the inverse transforms

$$
e^{-t} \cos t \text { and } e^{-t} \sin t
$$

of

$$
\frac{s+1}{(s+1)^{2}+1} \quad \text { and } \quad \frac{1}{(s+1)^{2}+1}
$$

respectively. Therefore, instead of (8.2.14) we write

$$
\begin{equation*}
F(s)=\frac{A}{s}+\frac{B(s+1)+C}{(s+1)^{2}+1} \tag{8.2.15}
\end{equation*}
$$

Finding a common denominator yields

$$
\begin{equation*}
F(s)=\frac{A\left[(s+1)^{2}+1\right]+B(s+1) s+C s}{s\left[(s+1)^{2}+1\right]} . \tag{8.2.16}
\end{equation*}
$$

If (8.2.13) and (8.2.16) are to be equivalent, then

$$
A\left[(s+1)^{2}+1\right]+B(s+1) s+C s=1-s(5+3 s)
$$

This is true for all $s$ if it's true for three distinct values of $s$. Choosing $s=0,-1$, and 1 yields the system

$$
\begin{aligned}
2 A & =1 \\
A-C & =3 \\
5 A+2 B+C & =-7 .
\end{aligned}
$$

Solving this system yields

$$
A=\frac{1}{2}, \quad B=-\frac{7}{2}, \quad C=-\frac{5}{2} .
$$

Hence, from (8.2.15),

$$
F(s)=\frac{1}{2 s}-\frac{7}{2} \frac{s+1}{(s+1)^{2}+1}-\frac{5}{2} \frac{1}{(s+1)^{2}+1}
$$

Therefore

$$
\begin{aligned}
\mathcal{L}^{-1}(F) & =\frac{1}{2} \mathcal{L}^{-1}\left(\frac{1}{s}\right)-\frac{7}{2} \mathcal{L}^{-1}\left(\frac{s+1}{(s+1)^{2}+1}\right)-\frac{5}{2} \mathcal{L}^{-1}\left(\frac{1}{(s+1)^{2}+1}\right) \\
& =\frac{1}{2}-\frac{7}{2} e^{-t} \cos t-\frac{5}{2} e^{-t} \sin t
\end{aligned}
$$

Example 8.2.9 Find the inverse Laplace transform of

$$
\begin{equation*}
F(s)=\frac{8+3 s}{\left(s^{2}+1\right)\left(s^{2}+4\right)} \tag{8.2.17}
\end{equation*}
$$

Solution The form for the partial fraction expansion is

$$
F(s)=\frac{A+B s}{s^{2}+1}+\frac{C+D s}{s^{2}+4}
$$

The coefficients $A, B, C$ and $D$ can be obtained by finding a common denominator and equating the resulting numerator to the numerator in (8.2.17). However, since there's no first power of $s$ in the denominator of (8.2.17), there's an easier way: the expansion of

$$
F_{1}(s)=\frac{1}{\left(s^{2}+1\right)\left(s^{2}+4\right)}
$$

can be obtained quickly by using Heaviside's method to expand

$$
\frac{1}{(x+1)(x+4)}=\frac{1}{3}\left(\frac{1}{x+1}-\frac{1}{x+4}\right)
$$

and then setting $x=s^{2}$ to obtain

$$
\frac{1}{\left(s^{2}+1\right)\left(s^{2}+4\right)}=\frac{1}{3}\left(\frac{1}{s^{2}+1}-\frac{1}{s^{2}+4}\right) .
$$

Multiplying this by $8+3 s$ yields

$$
F(s)=\frac{8+3 s}{\left(s^{2}+1\right)\left(s^{2}+4\right)}=\frac{1}{3}\left(\frac{8+3 s}{s^{2}+1}-\frac{8+3 s}{s^{2}+4}\right) .
$$

Therefore

$$
\mathcal{L}^{-1}(F)=\frac{8}{3} \sin t+\cos t-\frac{4}{3} \sin 2 t-\cos 2 t .
$$

## USING TECHNOLOGY

Some software packages that do symbolic algebra can find partial fraction expansions very easily. We recommend that you use such a package if one is available to you, but only after you've done enough partial fraction expansions on your own to master the technique.

### 8.2 Exercises

1. Use the table of Laplace transforms to find the inverse Laplace transform.
(a) $\frac{3}{(s-7)^{4}}$
(b) $\frac{2 s-4}{s^{2}-4 s+13}$
(c) $\frac{1}{s^{2}+4 s+20}$
(d) $\frac{2}{s^{2}+9}$
(e) $\frac{s^{2}-1}{\left(s^{2}+1\right)^{2}}$
(f) $\frac{1}{(s-2)^{2}-4}$
(g) $\frac{12 s-24}{\left(s^{2}-4 s+85\right)^{2}}$
(h) $\frac{2}{(s-3)^{2}-9}$
(i) $\frac{s^{2}-4 s+3}{\left(s^{2}-4 s+5\right)^{2}}$
2. Use Theorem 8.2.1 and the table of Laplace transforms to find the inverse Laplace transform.
(a) $\frac{2 s+3}{(s-7)^{4}}$
(b) $\frac{s^{2}-1}{(s-2)^{6}}$
(c) $\frac{s+5}{s^{2}+6 s+18}$
(d) $\frac{2 s+1}{s^{2}+9}$
(e) $\frac{s}{s^{2}+2 s+1}$
(f) $\frac{s+1}{s^{2}-9}$
(g) $\frac{s^{3}+2 s^{2}-s-3}{(s+1)^{4}}$
(h) $\frac{2 s+3}{(s-1)^{2}+4}$
(i) $\frac{1}{s}-\frac{s}{s^{2}+1}$
(j) $\frac{3 s+4}{s^{2}-1}$
(k) $\frac{3}{s-1}+\frac{4 s+1}{s^{2}+9}$
(l) $\frac{3}{(s+2)^{2}}-\frac{2 s+6}{s^{2}+4}$
3. Use Heaviside's method to find the inverse Laplace transform.
(a) $\frac{3-(s+1)(s-2)}{(s+1)(s+2)(s-2)}$
(b) $\frac{7+(s+4)(18-3 s)}{(s-3)(s-1)(s+4)}$
(c) $\frac{2+(s-2)(3-2 s)}{(s-2)(s+2)(s-3)}$
(d) $\frac{3-(s-1)(s+1)}{(s+4)(s-2)(s-1)}$
(e) $\frac{3+(s-2)\left(10-2 s-s^{2}\right)}{(s-2)(s+2)(s-1)(s+3)}$
(f) $\frac{3+(s-3)\left(2 s^{2}+s-21\right)}{(s-3)(s-1)(s+4)(s-2)}$
4. Find the inverse Laplace transform.
(a) $\frac{2+3 s}{\left(s^{2}+1\right)(s+2)(s+1)}$
(b) $\frac{3 s^{2}+2 s+1}{\left(s^{2}+1\right)\left(s^{2}+2 s+2\right)}$
(c) $\frac{3 s+2}{(s-2)\left(s^{2}+2 s+5\right)}$
(d) $\frac{3 s^{2}+2 s+1}{(s-1)^{2}(s+2)(s+3)}$
(e) $\frac{2 s^{2}+s+3}{(s-1)^{2}(s+2)^{2}}$
(f) $\frac{3 s+2}{\left(s^{2}+1\right)(s-1)^{2}}$
5. Use the method of Example 8.2.9 to find the inverse Laplace transform.
(a) $\frac{3 s+2}{\left(s^{2}+4\right)\left(s^{2}+9\right)}$
(b) $\frac{-4 s+1}{\left(s^{2}+1\right)\left(s^{2}+16\right)}$
(c) $\frac{5 s+3}{\left(s^{2}+1\right)\left(s^{2}+4\right)}$
(d) $\frac{-s+1}{\left(4 s^{2}+1\right)\left(s^{2}+1\right)}$
(e) $\frac{17 s-34}{\left(s^{2}+16\right)\left(16 s^{2}+1\right)}$
(f) $\frac{2 s-1}{\left(4 s^{2}+1\right)\left(9 s^{2}+1\right)}$
6. Find the inverse Laplace transform.
(a) $\frac{17 s-15}{\left(s^{2}-2 s+5\right)\left(s^{2}+2 s+10\right)}$
(b) $\frac{8 s+56}{\left(s^{2}-6 s+13\right)\left(s^{2}+2 s+5\right)}$
(c) $\frac{s+9}{\left(s^{2}+4 s+5\right)\left(s^{2}-4 s+13\right)}$
(d) $\frac{3 s-2}{\left(s^{2}-4 s+5\right)\left(s^{2}-6 s+13\right)}$
(e) $\frac{3 s-1}{\left(s^{2}-2 s+2\right)\left(s^{2}+2 s+5\right)}$
(f) $\frac{20 s+40}{\left(4 s^{2}-4 s+5\right)\left(4 s^{2}+4 s+5\right)}$
7. Find the inverse Laplace transform.
(a) $\frac{1}{s\left(s^{2}+1\right)}$
(b) $\frac{1}{(s-1)\left(s^{2}-2 s+17\right)}$
(c) $\frac{3 s+2}{(s-2)\left(s^{2}+2 s+10\right)}$
(d) $\frac{34-17 s}{(2 s-1)\left(s^{2}-2 s+5\right)}$
(e) $\frac{s+2}{(s-3)\left(s^{2}+2 s+5\right)}$
(f) $\frac{2 s-2}{(s-2)\left(s^{2}+2 s+10\right)}$
8. Find the inverse Laplace transform.
(a) $\frac{2 s+1}{\left(s^{2}+1\right)(s-1)(s-3)}$
(b) $\frac{s+2}{\left(s^{2}+2 s+2\right)\left(s^{2}-1\right)}$
(c) $\frac{2 s-1}{\left(s^{2}-2 s+2\right)(s+1)(s-2)}$
(d) $\frac{s-6}{\left(s^{2}-1\right)\left(s^{2}+4\right)}$
(e) $\frac{2 s-3}{s(s-2)\left(s^{2}-2 s+5\right)}$
(f) $\frac{5 s-15}{\left(s^{2}-4 s+13\right)(s-2)(s-1)}$
9. Given that $f(t) \leftrightarrow F(s)$, find the inverse Laplace transform of $F(a s-b)$, where $a>0$.
10. (a) If $s_{1}, s_{2}, \ldots, s_{n}$ are distinct and $P$ is a polynomial of degree less than $n$, then

$$
\frac{P(s)}{\left(s-s_{1}\right)\left(s-s_{2}\right) \cdots\left(s-s_{n}\right)}=\frac{A_{1}}{s-s_{1}}+\frac{A_{2}}{s-s_{2}}+\cdots+\frac{A_{n}}{s-s_{n}} .
$$

Multiply through by $s-s_{i}$ to show that $A_{i}$ can be obtained by ignoring the factor $s-s_{i}$ on the left and setting $s=s_{i}$ elsewhere.
(b) Suppose $P$ and $Q_{1}$ are polynomials such that degree $(P) \leq \operatorname{degree}\left(Q_{1}\right)$ and $Q_{1}\left(s_{1}\right) \neq 0$. Show that the coefficient of $1 /\left(s-s_{1}\right)$ in the partial fraction expansion of

$$
F(s)=\frac{P(s)}{\left(s-s_{1}\right) Q_{1}(s)}
$$

is $P\left(s_{1}\right) / Q_{1}\left(s_{1}\right)$.
(c) Explain how the results of (a) and (b) are related.

### 8.3 SOLUTION OF INITIAL VALUE PROBLEMS

Laplace Transforms of Derivatives
In the rest of this chapter we'll use the Laplace transform to solve initial value problems for constant coefficient second order equations. To do this, we must know how the Laplace transform of $f^{\prime}$ is related to the Laplace transform of $f$. The next theorem answers this question.

Theorem 8.3.1 Suppose $f$ is continuous on $[0, \infty)$ and of exponential order $s_{0}$, and $f^{\prime}$ is piecewise continuous on $[0, \infty)$. Then $f$ and $f^{\prime}$ have Laplace transforms for $s>s_{0}$, and

$$
\begin{equation*}
\mathcal{L}\left(f^{\prime}\right)=s \mathcal{L}(f)-f(0) . \tag{8.3.1}
\end{equation*}
$$

## Proof

We know from Theorem 8.1.6 that $\mathcal{L}(f)$ is defined for $s>s_{0}$. We first consider the case where $f^{\prime}$ is continuous on $[0, \infty)$. Integration by parts yields

$$
\begin{align*}
\int_{0}^{T} e^{-s t} f^{\prime}(t) d t & =\left.e^{-s t} f(t)\right|_{0} ^{T}+s \int_{0}^{T} e^{-s t} f(t) d t  \tag{8.3.2}\\
& =e^{-s T} f(T)-f(0)+s \int_{0}^{T} e^{-s t} f(t) d t
\end{align*}
$$

for any $T>0$. Since $f$ is of exponential order $s_{0}, \lim _{T \rightarrow \infty} e^{-s T} f(T)=0$ and the last integral in (8.3.2) converges as $T \rightarrow \infty$ if $s>s_{0}$. Therefore

$$
\begin{aligned}
\int_{0}^{\infty} e^{-s t} f^{\prime}(t) d t & =-f(0)+s \int_{0}^{\infty} e^{-s t} f(t) d t \\
& =-f(0)+s \mathcal{L}(f)
\end{aligned}
$$

which proves (8.3.1). Now suppose $T>0$ and $f^{\prime}$ is only piecewise continuous on $[0, T]$, with discontinuities at $t_{1}<t_{2}<\cdots<t_{n-1}$. For convenience, let $t_{0}=0$ and $t_{n}=T$. Integrating by parts yields

$$
\begin{aligned}
\int_{t_{i-1}}^{t_{i}} e^{-s t} f^{\prime}(t) d t & =\left.e^{-s t} f(t)\right|_{t_{i-1}} ^{t_{i}}+s \int_{t_{i-1}}^{t_{i}} e^{-s t} f(t) d t \\
& =e^{-s t_{i}} f\left(t_{i}\right)-e^{-s t_{i-1}} f\left(t_{i-1}\right)+s \int_{t_{i-1}}^{t_{i}} e^{-s t} f(t) d t
\end{aligned}
$$

Summing both sides of this equation from $i=1$ to $n$ and noting that

$$
\begin{gathered}
\left(e^{-s t_{1}} f\left(t_{1}\right)-e^{-s t_{0}} f\left(t_{0}\right)\right)+\left(e^{-s t_{2}} f\left(t_{2}\right)-e^{-s t_{1}} f\left(t_{1}\right)\right)+\cdots+\left(e^{-s t_{N}} f\left(t_{N}\right)-e^{-s t_{N-1}} f\left(t_{N-1}\right)\right) \\
=e^{-s t_{N}} f\left(t_{N}\right)-e^{-s t_{0}} f\left(t_{0}\right)=e^{-s T} f(T)-f(0)
\end{gathered}
$$

yields (8.3.2), so (8.3.1) follows as before.
Example 8.3.1 In Example 8.1.4 we saw that

$$
\mathcal{L}(\cos \omega t)=\frac{s}{s^{2}+\omega^{2}}
$$

Applying (8.3.1) with $f(t)=\cos \omega t$ shows that

$$
\mathcal{L}(-\omega \sin \omega t)=s \frac{s}{s^{2}+\omega^{2}}-1=-\frac{\omega^{2}}{s^{2}+\omega^{2}}
$$

Therefore

$$
\mathcal{L}(\sin \omega t)=\frac{\omega}{s^{2}+\omega^{2}}
$$

which agrees with the corresponding result obtained in 8.1.4.
In Section 2.1 we showed that the solution of the initial value problem

$$
\begin{equation*}
y^{\prime}=a y, \quad y(0)=y_{0} \tag{8.3.3}
\end{equation*}
$$

is $y=y_{0} e^{a t}$. We'll now obtain this result by using the Laplace transform.
Let $Y(s)=\mathcal{L}(y)$ be the Laplace transform of the unknown solution of (8.3.3). Taking Laplace transforms of both sides of (8.3.3) yields

$$
\mathcal{L}\left(y^{\prime}\right)=\mathcal{L}(a y)
$$

which, by Theorem 8.3.1, can be rewritten as

$$
s \mathcal{L}(y)-y(0)=a \mathcal{L}(y)
$$

or

$$
s Y(s)-y_{0}=a Y(s)
$$

Solving for $Y(s)$ yields

$$
Y(s)=\frac{y_{0}}{s-a}
$$

so

$$
y=\mathcal{L}^{-1}(Y(s))=\mathcal{L}^{-1}\left(\frac{y_{0}}{s-a}\right)=y_{0} \mathcal{L}^{-1}\left(\frac{1}{s-a}\right)=y_{0} e^{a t}
$$

which agrees with the known result.
We need the next theorem to solve second order differential equations using the Laplace transform.
Theorem 8.3.2 Suppose $f$ and $f^{\prime}$ are continuous on $[0, \infty)$ and of exponential order $s_{0}$, and that $f^{\prime \prime}$ is piecewise continuous on $[0, \infty)$. Then $f, f^{\prime}$, and $f^{\prime \prime}$ have Laplace transforms for $s>s_{0}$,

$$
\begin{equation*}
\mathcal{L}\left(f^{\prime}\right)=s \mathcal{L}(f)-f(0), \tag{8.3.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{L}\left(f^{\prime \prime}\right)=s^{2} \mathcal{L}(f)-f^{\prime}(0)-s f(0) \tag{8.3.5}
\end{equation*}
$$

Proof Theorem 8.3.1 implies that $\mathcal{L}\left(f^{\prime}\right)$ exists and satisfies (8.3.4) for $s>s_{0}$. To prove that $\mathcal{L}\left(f^{\prime \prime}\right)$ exists and satisfies (8.3.5) for $s>s_{0}$, we first apply Theorem 8.3.1 to $g=f^{\prime}$. Since $g$ satisfies the hypotheses of Theorem 8.3.1, we conclude that $\mathcal{L}\left(g^{\prime}\right)$ is defined and satisfies

$$
\mathcal{L}\left(g^{\prime}\right)=s \mathcal{L}(g)-g(0)
$$

for $s>s_{0}$. However, since $g^{\prime}=f^{\prime \prime}$, this can be rewritten as

$$
\mathcal{L}\left(f^{\prime \prime}\right)=s \mathcal{L}\left(f^{\prime}\right)-f^{\prime}(0)
$$

Substituting (8.3.4) into this yields (8.3.5).
Solving Second Order Equations with the Laplace Transform
We'll now use the Laplace transform to solve initial value problems for second order equations.
Example 8.3.2 Use the Laplace transform to solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}-6 y^{\prime}+5 y=3 e^{2 t}, \quad y(0)=2, \quad y^{\prime}(0)=3 . \tag{8.3.6}
\end{equation*}
$$

Solution Taking Laplace transforms of both sides of the differential equation in (8.3.6) yields

$$
\mathcal{L}\left(y^{\prime \prime}-6 y^{\prime}+5 y\right)=\mathcal{L}\left(3 e^{2 t}\right)=\frac{3}{s-2}
$$

which we rewrite as

$$
\begin{equation*}
\mathcal{L}\left(y^{\prime \prime}\right)-6 \mathcal{L}\left(y^{\prime}\right)+5 \mathcal{L}(y)=\frac{3}{s-2} \tag{8.3.7}
\end{equation*}
$$

Now denote $\mathcal{L}(y)=Y(s)$. Theorem 8.3.2 and the initial conditions in (8.3.6) imply that

$$
\mathcal{L}\left(y^{\prime}\right)=s Y(s)-y(0)=s Y(s)-2
$$

and

$$
\mathcal{L}\left(y^{\prime \prime}\right)=s^{2} Y(s)-y^{\prime}(0)-s y(0)=s^{2} Y(s)-3-2 s
$$

Substituting from the last two equations into (8.3.7) yields

$$
\left(s^{2} Y(s)-3-2 s\right)-6(s Y(s)-2)+5 Y(s)=\frac{3}{s-2}
$$

Therefore

$$
\begin{equation*}
\left(s^{2}-6 s+5\right) Y(s)=\frac{3}{s-2}+(3+2 s)+6(-2) \tag{8.3.8}
\end{equation*}
$$

so

$$
(s-5)(s-1) Y(s)=\frac{3+(s-2)(2 s-9)}{s-2}
$$

and

$$
Y(s)=\frac{3+(s-2)(2 s-9)}{(s-2)(s-5)(s-1)}
$$

Heaviside's method yields the partial fraction expansion

$$
Y(s)=-\frac{1}{s-2}+\frac{1}{2} \frac{1}{s-5}+\frac{5}{2} \frac{1}{s-1}
$$

and taking the inverse transform of this yields

$$
y=-e^{2 t}+\frac{1}{2} e^{5 t}+\frac{5}{2} e^{t}
$$

as the solution of (8.3.6).
It isn't necessary to write all the steps that we used to obtain (8.3.8). To see how to avoid this, let's apply the method of Example 8.3.2 to the general initial value problem

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} . \tag{8.3.9}
\end{equation*}
$$

Taking Laplace transforms of both sides of the differential equation in (8.3.9) yields

$$
\begin{equation*}
a \mathcal{L}\left(y^{\prime \prime}\right)+b \mathcal{L}\left(y^{\prime}\right)+c \mathcal{L}(y)=F(s) \tag{8.3.10}
\end{equation*}
$$

Now let $Y(s)=\mathcal{L}(y)$. Theorem 8.3.2 and the initial conditions in (8.3.9) imply that

$$
\mathcal{L}\left(y^{\prime}\right)=s Y(s)-k_{0} \quad \text { and } \quad \mathcal{L}\left(y^{\prime \prime}\right)=s^{2} Y(s)-k_{1}-k_{0} s
$$

Substituting these into (8.3.10) yields

$$
\begin{equation*}
a\left(s^{2} Y(s)-k_{1}-k_{0} s\right)+b\left(s Y(s)-k_{0}\right)+c Y(s)=F(s) \tag{8.3.11}
\end{equation*}
$$

The coefficient of $Y(s)$ on the left is the characteristic polynomial

$$
p(s)=a s^{2}+b s+c
$$

of the complementary equation for (8.3.9). Using this and moving the terms involving $k_{0}$ and $k_{1}$ to the right side of (8.3.11) yields

$$
\begin{equation*}
p(s) Y(s)=F(s)+a\left(k_{1}+k_{0} s\right)+b k_{0} . \tag{8.3.12}
\end{equation*}
$$

This equation corresponds to (8.3.8) of Example 8.3.2. Having established the form of this equation in the general case, it is preferable to go directly from the initial value problem to this equation. You may find it easier to remember (8.3.12) rewritten as

$$
\begin{equation*}
p(s) Y(s)=F(s)+a\left(y^{\prime}(0)+s y(0)\right)+b y(0) . \tag{8.3.13}
\end{equation*}
$$

Example 8.3.3 Use the Laplace transform to solve the initial value problem

$$
\begin{equation*}
2 y^{\prime \prime}+3 y^{\prime}+y=8 e^{-2 t}, \quad y(0)=-4, y^{\prime}(0)=2 . \tag{8.3.14}
\end{equation*}
$$

Solution The characteristic polynomial is

$$
p(s)=2 s^{2}+3 s+1=(2 s+1)(s+1)
$$

and

$$
F(s)=\mathcal{L}\left(8 e^{-2 t}\right)=\frac{8}{s+2},
$$

so (8.3.13) becomes

$$
(2 s+1)(s+1) Y(s)=\frac{8}{s+2}+2(2-4 s)+3(-4)
$$

Solving for $Y(s)$ yields

$$
Y(s)=\frac{4(1-(s+2)(s+1))}{(s+1 / 2)(s+1)(s+2)}
$$

Heaviside's method yields the partial fraction expansion

$$
Y(s)=\frac{4}{3} \frac{1}{s+1 / 2}-\frac{8}{s+1}+\frac{8}{3} \frac{1}{s+2}
$$

so the solution of (8.3.14) is

$$
y=\mathcal{L}^{-1}(Y(s))=\frac{4}{3} e^{-t / 2}-8 e^{-t}+\frac{8}{3} e^{-2 t}
$$

(Figure 8.3.1).
Example 8.3.4 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+2 y^{\prime}+2 y=1, \quad y(0)=-3, y^{\prime}(0)=1 . \tag{8.3.15}
\end{equation*}
$$



Figure 8.3.1 $y=\frac{4}{3} e^{-t / 2}-8 e^{-t}+\frac{8}{3} e^{-2 t}$


Figure 8.3.2 $y=\frac{1}{2}-\frac{7}{2} e^{-t} \cos t-\frac{5}{2} e^{-t} \sin t$

Solution The characteristic polynomial is

$$
p(s)=s^{2}+2 s+2=(s+1)^{2}+1
$$

and

$$
F(s)=\mathcal{L}(1)=\frac{1}{s}
$$

so (8.3.13) becomes

$$
\left[(s+1)^{2}+1\right] Y(s)=\frac{1}{s}+1 \cdot(1-3 s)+2(-3)
$$

Solving for $Y(s)$ yields

$$
Y(s)=\frac{1-s(5+3 s)}{s\left[(s+1)^{2}+1\right]}
$$

In Example 8.2.8 we found the inverse transform of this function to be

$$
y=\frac{1}{2}-\frac{7}{2} e^{-t} \cos t-\frac{5}{2} e^{-t} \sin t
$$

(Figure 8.3.2), which is therefore the solution of (8.3.15).
REMARK: In our examples we applied Theorems 8.3.1 and 8.3.2 without verifying that the unknown function $y$ satisfies their hypotheses. This is characteristic of the formal manipulative way in which the Laplace transform is used to solve differential equations. Any doubts about the validity of the method for solving a given equation can be resolved by verifying that the resulting function $y$ is the solution of the given problem.

### 8.3 Exercises

In Exercises 1-31 use the Laplace transform to solve the initial value problem.

1. $y^{\prime \prime}+3 y^{\prime}+2 y=e^{t}, \quad y(0)=1, \quad y^{\prime}(0)=-6$
2. $y^{\prime \prime}-y^{\prime}-6 y=2, \quad y(0)=1, \quad y^{\prime}(0)=0$
3. $y^{\prime \prime}+y^{\prime}-2 y=2 e^{3 t}, \quad y(0)=-1, \quad y^{\prime}(0)=4$
4. $y^{\prime \prime}-4 y=2 e^{3 t}, \quad y(0)=1, \quad y^{\prime}(0)=-1$
5. $y^{\prime \prime}+y^{\prime}-2 y=e^{3 t}, \quad y(0)=1, \quad y^{\prime}(0)=-1$
6. $y^{\prime \prime}+3 y^{\prime}+2 y=6 e^{t}, \quad y(0)=1, \quad y^{\prime}(0)=-1$
7. $y^{\prime \prime}+y=\sin 2 t, \quad y(0)=0, \quad y^{\prime}(0)=1$
8. $y^{\prime \prime}-3 y^{\prime}+2 y=2 e^{3 t}, \quad y(0)=1, \quad y^{\prime}(0)=-1$
9. $y^{\prime \prime}-3 y^{\prime}+2 y=e^{4 t}, \quad y(0)=1, \quad y^{\prime}(0)=-2$
10. $y^{\prime \prime}-3 y^{\prime}+2 y=e^{3 t}, \quad y(0)=-1, \quad y^{\prime}(0)=-4$
11. $y^{\prime \prime}+3 y^{\prime}+2 y=2 e^{t}, \quad y(0)=0, \quad y^{\prime}(0)=-1$
12. $y^{\prime \prime}+y^{\prime}-2 y=-4, \quad y(0)=2, \quad y^{\prime}(0)=3$
13. $y^{\prime \prime}+4 y=4, \quad y(0)=0, \quad y^{\prime}(0)=1$
14. $y^{\prime \prime}-y^{\prime}-6 y=2, \quad y(0)=1, \quad y^{\prime}(0)=0$
15. $y^{\prime \prime}+3 y^{\prime}+2 y=e^{t}, \quad y(0)=0, \quad y^{\prime}(0)=1$
16. $\quad y^{\prime \prime}-y=1, \quad y(0)=1, \quad y^{\prime}(0)=0$
17. $y^{\prime \prime}+4 y=3 \sin t, \quad y(0)=1, \quad y^{\prime}(0)=-1$
18. $y^{\prime \prime}+y^{\prime}=2 e^{3 t}, \quad y(0)=-1, \quad y^{\prime}(0)=4$
19. $y^{\prime \prime}+y=1, \quad y(0)=2, \quad y^{\prime}(0)=0$
20. $y^{\prime \prime}+y=t, \quad y(0)=0, \quad y^{\prime}(0)=2$
21. $y^{\prime \prime}+y=t-3 \sin 2 t, \quad y(0)=1, \quad y^{\prime}(0)=-3$
22. $y^{\prime \prime}+5 y^{\prime}+6 y=2 e^{-t}, \quad y(0)=1, \quad y^{\prime}(0)=3$
23. $y^{\prime \prime}+2 y^{\prime}+y=6 \sin t-4 \cos t, \quad y(0)=-1, y^{\prime}(0)=1$
24. $y^{\prime \prime}-2 y^{\prime}-3 y=10 \cos t, \quad y(0)=2, \quad y^{\prime}(0)=7$
25. $y^{\prime \prime}+y=4 \sin t+6 \cos t, \quad y(0)=-6, y^{\prime}(0)=2$
26. $y^{\prime \prime}+4 y=8 \sin 2 t+9 \cos t, \quad y(0)=1, \quad y^{\prime}(0)=0$
27. $y^{\prime \prime}-5 y^{\prime}+6 y=10 e^{t} \cos t, \quad y(0)=2, \quad y^{\prime}(0)=1$
28. $y^{\prime \prime}+2 y^{\prime}+2 y=2 t, \quad y(0)=2, \quad y^{\prime}(0)=-7$
29. $y^{\prime \prime}-2 y^{\prime}+2 y=5 \sin t+10 \cos t, \quad y(0)=1, y^{\prime}(0)=2$
30. $y^{\prime \prime}+4 y^{\prime}+13 y=10 e^{-t}-36 e^{t}, \quad y(0)=0, y^{\prime}(0)=-16$
31. $y^{\prime \prime}+4 y^{\prime}+5 y=e^{-t}(\cos t+3 \sin t), \quad y(0)=0, \quad y^{\prime}(0)=4$
32. $2 y^{\prime \prime}-3 y^{\prime}-2 y=4 e^{t}, \quad y(0)=1, y^{\prime}(0)=-2$
33. $6 y^{\prime \prime}-y^{\prime}-y=3 e^{2 t}, \quad y(0)=0, y^{\prime}(0)=0$
34. $2 y^{\prime \prime}+2 y^{\prime}+y=2 t, \quad y(0)=1, y^{\prime}(0)=-1$
35. $4 y^{\prime \prime}-4 y^{\prime}+5 y=4 \sin t-4 \cos t, \quad y(0)=0, y^{\prime}(0)=11 / 17$
36. $4 y^{\prime \prime}+4 y^{\prime}+y=3 \sin t+\cos t, \quad y(0)=2, y^{\prime}(0)=-1$
37. $9 y^{\prime \prime}+6 y^{\prime}+y=3 e^{3 t}, \quad y(0)=0, y^{\prime}(0)=-3$
38. Suppose $a, b$, and $c$ are constants and $a \neq 0$. Let

$$
y_{1}=\mathcal{L}^{-1}\left(\frac{a s+b}{a s^{2}+b s+c}\right) \quad \text { and } \quad y_{2}=\mathcal{L}^{-1}\left(\frac{a}{a s^{2}+b s+c}\right) .
$$

Show that

$$
y_{1}(0)=1, \quad y_{1}^{\prime}(0)=0 \quad \text { and } \quad y_{2}(0)=0, \quad y_{2}^{\prime}(0)=1 .
$$

Hint: Use the Laplace transform to solve the initial value problems

$$
\begin{aligned}
& a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y(0)=1, \quad y^{\prime}(0)=0 \\
& a y^{\prime \prime}+b y^{\prime}+c y=0, \quad y(0)=0, \quad y^{\prime}(0)=1
\end{aligned}
$$

### 8.4 THE UNIT STEP FUNCTION

In the next section we'll consider initial value problems

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1},
$$

where $a, b$, and $c$ are constants and $f$ is piecewise continuous. In this section we'll develop procedures for using the table of Laplace transforms to find Laplace transforms of piecewise continuous functions, and to find the piecewise continuous inverses of Laplace transforms.

Example 8.4.1 Use the table of Laplace transforms to find the Laplace transform of

$$
f(t)=\left\{\begin{array}{cl}
2 t+1, & 0 \leq t<2  \tag{8.4.1}\\
3 t, & t \geq 2
\end{array}\right.
$$

(Figure 8.4.1).

Solution Since the formula for $f$ changes at $t=2$, we write

$$
\begin{align*}
\mathcal{L}(f) & =\int_{0}^{\infty} e^{-s t} f(t) d t \\
& =\int_{0}^{2} e^{-s t}(2 t+1) d t+\int_{2}^{\infty} e^{-s t}(3 t) d t \tag{8.4.2}
\end{align*}
$$

To relate the first term to a Laplace transform, we add and subtract

$$
\int_{2}^{\infty} e^{-s t}(2 t+1) d t
$$

in (8.4.2) to obtain

$$
\begin{align*}
\mathcal{L}(f) & =\int_{0}^{\infty} e^{-s t}(2 t+1) d t+\int_{2}^{\infty} e^{-s t}(3 t-2 t-1) d t \\
& =\int_{0}^{\infty} e^{-s t}(2 t+1) d t+\int_{2}^{\infty} e^{-s t}(t-1) d t  \tag{8.4.3}\\
& =\mathcal{L}(2 t+1)+\int_{2}^{\infty} e^{-s t}(t-1) d t .
\end{align*}
$$

To relate the last integral to a Laplace transform, we make the change of variable $x=t-2$ and rewrite the integral as

$$
\begin{aligned}
\int_{2}^{\infty} e^{-s t}(t-1) d t & =\int_{0}^{\infty} e^{-s(x+2)}(x+1) d x \\
& =e^{-2 s} \int_{0}^{\infty} e^{-s x}(x+1) d x
\end{aligned}
$$

Since the symbol used for the variable of integration has no effect on the value of a definite integral, we can now replace $x$ by the more standard $t$ and write

$$
\int_{2}^{\infty} e^{-s t}(t-1) d t=e^{-2 s} \int_{0}^{\infty} e^{-s t}(t+1) d t=e^{-2 s} \mathcal{L}(t+1)
$$

This and (8.4.3) imply that

$$
\mathcal{L}(f)=\mathcal{L}(2 t+1)+e^{-2 s} \mathcal{L}(t+1)
$$

Now we can use the table of Laplace transforms to find that

$$
\mathcal{L}(f)=\frac{2}{s^{2}}+\frac{1}{s}+e^{-2 s}\left(\frac{1}{s^{2}}+\frac{1}{s}\right) .
$$



Figure 8.4.1 The piecewise continuous function (8.4.1)


Figure 8.4.2 $y=u(t-\tau)$

## Laplace Transforms of Piecewise Continuous Functions

We'll now develop the method of Example 8.4.1 into a systematic way to find the Laplace transform of a piecewise continuous function. It is convenient to introduce the unit step function, defined as

$$
u(t)= \begin{cases}0, & t<0  \tag{8.4.4}\\ 1, & t \geq 0\end{cases}
$$

Thus, $u(t)$ "steps" from the constant value 0 to the constant value 1 at $t=0$. If we replace $t$ by $t-\tau$ in (8.4.4), then

$$
u(t-\tau)= \begin{cases}0, & t<\tau \\ 1, & t \geq \tau\end{cases}
$$

that is, the step now occurs at $t=\tau$ (Figure 8.4.2).
The step function enables us to represent piecewise continuous functions conveniently. For example, consider the function

$$
f(t)= \begin{cases}f_{0}(t), & 0 \leq t<t_{1}  \tag{8.4.5}\\ f_{1}(t), & t \geq t_{1}\end{cases}
$$

where we assume that $f_{0}$ and $f_{1}$ are defined on $[0, \infty)$, even though they equal $f$ only on the indicated intervals. This assumption enables us to rewrite (8.4.5) as

$$
\begin{equation*}
f(t)=f_{0}(t)+u\left(t-t_{1}\right)\left(f_{1}(t)-f_{0}(t)\right) \tag{8.4.6}
\end{equation*}
$$

To verify this, note that if $t<t_{1}$ then $u\left(t-t_{1}\right)=0$ and (8.4.6) becomes

$$
f(t)=f_{0}(t)+(0)\left(f_{1}(t)-f_{0}(t)\right)=f_{0}(t)
$$

If $t \geq t_{1}$ then $u\left(t-t_{1}\right)=1$ and (8.4.6) becomes

$$
f(t)=f_{0}(t)+(1)\left(f_{1}(t)-f_{0}(t)\right)=f_{1}(t)
$$

We need the next theorem to show how (8.4.6) can be used to find $\mathcal{L}(f)$.
Theorem 8.4.1 Let $g$ be defined on $[0, \infty)$. Suppose $\tau \geq 0$ and $\mathcal{L}(g(t+\tau))$ exists for $s>s_{0}$. Then $\mathcal{L}(u(t-\tau) g(t))$ exists for $s>s_{0}$, and

$$
\mathcal{L}(u(t-\tau) g(t))=e^{-s \tau} \mathcal{L}(g(t+\tau)) .
$$

Proof By definition,

$$
\mathcal{L}(u(t-\tau) g(t))=\int_{0}^{\infty} e^{-s t} u(t-\tau) g(t) d t
$$

From this and the definition of $u(t-\tau)$,

$$
\mathcal{L}(u(t-\tau) g(t))=\int_{0}^{\tau} e^{-s t}(0) d t+\int_{\tau}^{\infty} e^{-s t} g(t) d t
$$

The first integral on the right equals zero. Introducing the new variable of integration $x=t-\tau$ in the second integral yields

$$
\mathcal{L}(u(t-\tau) g(t))=\int_{0}^{\infty} e^{-s(x+\tau)} g(x+\tau) d x=e^{-s \tau} \int_{0}^{\infty} e^{-s x} g(x+\tau) d x
$$

Changing the name of the variable of integration in the last integral from $x$ to $t$ yields

$$
\mathcal{L}(u(t-\tau) g(t))=e^{-s \tau} \int_{0}^{\infty} e^{-s t} g(t+\tau) d t=e^{-s \tau} \mathcal{L}(g(t+\tau))
$$

Example 8.4.2 Find

$$
\mathcal{L}\left(u(t-1)\left(t^{2}+1\right)\right) .
$$

Solution Here $\tau=1$ and $g(t)=t^{2}+1$, so

$$
g(t+1)=(t+1)^{2}+1=t^{2}+2 t+2 .
$$

Since

$$
\mathcal{L}(g(t+1))=\frac{2}{s^{3}}+\frac{2}{s^{2}}+\frac{2}{s}
$$

Theorem 8.4.1 implies that

$$
\mathcal{L}\left(u(t-1)\left(t^{2}+1\right)\right)=e^{-s}\left(\frac{2}{s^{3}}+\frac{2}{s^{2}}+\frac{2}{s}\right) .
$$

Example 8.4.3 Use Theorem 8.4.1 to find the Laplace transform of the function

$$
f(t)=\left\{\begin{array}{cl}
2 t+1, & 0 \leq t<2 \\
3 t, & t \geq 2
\end{array}\right.
$$

from Example 8.4.1.

Solution We first write $f$ in the form (8.4.6) as

$$
f(t)=2 t+1+u(t-2)(t-1)
$$

Therefore

$$
\begin{aligned}
\mathcal{L}(f) & =\mathcal{L}(2 t+1)+\mathcal{L}(u(t-2)(t-1)) \\
& =\mathcal{L}(2 t+1)+e^{-2 s} \mathcal{L}(t+1) \quad \text { (from Theorem 8.4.1) } \\
& =\frac{2}{s^{2}}+\frac{1}{s}+e^{-2 s}\left(\frac{1}{s^{2}}+\frac{1}{s}\right)
\end{aligned}
$$

which is the result obtained in Example 8.4.1.
Formula (8.4.6) can be extended to more general piecewise continuous functions. For example, we can write

$$
f(t)= \begin{cases}f_{0}(t), & 0 \leq t<t_{1} \\ f_{1}(t), & t_{1} \leq t<t_{2} \\ f_{2}(t), & t \geq t_{2}\end{cases}
$$

as

$$
f(t)=f_{0}(t)+u\left(t-t_{1}\right)\left(f_{1}(t)-f_{0}(t)\right)+u\left(t-t_{2}\right)\left(f_{2}(t)-f_{1}(t)\right)
$$

if $f_{0}, f_{1}$, and $f_{2}$ are all defined on $[0, \infty)$.
Example 8.4.4 Find the Laplace transform of

$$
f(t)=\left\{\begin{array}{cl}
1, & 0 \leq t<2  \tag{8.4.7}\\
-2 t+1, & 2 \leq t<3 \\
3 t, & 3 \leq t<5 \\
t-1, & t \geq 5
\end{array}\right.
$$

(Figure 8.4.3).

Solution In terms of step functions,

$$
\begin{aligned}
f(t)= & 1+u(t-2)(-2 t+1-1)+u(t-3)(3 t+2 t-1) \\
& +u(t-5)(t-1-3 t)
\end{aligned}
$$

or

$$
f(t)=1-2 u(t-2) t+u(t-3)(5 t-1)-u(t-5)(2 t+1)
$$



Figure 8.4.3 The piecewise contnuous function (8.4.7)

Now Theorem 8.4.1 implies that

$$
\begin{aligned}
\mathcal{L}(f) & =\mathcal{L}(1)-2 e^{-2 s} \mathcal{L}(t+2)+e^{-3 s} \mathcal{L}(5(t+3)-1)-e^{-5 s} \mathcal{L}(2(t+5)+1) \\
& =\mathcal{L}(1)-2 e^{-2 s} \mathcal{L}(t+2)+e^{-3 s} \mathcal{L}(5 t+14)-e^{-5 s} \mathcal{L}(2 t+11) \\
& =\frac{1}{s}-2 e^{-2 s}\left(\frac{1}{s^{2}}+\frac{2}{s}\right)+e^{-3 s}\left(\frac{5}{s^{2}}+\frac{14}{s}\right)-e^{-5 s}\left(\frac{2}{s^{2}}+\frac{11}{s}\right) .
\end{aligned}
$$

The trigonometric identities

$$
\begin{align*}
& \sin (A+B)=\sin A \cos B+\cos A \sin B  \tag{8.4.8}\\
& \cos (A+B)=\cos A \cos B-\sin A \sin B \tag{8.4.9}
\end{align*}
$$

are useful in problems that involve shifting the arguments of trigonometric functions. We'll use these identities in the next example.

Example 8.4.5 Find the Laplace transform of

$$
f(t)=\left\{\begin{array}{cl}
\sin t, & 0 \leq t<\frac{\pi}{2}  \tag{8.4.10}\\
\cos t-3 \sin t, & \frac{\pi}{2} \leq t<\pi \\
3 \cos t, & t \geq \pi
\end{array}\right.
$$

(Figure 8.4.4).

Solution In terms of step functions,

$$
f(t)=\sin t+u(t-\pi / 2)(\cos t-4 \sin t)+u(t-\pi)(2 \cos t+3 \sin t)
$$

Now Theorem 8.4.1 implies that

$$
\begin{gather*}
\mathcal{L}(f)=\mathcal{L}(\sin t)+e^{-\frac{\pi}{2} s} \mathcal{L}\left(\cos \left(t+\frac{\pi}{2}\right)-4 \sin \left(t+\frac{\pi}{2}\right)\right)  \tag{8.4.11}\\
+e^{-\pi s} \mathcal{L}(2 \cos (t+\pi)+3 \sin (t+\pi))
\end{gather*}
$$

Since

$$
\cos \left(t+\frac{\pi}{2}\right)-4 \sin \left(t+\frac{\pi}{2}\right)=-\sin t-4 \cos t
$$



Figure 8.4.4 The piecewise continuous function (8.4.10)
and

$$
2 \cos (t+\pi)+3 \sin (t+\pi)=-2 \cos t-3 \sin t
$$

we see from (8.4.11) that

$$
\begin{aligned}
\mathcal{L}(f) & =\mathcal{L}(\sin t)-e^{-\pi s / 2} \mathcal{L}(\sin t+4 \cos t)-e^{-\pi s} \mathcal{L}(2 \cos t+3 \sin t) \\
& =\frac{1}{s^{2}+1}-e^{-\frac{\pi}{2} s}\left(\frac{1+4 s}{s^{2}+1}\right)-e^{-\pi s}\left(\frac{3+2 s}{s^{2}+1}\right)
\end{aligned}
$$

The Second Shifting Theorem
Replacing $g(t)$ by $g(t-\tau)$ in Theorem 8.4.1 yields the next theorem.
Theorem 8.4.2 [Second Shifting Theorem] If $\tau \geq 0$ and $\mathcal{L}(g)$ exists for $s>s_{0}$ then $\mathcal{L}(u(t-\tau) g(t-\tau))$ exists for $s>s_{0}$ and

$$
\mathcal{L}(u(t-\tau) g(t-\tau))=e^{-s \tau} \mathcal{L}(g(t)),
$$

or, equivalently,

$$
\begin{equation*}
\text { if } g(t) \leftrightarrow G(s) \text {, then } u(t-\tau) g(t-\tau) \leftrightarrow e^{-s \tau} G(s) \text {. } \tag{8.4.12}
\end{equation*}
$$

REmARK: Recall that the First Shifting Theorem (Theorem 8.1.3 states that multiplying a function by $e^{a t}$ corresponds to shifting the argument of its transform by $a$ units. Theorem 8.4.2 states that multiplying a Laplace transform by the exponential $e^{-\tau s}$ corresponds to shifting the argument of the inverse transform by $\tau$ units.
Example 8.4.6 Use (8.4.12) to find

$$
\mathcal{L}^{-1}\left(\frac{e^{-2 s}}{s^{2}}\right) .
$$

Solution To apply (8.4.12) we let $\tau=2$ and $G(s)=1 / s^{2}$. Then $g(t)=t$ and (8.4.12) implies that

$$
\mathcal{L}^{-1}\left(\frac{e^{-2 s}}{s^{2}}\right)=u(t-2)(t-2) .
$$

Example 8.4.7 Find the inverse Laplace transform $h$ of

$$
H(s)=\frac{1}{s^{2}}-e^{-s}\left(\frac{1}{s^{2}}+\frac{2}{s}\right)+e^{-4 s}\left(\frac{4}{s^{3}}+\frac{1}{s}\right),
$$

and find distinct formulas for $h$ on appropriate intervals.

Solution Let

$$
G_{0}(s)=\frac{1}{s^{2}}, \quad G_{1}(s)=\frac{1}{s^{2}}+\frac{2}{s}, \quad G_{2}(s)=\frac{4}{s^{3}}+\frac{1}{s}
$$

Then

$$
g_{0}(t)=t, g_{1}(t)=t+2, g_{2}(t)=2 t^{2}+1
$$

Hence, (8.4.12) and the linearity of $\mathcal{L}^{-1}$ imply that

$$
\begin{aligned}
h(t) & =\mathcal{L}^{-1}\left(G_{0}(s)\right)-\mathcal{L}^{-1}\left(e^{-s} G_{1}(s)\right)+\mathcal{L}^{-1}\left(e^{-4 s} G_{2}(s)\right) \\
& =t-u(t-1)[(t-1)+2]+u(t-4)\left[2(t-4)^{2}+1\right] \\
& =t-u(t-1)(t+1)+u(t-4)\left(2 t^{2}-16 t+33\right),
\end{aligned}
$$

which can also be written as

$$
h(t)=\left\{\begin{array}{cl}
t, & 0 \leq t<1 \\
-1, & 1 \leq t<4 \\
2 t^{2}-16 t+32, & t \geq 4
\end{array}\right.
$$

Example 8.4.8 Find the inverse transform of

$$
H(s)=\frac{2 s}{s^{2}+4}-e^{-\frac{\pi}{2} s} \frac{3 s+1}{s^{2}+9}+e^{-\pi s} \frac{s+1}{s^{2}+6 s+10} .
$$

## Solution Let

$$
G_{0}(s)=\frac{2 s}{s^{2}+4}, \quad G_{1}(s)=-\frac{(3 s+1)}{s^{2}+9}
$$

and

$$
G_{2}(s)=\frac{s+1}{s^{2}+6 s+10}=\frac{(s+3)-2}{(s+3)^{2}+1}
$$

Then

$$
g_{0}(t)=2 \cos 2 t, \quad g_{1}(t)=-3 \cos 3 t-\frac{1}{3} \sin 3 t
$$

and

$$
g_{2}(t)=e^{-3 t}(\cos t-2 \sin t)
$$

Therefore (8.4.12) and the linearity of $\mathcal{L}^{-1}$ imply that

$$
\begin{aligned}
h(t)= & 2 \cos 2 t-u(t-\pi / 2)\left[3 \cos 3(t-\pi / 2)+\frac{1}{3} \sin 3\left(t-\frac{\pi}{2}\right)\right] \\
& +u(t-\pi) e^{-3(t-\pi)}[\cos (t-\pi)-2 \sin (t-\pi)]
\end{aligned}
$$

Using the trigonometric identities (8.4.8) and (8.4.9), we can rewrite this as

$$
\begin{align*}
h(t)= & 2 \cos 2 t+u(t-\pi / 2)\left(3 \sin 3 t-\frac{1}{3} \cos 3 t\right)  \tag{8.4.13}\\
& -u(t-\pi) e^{-3(t-\pi)}(\cos t-2 \sin t)
\end{align*}
$$

(Figure 8.4.5).


Figure 8.4.5 The piecewise continouous function (8.4.13)

### 8.4 Exercises

In Exercises 1-6 find the Laplace transform by the method of Example 8.4.1. Then express the given function $f$ in terms of unit step functions as in Eqn. (8.4.6), and use Theorem 8.4.1 to find $\mathcal{L}(f)$. Where indicated by $\mathrm{C} / \mathrm{G}$, graph $f$.

1. $f(t)= \begin{cases}1, & 0 \leq t<4, \\ t, & t \geq 4 .\end{cases}$
2. $f(t)= \begin{cases}t, & 0 \leq t<1, \\ 1, & t \geq 1 .\end{cases}$
3. $\mathrm{C} / \mathrm{G} f(t)=\left\{\begin{array}{cl}2 t-1, & 0 \leq t<2, \\ t, & t \geq 2 .\end{array}\right.$ 4. $\quad \mathrm{C} / \mathrm{G} \quad f(t)=\left\{\begin{array}{cl}1, & 0 \leq t<1, \\ t+2, & t \geq 1 .\end{array}\right.$
4. $f(t)=\left\{\begin{array}{cl}t-1, & 0 \leq t<2, \\ 4, & t \geq 2 .\end{array}\right.$
5. $f(t)=\left\{\begin{array}{cl}t^{2}, & 0 \leq t<1, \\ 0, & t \geq 1 .\end{array}\right.$

In Exercises 7-18 express the given function $f$ in terms of unit step functions and use Theorem 8.4.1 to find $\mathcal{L}(f)$. Where indicated by $C / G$, graph $f$.
7. $f(t)=\left\{\begin{array}{cl}0, & \\ 0 \leq t<2, \\ t^{2}+3 t, & \\ t \geq 2 .\end{array}\right.$
8. $f(t)=\left\{\begin{array}{cl}t^{2}+2, & 0 \leq t<1, \\ t, & t \geq 1 .\end{array}\right.$
9. $f(t)=\left\{\begin{array}{cl}t e^{t}, & 0 \leq t<1, \\ e^{t}, & t \geq 1 .\end{array}\right.$
10. $f(t)= \begin{cases}e^{-t}, & 0 \leq t<1, \\ e^{-2 t}, & t \geq 1 .\end{cases}$
11. $f(t)=\left\{\begin{array}{cl}-t, & 0 \leq t<2, \\ t-4, & 2 \leq t<3, \\ 1, & t \geq 3 .\end{array}\right.$
12. $f(t)= \begin{cases}0, & 0 \leq t<1, \\ t, & 1 \leq t<2, \\ 0, & t \geq 2 .\end{cases}$
13. $f(t)=\left\{\begin{array}{cl}t, & 0 \leq t<1, \\ t^{2}, & 1 \leq t<2, \\ 0, & t \geq 2 .\end{array} \quad\right.$ 14. $\quad f(t)=\left\{\begin{array}{cl}t, & 0 \leq t<1, \\ 2-t, & 1 \leq t<2, \\ 6, & t>2 .\end{array}\right.$
15. $\mathrm{C} / \mathrm{G} ~ f(t)=\left\{\begin{aligned} \sin t, & 0 \leq t<\frac{\pi}{2}, \\ 2 \sin t, & \frac{\pi}{2} \leq t<\pi, \\ \cos t, & t \geq \pi .\end{aligned}\right.$
16. $\mathrm{C} / \mathrm{G} f(t)=\left\{\begin{array}{cl}2, & 0 \leq t<1, \\ -2 t+2, & 1 \leq t<3, \\ 3 t, & t \geq 3 .\end{array}\right.$
17. $\mathrm{C} / \mathrm{G} f(t)=\left\{\begin{array}{cl}3, & 0 \leq t<2, \\ 3 t+2, & 2 \leq t<4, \\ 4 t, & t \geq 4 .\end{array}\right.$
18. $\mathrm{C} / \mathrm{G} f(t)= \begin{cases}(t+1)^{2}, & 0 \leq t<1, \\ (t+2)^{2}, & t \geq 1 .\end{cases}$

In Exercises 19-28 use Theorem 8.4.2 to express the inverse transforms in terms of step functions, and then find distinct formulas the for inverse transforms on the appropriate intervals, as in Example 8.4.7. Where indicated by $C / G$, graph the inverse transform.
19. $H(s)=\frac{e^{-2 s}}{s-2}$
20. $H(s)=\frac{e^{-s}}{s(s+1)}$
21. $\mathrm{C} / \mathrm{G} H(s)=\frac{e^{-s}}{s^{3}}+\frac{e^{-2 s}}{s^{2}}$
22. $\mathrm{C} / \mathrm{G} H(s)=\left(\frac{2}{s}+\frac{1}{s^{2}}\right)+e^{-s}\left(\frac{3}{s}-\frac{1}{s^{2}}\right)+e^{-3 s}\left(\frac{1}{s}+\frac{1}{s^{2}}\right)$
23. $H(s)=\left(\frac{5}{s}-\frac{1}{s^{2}}\right)+e^{-3 s}\left(\frac{6}{s}+\frac{7}{s^{2}}\right)+\frac{3 e^{-6 s}}{s^{3}}$
24. $H(s)=\frac{e^{-\pi s}(1-2 s)}{s^{2}+4 s+5}$
25. $\mathrm{C} / \mathrm{G} H(s)=\left(\frac{1}{s}-\frac{s}{s^{2}+1}\right)+e^{-\frac{\pi}{2} s}\left(\frac{3 s-1}{s^{2}+1}\right)$
26. $H(s)=e^{-2 s}\left[\frac{3(s-3)}{(s+1)(s-2)}-\frac{s+1}{(s-1)(s-2)}\right]$
27. $H(s)=\frac{1}{s}+\frac{1}{s^{2}}+e^{-s}\left(\frac{3}{s}+\frac{2}{s^{2}}\right)+e^{-3 s}\left(\frac{4}{s}+\frac{3}{s^{2}}\right)$
28. $H(s)=\frac{1}{s}-\frac{2}{s^{3}}+e^{-2 s}\left(\frac{3}{s}-\frac{1}{s^{3}}\right)+\frac{e^{-4 s}}{s^{2}}$
29. Find $\mathcal{L}(u(t-\tau))$.
30. Let $\left\{t_{m}\right\}_{m=0}^{\infty}$ be a sequence of points such that $t_{0}=0, t_{m+1}>t_{m}$, and $\lim _{m \rightarrow \infty} t_{m}=\infty$. For each nonnegative integer $m$, let $f_{m}$ be continuous on $\left[t_{m}, \infty\right)$, and let $f$ be defined on $[0, \infty)$ by

$$
f(t)=f_{m}(t), t_{m} \leq t<t_{m+1} \quad(m=0,1, \ldots)
$$

Show that $f$ is piecewise continuous on $[0, \infty)$ and that it has the step function representation

$$
f(t)=f_{0}(t)+\sum_{m=1}^{\infty} u\left(t-t_{m}\right)\left(f_{m}(t)-f_{m-1}(t)\right), 0 \leq t<\infty .
$$

How do we know that the series on the right converges for all $t$ in $[0, \infty)$ ?
31. In addition to the assumptions of Exercise 30, assume that

$$
\begin{equation*}
\left|f_{m}(t)\right| \leq M e^{s_{0} t}, t \geq t_{m}, m=0,1, \ldots \tag{A}
\end{equation*}
$$

and that the series

$$
\begin{equation*}
\sum_{m=0}^{\infty} e^{-\rho t_{m}} \tag{B}
\end{equation*}
$$

converges for some $\rho>0$. Using the steps listed below, show that $\mathcal{L}(f)$ is defined for $s>s_{0}$ and

$$
\begin{equation*}
\mathcal{L}(f)=\mathcal{L}\left(f_{0}\right)+\sum_{m=1}^{\infty} e^{-s t_{m}} \mathcal{L}\left(g_{m}\right) \tag{C}
\end{equation*}
$$

for $s>s_{0}+\rho$, where

$$
g_{m}(t)=f_{m}\left(t+t_{m}\right)-f_{m-1}\left(t+t_{m}\right)
$$

(a) Use (A) and Theorem 8.1.6 to show that

$$
\begin{equation*}
\mathcal{L}(f)=\sum_{m=0}^{\infty} \int_{t_{m}}^{t_{m+1}} e^{-s t} f_{m}(t) d t \tag{D}
\end{equation*}
$$

is defined for $s>s_{0}$.
(b) Show that (D) can be rewritten as

$$
\begin{equation*}
\mathcal{L}(f)=\sum_{m=0}^{\infty}\left(\int_{t_{m}}^{\infty} e^{-s t} f_{m}(t) d t-\int_{t_{m+1}}^{\infty} e^{-s t} f_{m}(t) d t\right) \tag{E}
\end{equation*}
$$

(c) Use (A), the assumed convergence of (B), and the comparison test to show that the series

$$
\sum_{m=0}^{\infty} \int_{t_{m}}^{\infty} e^{-s t} f_{m}(t) d t \quad \text { and } \quad \sum_{m=0}^{\infty} \int_{t_{m+1}}^{\infty} e^{-s t} f_{m}(t) d t
$$

both converge (absolutely) if $s>s_{0}+\rho$.
(d) Show that (E) can be rewritten as

$$
\mathcal{L}(f)=\mathcal{L}\left(f_{0}\right)+\sum_{m=1}^{\infty} \int_{t_{m}}^{\infty} e^{-s t}\left(f_{m}(t)-f_{m-1}(t)\right) d t
$$

if $s>s_{0}+\rho$.
(e) Complete the proof of (C).
32. Suppose $\left\{t_{m}\right\}_{m=0}^{\infty}$ and $\left\{f_{m}\right\}_{m=0}^{\infty}$ satisfy the assumptions of Exercises 30 and 31, and there's a positive constant $K$ such that $t_{m} \geq K m$ for $m$ sufficiently large. Show that the series (B) of Exercise 31 converges for any $\rho>0$, and conclude from this that (C) of Exercise 31 holds for $s>s_{0}$.

In Exercises 33-36 find the step function representation of $f$ and use the result of Exercise 32 to find $\mathcal{L}(f)$. Hint: You will need formulas related to the formula for the sum of a geometric series.
33. $f(t)=m+1, m \leq t<m+1(m=0,1,2, \ldots)$
34. $f(t)=(-1)^{m}, m \leq t<m+1(m=0,1,2, \ldots)$
35. $f(t)=(m+1)^{2}, m \leq t<m+1(m=0,1,2, \ldots)$
36. $f(t)=(-1)^{m} m, m \leq t<m+1(m=0,1,2, \ldots)$

### 8.5 CONSTANT COEEFFICIENT EQUATIONS WITH PIECEWISE CONTINUOUS FORCING FUNCTIONS

We'll now consider initial value problems of the form

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} \tag{8.5.1}
\end{equation*}
$$

where $a, b$, and $c$ are constants $(a \neq 0)$ and $f$ is piecewise continuous on $[0, \infty)$. Problems of this kind occur in situations where the input to a physical system undergoes instantaneous changes, as when a switch is turned on or off or the forces acting on the system change abruptly.

It can be shown (Exercises 23 and 24) that the differential equation in (8.5.1) has no solutions on an open interval that contains a jump discontinuity of $f$. Therefore we must define what we mean by a solution of (8.5.1) on $[0, \infty)$ in the case where $f$ has jump discontinuities. The next theorem motivates our definition. We omit the proof.

Theorem 8.5.1 Suppose $a, b$, and c are constants $(a \neq 0)$, and $f$ is piecewise continuous on $[0, \infty)$. with jump discontinuities at $t_{1}, \ldots, t_{n}$, where

$$
0<t_{1}<\cdots<t_{n} .
$$

Let $k_{0}$ and $k_{1}$ be arbitrary real numbers. Then there is a unique function $y$ defined on $[0, \infty)$ with these properties:
(a) $y(0)=k_{0}$ and $y^{\prime}(0)=k_{1}$.
(b) $y$ and $y^{\prime}$ are continuous on $[0, \infty)$.
(c) $y^{\prime \prime}$ is defined on every open subinterval of $[0, \infty)$ that does not contain any of the points $t_{1}, \ldots, t_{n}$, and

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t)
$$

on every such subinterval.
(d) $y^{\prime \prime}$ has limits from the right and left at $t_{1}, \ldots, t_{n}$.

We define the function $y$ of Theorem 8.5 .1 to be the solution of the initial value problem (8.5.1).
We begin by considering initial value problems of the form

$$
a y^{\prime \prime}+b y^{\prime}+c y=\left\{\begin{array}{ll}
f_{0}(t), & 0 \leq t<t_{1},  \tag{8.5.2}\\
f_{1}(t), & t \geq t_{1},
\end{array} \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}\right.
$$

where the forcing function has a single jump discontinuity at $t_{1}$.
We can solve (8.5.2) by the these steps:
Step 1. Find the solution $y_{0}$ of the initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=f_{0}(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

Step 2. Compute $c_{0}=y_{0}\left(t_{1}\right)$ and $c_{1}=y_{0}^{\prime}\left(t_{1}\right)$.
Step 3. Find the solution $y_{1}$ of the initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=f_{1}(t), \quad y\left(t_{1}\right)=c_{0}, \quad y^{\prime}\left(t_{1}\right)=c_{1} .
$$

Step 4. Obtain the solution $y$ of (8.5.2) as

$$
y= \begin{cases}y_{0}(t), & 0 \leq t<t_{1} \\ y_{1}(t), & t \geq t_{1} .\end{cases}
$$

It is shown in Exercise 23 that $y^{\prime}$ exists and is continuous at $t_{1}$. The next example illustrates this procedure.


Figure 8.5.1 Graph of (8.5.4)

Example 8.5.1 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+y=f(t), \quad y(0)=2, y^{\prime}(0)=-1, \tag{8.5.3}
\end{equation*}
$$

where

$$
f(t)=\left\{\begin{aligned}
1, & 0 \leq t<\frac{\pi}{2} \\
-1, & t \geq \frac{\pi}{2}
\end{aligned}\right.
$$

Solution The initial value problem in Step 1 is

$$
y^{\prime \prime}+y=1, \quad y(0)=2, \quad y^{\prime}(0)=-1 .
$$

We leave it to you to verify that its solution is

$$
y_{0}=1+\cos t-\sin t .
$$

Doing Step 2 yields $y_{0}(\pi / 2)=0$ and $y_{0}^{\prime}(\pi / 2)=-1$, so the second initial value problem is

$$
y^{\prime \prime}+y=-1, \quad y\left(\frac{\pi}{2}\right)=0, y^{\prime}\left(\frac{\pi}{2}\right)=-1 .
$$

We leave it to you to verify that the solution of this problem is

$$
y_{1}=-1+\cos t+\sin t .
$$

Hence, the solution of (8.5.3) is

$$
y=\left\{\begin{align*}
1+\cos t-\sin t, & 0 \leq t<\frac{\pi}{2}  \tag{8.5.4}\\
-1+\cos t+\sin t, & t \geq \frac{\pi}{2}
\end{align*}\right.
$$

(Figure:8.5.1).
If $f_{0}$ and $f_{1}$ are defined on $[0, \infty)$, we can rewrite (8.5.2) as

$$
a y^{\prime \prime}+b y^{\prime}+c y=f_{0}(t)+u\left(t-t_{1}\right)\left(f_{1}(t)-f_{0}(t)\right), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

and apply the method of Laplace transforms. We'll now solve the problem considered in Example 8.5.1 by this method.

Example 8.5.2 Use the Laplace transform to solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+y=f(t), \quad y(0)=2, \quad y^{\prime}(0)=-1, \tag{8.5.5}
\end{equation*}
$$

where

$$
f(t)=\left\{\begin{aligned}
1, & 0 \leq t<\frac{\pi}{2} \\
-1, & t \geq \frac{\pi}{2}
\end{aligned}\right.
$$

## Solution Here

$$
f(t)=1-2 u\left(t-\frac{\pi}{2}\right),
$$

so Theorem 8.4.1 (with $g(t)=1$ ) implies that

$$
\mathcal{L}(f)=\frac{1-2 e^{-\pi s / 2}}{s}
$$

Therefore, transforming (8.5.5) yields

$$
\left(s^{2}+1\right) Y(s)=\frac{1-2 e^{-\pi s / 2}}{s}-1+2 s,
$$

so

$$
\begin{equation*}
Y(s)=\left(1-2 e^{-\pi s / 2}\right) G(s)+\frac{2 s-1}{s^{2}+1} \tag{8.5.6}
\end{equation*}
$$

with

$$
G(s)=\frac{1}{s\left(s^{2}+1\right)}
$$

The form for the partial fraction expansion of $G$ is

$$
\begin{equation*}
\frac{1}{s\left(s^{2}+1\right)}=\frac{A}{s}+\frac{B s+C}{s^{2}+1} . \tag{8.5.7}
\end{equation*}
$$

Multiplying through by $s\left(s^{2}+1\right)$ yields

$$
A\left(s^{2}+1\right)+(B s+C) s=1,
$$

or

$$
(A+B) s^{2}+C s+A=1
$$

Equating coefficients of like powers of $s$ on the two sides of this equation shows that $A=1, B=-A=$ -1 and $C=0$. Hence, from (8.5.7),

$$
G(s)=\frac{1}{s}-\frac{s}{s^{2}+1} .
$$

Therefore

$$
g(t)=1-\cos t .
$$

From this, (8.5.6), and Theorem 8.4.2,

$$
y=1-\cos t-2 u\left(t-\frac{\pi}{2}\right)\left(1-\cos \left(t-\frac{\pi}{2}\right)\right)+2 \cos t-\sin t
$$

Simplifying this (recalling that $\cos (t-\pi / 2)=\sin t$ ) yields

$$
y=1+\cos t-\sin t-2 u\left(t-\frac{\pi}{2}\right)(1-\sin t)
$$

or

$$
y=\left\{\begin{aligned}
1+\cos t-\sin t, & 0 \leq t<\frac{\pi}{2}, \\
-1+\cos t+\sin t, & t \geq \frac{\pi}{2},
\end{aligned}\right.
$$

which is the result obtained in Example 8.5.1.
REMARK: It isn't obvious that using the Laplace transform to solve (8.5.2) as we did in Example 8.5.2 yields a function $y$ with the properties stated in Theorem 8.5.1; that is, such that $y$ and $y^{\prime}$ are continuous on $[0, \infty)$ and $y^{\prime \prime}$ has limits from the right and left at $t_{1}$. However, this is true if $f_{0}$ and $f_{1}$ are continuous and of exponential order on $[0, \infty)$. A proof is sketched in Exercises 8.6.11-8.613.

Example 8.5.3 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}-y=f(t), \quad y(0)=-1, y^{\prime}(0)=2, \tag{8.5.8}
\end{equation*}
$$

where

$$
f(t)= \begin{cases}t, & 0 \leq t<1 \\ 1, & t \geq 1\end{cases}
$$

Solution Here

$$
f(t)=t-u(t-1)(t-1)
$$

so

$$
\begin{aligned}
\mathcal{L}(f) & =\mathcal{L}(t)-\mathcal{L}(u(t-1)(t-1)) \\
& =\mathcal{L}(t)-e^{-s} \mathcal{L}(t)(\text { from Theorem 8.4.1) } \\
& =\frac{1}{s^{2}}-\frac{e^{-s}}{s^{2}}
\end{aligned}
$$

Since transforming (8.5.8) yields

$$
\left(s^{2}-1\right) Y(s)=\mathcal{L}(f)+2-s
$$

we see that

$$
\begin{equation*}
Y(s)=\left(1-e^{-s}\right) H(s)+\frac{2-s}{s^{2}-1} \tag{8.5.9}
\end{equation*}
$$

where

$$
H(s)=\frac{1}{s^{2}\left(s^{2}-1\right)}=\frac{1}{s^{2}-1}-\frac{1}{s^{2}}
$$

therefore

$$
\begin{equation*}
h(t)=\sinh t-t . \tag{8.5.10}
\end{equation*}
$$

Since

$$
\mathcal{L}^{-1}\left(\frac{2-s}{s^{2}-1}\right)=2 \sinh t-\cosh t
$$

we conclude from (8.5.9), (8.5.10), and Theorem 8.4.1 that

$$
y=\sinh t-t-u(t-1)(\sinh (t-1)-t+1)+2 \sinh t-\cosh t
$$

or

$$
\begin{equation*}
y=3 \sinh t-\cosh t-t-u(t-1)(\sinh (t-1)-t+1) \tag{8.5.11}
\end{equation*}
$$

We leave it to you to verify that $y$ and $y^{\prime}$ are continuous and $y^{\prime \prime}$ has limits from the right and left at $t_{1}=1$.
Example 8.5.4 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+y=f(t), \quad y(0)=0, y^{\prime}(0)=0 \tag{8.5.12}
\end{equation*}
$$

where

$$
f(t)=\left\{\begin{array}{cl}
0, & 0 \leq t<\frac{\pi}{4} \\
\cos 2 t, & \frac{\pi}{4} \leq t<\pi \\
0, & t \geq \pi
\end{array}\right.
$$

Solution Here

$$
f(t)=u(t-\pi / 4) \cos 2 t-u(t-\pi) \cos 2 t
$$

so

$$
\begin{aligned}
\mathcal{L}(f) & =\mathcal{L}(u(t-\pi / 4) \cos 2 t)-\mathcal{L}(u(t-\pi) \cos 2 t) \\
& =e^{-\pi s / 4} \mathcal{L}(\cos 2(t+\pi / 4))-e^{-\pi s} \mathcal{L}(\cos 2(t+\pi)) \\
& =-e^{-\pi s / 4} \mathcal{L}(\sin 2 t)-e^{-\pi s} \mathcal{L}(\cos 2 t) \\
& =-\frac{2 e^{-\pi s / 4}}{s^{2}+4}-\frac{s e^{-\pi s}}{s^{2}+4}
\end{aligned}
$$

Since transforming (8.5.12) yields

$$
\left(s^{2}+1\right) Y(s)=\mathcal{L}(f),
$$

we see that

$$
\begin{equation*}
Y(s)=e^{-\pi s / 4} H_{1}(s)+e^{-\pi s} H_{2}(s), \tag{8.5.13}
\end{equation*}
$$

where

$$
\begin{equation*}
H_{1}(s)=-\frac{2}{\left(s^{2}+1\right)\left(s^{2}+4\right)} \quad \text { and } \quad H_{2}(s)=-\frac{s}{\left(s^{2}+1\right)\left(s^{2}+4\right)} \tag{8.5.14}
\end{equation*}
$$

To simplify the required partial fraction expansions, we first write

$$
\frac{1}{(x+1)(x+4)}=\frac{1}{3}\left[\frac{1}{x+1}-\frac{1}{x+4}\right] .
$$

Setting $x=s^{2}$ and substituting the result in (8.5.14) yields

$$
H_{1}(s)=-\frac{2}{3}\left[\frac{1}{s^{2}+1}-\frac{1}{s^{2}+4}\right] \quad \text { and } \quad H_{2}(s)=-\frac{1}{3}\left[\frac{s}{s^{2}+1}-\frac{s}{s^{2}+4}\right]
$$

The inverse transforms are

$$
h_{1}(t)=-\frac{2}{3} \sin t+\frac{1}{3} \sin 2 t \quad \text { and } \quad h_{2}(t)=-\frac{1}{3} \cos t+\frac{1}{3} \cos 2 t .
$$

From (8.5.13) and Theorem 8.4.2,

$$
\begin{equation*}
y=u\left(t-\frac{\pi}{4}\right) h_{1}\left(t-\frac{\pi}{4}\right)+u(t-\pi) h_{2}(t-\pi) . \tag{8.5.15}
\end{equation*}
$$

Since

$$
\begin{aligned}
h_{1}\left(t-\frac{\pi}{4}\right) & =-\frac{2}{3} \sin \left(t-\frac{\pi}{4}\right)+\frac{1}{3} \sin 2\left(t-\frac{\pi}{4}\right) \\
& =-\frac{\sqrt{2}}{3}(\sin t-\cos t)-\frac{1}{3} \cos 2 t
\end{aligned}
$$

and

$$
\begin{aligned}
h_{2}(t-\pi) & =-\frac{1}{3} \cos (t-\pi)+\frac{1}{3} \cos 2(t-\pi) \\
& =\frac{1}{3} \cos t+\frac{1}{3} \cos 2 t
\end{aligned}
$$

(8.5.15) can be rewritten as

$$
y=-\frac{1}{3} u\left(t-\frac{\pi}{4}\right)(\sqrt{2}(\sin t-\cos t)+\cos 2 t)+\frac{1}{3} u(t-\pi)(\cos t+\cos 2 t)
$$

or

$$
y=\left\{\begin{array}{cl}
0, & 0 \leq t<\frac{\pi}{4}  \tag{8.5.16}\\
-\frac{\sqrt{2}}{3}(\sin t-\cos t)-\frac{1}{3} \cos 2 t, & \frac{\pi}{4} \leq t<\pi \\
-\frac{\sqrt{2}}{3} \sin t+\frac{1+\sqrt{2}}{3} \cos t, & t \geq \pi
\end{array}\right.
$$

We leave it to you to verify that $y$ and $y^{\prime}$ are continuous and $y^{\prime \prime}$ has limits from the right and left at $t_{1}=\pi / 4$ and $t_{2}=\pi$ (Figure 8.5.2).

### 8.5 Exercises

In Exercises 1-20 use the Laplace transform to solve the initial value problem. Where indicated by C/G, graph the solution.

1. $y^{\prime \prime}+y=\left\{\begin{array}{ll}3, & 0 \leq t<\pi, \\ 0, & t \geq \pi,\end{array} \quad y(0)=0, \quad y^{\prime}(0)=0\right.$


Figure 8.5.2 Graph of (8.5.16)
2. $y^{\prime \prime}+y=\left\{\begin{array}{cl}3, & 0 \leq t<4, \\ ; 2 t-5, & t>4,\end{array} \quad y(0)=1, \quad y^{\prime}(0)=0\right.$
3. $y^{\prime \prime}-2 y^{\prime}=\left\{\begin{array}{ll}4, & 0 \leq t<1, \\ 6, & t \geq 1,\end{array} \quad y(0)=-6, \quad y^{\prime}(0)=1\right.$
4. $y^{\prime \prime}-y=\left\{\begin{array}{cl}e^{2 t}, & 0 \leq t<2, \\ 1, & t \geq 2,\end{array} \quad y(0)=3, \quad y^{\prime}(0)=-1\right.$
5. $y^{\prime \prime}-3 y^{\prime}+2 y=\left\{\begin{array}{rl}0, & 0 \leq t<1, \\ 1, & 1 \leq t<2, \\ -1, & t \geq 2,\end{array} \quad y(0)=-3, \quad y^{\prime}(0)=1\right.$
6. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y=\left\{\begin{array}{cl}|\sin t|, & 0 \leq t<2 \pi, \\ 0, & t \geq 2 \pi,\end{array} \quad y(0)=-3, \quad y^{\prime}(0)=1\right.$
7. $y^{\prime \prime}-5 y^{\prime}+4 y=\left\{\begin{array}{rl}1, & 0 \leq t<1 \\ -1, & 1 \leq t<2, \\ 0, & t \geq 2,\end{array} \quad y(0)=3, \quad y^{\prime}(0)=-5\right.$
8. $y^{\prime \prime}+9 y=\left\{\begin{array}{ll}\cos t, & 0 \leq t<\frac{3 \pi}{2}, \\ \sin t, & t \geq \frac{3 \pi}{2},\end{array} \quad y(0)=0, y^{\prime}(0)=0\right.$
9. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y=\left\{\begin{array}{ll}t, & 0 \leq t<\frac{\pi}{2}, \\ \pi, & t \geq \frac{\pi}{2},\end{array} \quad y(0)=0, \quad y^{\prime}(0)=0\right.$
10. $y^{\prime \prime}+y=\left\{\begin{array}{rl}t, & 0 \leq t<\pi, \\ -t, & t \geq \pi,\end{array} \quad y(0)=0, y^{\prime}(0)=0\right.$
11. $\quad y^{\prime \prime}-3 y^{\prime}+2 y=\left\{\begin{array}{cl}0, & 0 \leq t<2, \\ 2 t-4, & t \geq 2,\end{array}, \quad y(0)=0, \quad y^{\prime}(0)=0\right.$
12. $\quad y^{\prime \prime}+y=\left\{\begin{array}{cl}t, & 0 \leq t<2 \pi, \\ -2 t, & t \geq 2 \pi,\end{array} \quad y(0)=1, \quad y^{\prime}(0)=2\right.$
13. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+3 y^{\prime}+2 y=\left\{\begin{array}{rl}1, & 0 \leq t<2, \\ -1, & t \geq 2,\end{array} \quad y(0)=0, \quad y^{\prime}(0)=0\right.$
14. $\quad y^{\prime \prime}-4 y^{\prime}+3 y=\left\{\begin{array}{rl}-1, & 0 \leq t<1, \\ 1, & t \geq 1,\end{array} \quad y(0)=0, y^{\prime}(0)=0\right.$
15. $y^{\prime \prime}+2 y^{\prime}+y=\left\{\begin{array}{cl}e^{t}, & 0 \leq t<1, \\ e^{t}-1, & t \geq 1,\end{array} \quad y(0)=3, y^{\prime}(0)=-1\right.$
16. $y^{\prime \prime}+2 y^{\prime}+y=\left\{\begin{array}{cl}4 e^{t}, & 0 \leq t<1, \\ 0, & t \geq 1,\end{array} \quad y(0)=0, y^{\prime}(0)=0\right.$
17. $y^{\prime \prime}+3 y^{\prime}+2 y=\left\{\begin{array}{cl}e^{-t}, & 0 \leq t<1, \\ 0, & t \geq 1,\end{array} \quad y(0)=1, y^{\prime}(0)=-1\right.$
18. $\quad y^{\prime \prime}-4 y^{\prime}+4 y=\left\{\begin{array}{rl}e^{2 t}, & 0 \leq t<2, \\ -e^{2 t}, & t \geq 2,\end{array} \quad y(0)=0, y^{\prime}(0)=-1\right.$
19. $\mathrm{C} / \mathrm{G} y^{\prime \prime}=\left\{\begin{array}{cl}t^{2}, & 0 \leq t<1, \\ -t, & 1 \leq t<2, \quad y(0)=1, y^{\prime}(0)=0 \\ t+1, & t \geq 2,\end{array}\right.$
20. $\quad y^{\prime \prime}+2 y^{\prime}+2 y=\left\{\begin{aligned} 1, & 0 \leq t<2 \pi, \\ t, & 2 \pi \leq t<3 \pi, \quad y(0)=2, \quad y^{\prime}(0)=-1 \\ -1, & t \geq 3 \pi,\end{aligned}\right.$
21. Solve the initial value problem

$$
y^{\prime \prime}=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0,
$$

where

$$
f(t)=m+1, \quad m \leq t<m+1, \quad m=0,1,2, \ldots
$$

22. Solve the given initial value problem and find a formula that does not involve step functions and represents $y$ on each interval of continuity of $f$.
(a) $y^{\prime \prime}+y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$;
$f(t)=m+1, \quad m \pi \leq t<(m+1) \pi, \quad m=0,1,2, \ldots$.
(b) $y^{\prime \prime}+y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$;
$f(t)=(m+1) t, \quad 2 m \pi \leq t<2(m+1) \pi, \quad m=0,1,2, \ldots$ Hint: You'll need the formula

$$
1+2+\cdots+m=\frac{m(m+1)}{2}
$$

(c) $\quad y^{\prime \prime}+y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$;
$f(t)=(-1)^{m}, \quad m \pi \leq t<(m+1) \pi, \quad m=0,1,2, \ldots$
(d) $y^{\prime \prime}-y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$;
$f(t)=m+1, \quad m \leq t<(m+1), \quad m=0,1,2, \ldots$.
Hint: You will need the formula

$$
1+r+\cdots+r^{m}=\frac{1-r^{m+1}}{1-r}(r \neq 1)
$$

(e) $y^{\prime \prime}+2 y^{\prime}+2 y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$;
$f(t)=(m+1)(\sin t+2 \cos t), \quad 2 m \pi \leq t<2(m+1) \pi, \quad m=0,1,2, \ldots$.
(See the hint in (d).)
(f) $y^{\prime \prime}-3 y^{\prime}+2 y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$;
$f(t)=m+1, \quad m \leq t<m+1, \quad m=0,1,2, \ldots$.
(See the hints in (b) and (d).)
23. (a) Let $g$ be continuous on $(\alpha, \beta)$ and differentiable on the $\left(\alpha, t_{0}\right)$ and $\left(t_{0}, \beta\right)$. Suppose $A=$ $\lim _{t \rightarrow t_{0}-} g^{\prime}(t)$ and $B=\lim _{t \rightarrow t_{0}+} g^{\prime}(t)$ both exist. Use the mean value theorem to show that

$$
\lim _{t \rightarrow t_{0}-} \frac{g(t)-g\left(t_{0}\right)}{t-t_{0}}=A \quad \text { and } \quad \lim _{t \rightarrow t_{0}+} \frac{g(t)-g\left(t_{0}\right)}{t-t_{0}}=B
$$

(b) Conclude from (a) that $g^{\prime}\left(t_{0}\right)$ exists and $g^{\prime}$ is continuous at $t_{0}$ if $A=B$.
(c) Conclude from (a) that if $g$ is differentiable on $(\alpha, \beta)$ then $g^{\prime}$ can't have a jump discontinuity on $(\alpha, \beta)$.
24. (a) Let $a, b$, and $c$ be constants, with $a \neq 0$. Let $f$ be piecewise continuous on an interval $(\alpha, \beta)$, with a single jump discontinuity at a point $t_{0}$ in $(\alpha, \beta)$. Suppose $y$ and $y^{\prime}$ are continuous on $(\alpha, \beta)$ and $y^{\prime \prime}$ on $\left(\alpha, t_{0}\right)$ and $\left(t_{0}, \beta\right)$. Suppose also that

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=f(t) \tag{A}
\end{equation*}
$$

on $\left(\alpha, t_{0}\right)$ and $\left(t_{0}, \beta\right)$. Show that

$$
y^{\prime \prime}\left(t_{0}+\right)-y^{\prime \prime}\left(t_{0}-\right)=\frac{f\left(t_{0}+\right)-f\left(t_{0}-\right)}{a} \neq 0
$$

(b) Use (a) and Exercise 23(c) to show that (A) does not have solutions on any interval $(\alpha, \beta)$ that contains a jump discontinuity of $f$.
25. Suppose $P_{0}, P_{1}$, and $P_{2}$ are continuous and $P_{0}$ has no zeros on an open interval $(a, b)$, and that $F$ has a jump discontinuity at a point $t_{0}$ in $(a, b)$. Show that the differential equation

$$
P_{0}(t) y^{\prime \prime}+P_{1}(t) y^{\prime}+P_{2}(t) y=F(t)
$$

has no solutions on $(a, b)$.Hint: Generalize the result of Exercise 24 and use Exercise 23(c).
26. Let $0=t_{0}<t_{1}<\cdots<t_{n}$. Suppose $f_{m}$ is continuous on $\left[t_{m}, \infty\right)$ for $m=1, \ldots, n$. Let

$$
f(t)= \begin{cases}f_{m}(t), & t_{m} \leq t<t_{m+1}, \quad m=1, \ldots, n-1 \\ f_{n}(t), & t \geq t_{n}\end{cases}
$$

Show that the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

as defined following Theorem 8.5.1, is given by

$$
y=\left\{\begin{array}{cl}
z_{0}(t), & 0 \leq t<t_{1} \\
z_{0}(t)+z_{1}(t), & t_{1} \leq t<t_{2} \\
& \vdots \\
z_{0}+\cdots+z_{n-1}(t), & t_{n-1} \leq t<t_{n} \\
z_{0}+\cdots+z_{n}(t), & t \geq t_{n}
\end{array}\right.
$$

where $z_{0}$ is the solution of

$$
a z^{\prime \prime}+b z^{\prime}+c z=f_{0}(t), \quad z(0)=k_{0}, \quad z^{\prime}(0)=k_{1}
$$

and $z_{m}$ is the solution of

$$
a z^{\prime \prime}+b z^{\prime}+c z=f_{m}(t)-f_{m-1}(t), \quad z\left(t_{m}\right)=0, \quad z^{\prime}\left(t_{m}\right)=0
$$

for $m=1, \ldots, n$.

### 8.6 CONVOLUTION

In this section we consider the problem of finding the inverse Laplace transform of a product $H(s)=$ $F(s) G(s)$, where $F$ and $G$ are the Laplace transforms of known functions $f$ and $g$. To motivate our interest in this problem, consider the initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0
$$

Taking Laplace transforms yields

$$
\left(a s^{2}+b s+c\right) Y(s)=F(s)
$$

so

$$
\begin{equation*}
Y(s)=F(s) G(s) \tag{8.6.1}
\end{equation*}
$$

where

$$
G(s)=\frac{1}{a s^{2}+b s+c}
$$

Until now wen't been interested in the factorization indicated in (8.6.1), since we dealt only with differential equations with specific forcing functions. Hence, we could simply do the indicated multiplication in (8.6.1) and use the table of Laplace transforms to find $y=\mathcal{L}^{-1}(Y)$. However, this isn't possible if we want a formula for $y$ in terms of $f$, which may be unspecified.

To motivate the formula for $\mathcal{L}^{-1}(F G)$, consider the initial value problem

$$
\begin{equation*}
y^{\prime}-a y=f(t), \quad y(0)=0 \tag{8.6.2}
\end{equation*}
$$

which we first solve without using the Laplace transform. The solution of the differential equation in (8.6.2) is of the form $y=u e^{a t}$ where

$$
u^{\prime}=e^{-a t} f(t)
$$

Integrating this from 0 to $t$ and imposing the initial condition $u(0)=y(0)=0$ yields

$$
u=\int_{0}^{t} e^{-a \tau} f(\tau) d \tau
$$

Therefore

$$
\begin{equation*}
y(t)=e^{a t} \int_{0}^{t} e^{-a \tau} f(\tau) d \tau=\int_{0}^{t} e^{a(t-\tau)} f(\tau) d \tau \tag{8.6.3}
\end{equation*}
$$

Now we'll use the Laplace transform to solve (8.6.2) and compare the result to (8.6.3). Taking Laplace transforms in (8.6.2) yields

$$
(s-a) Y(s)=F(s),
$$

so

$$
Y(s)=F(s) \frac{1}{s-a}
$$

which implies that

$$
\begin{equation*}
y(t)=\mathcal{L}^{-1}\left(F(s) \frac{1}{s-a}\right) \tag{8.6.4}
\end{equation*}
$$

If we now let $g(t)=e^{a t}$, so that

$$
G(s)=\frac{1}{s-a}
$$

then (8.6.3) and (8.6.4) can be written as

$$
y(t)=\int_{0}^{t} f(\tau) g(t-\tau) d \tau
$$

and

$$
y=\mathcal{L}^{-1}(F G)
$$

respectively. Therefore

$$
\begin{equation*}
\mathcal{L}^{-1}(F G)=\int_{0}^{t} f(\tau) g(t-\tau) d \tau \tag{8.6.5}
\end{equation*}
$$

in this case.
This motivates the next definition.
Definition 8.6.1 The convolution $f * g$ of two functions $f$ and $g$ is defined by

$$
(f * g)(t)=\int_{0}^{t} f(\tau) g(t-\tau) d \tau
$$

It can be shown (Exercise 6) that $f * g=g * f$; that is,

$$
\int_{0}^{t} f(t-\tau) g(\tau) d \tau=\int_{0}^{t} f(\tau) g(t-\tau) d \tau
$$

Eqn. (8.6.5) shows that $\mathcal{L}^{-1}(F G)=f * g$ in the special case where $g(t)=e^{a t}$. This next theorem states that this is true in general.


Figure 8.6.1

Theorem 8.6.2 [The Convolution Theorem] If $\mathcal{L}(f)=F$ and $\mathcal{L}(g)=G$, then

$$
\mathcal{L}(f * g)=F G .
$$

A complete proof of the convolution theorem is beyond the scope of this book. However, we'll assume that $f * g$ has a Laplace transform and verify the conclusion of the theorem in a purely computational way. By the definition of the Laplace transform,

$$
\mathcal{L}(f * g)=\int_{0}^{\infty} e^{-s t}(f * g)(t) d t=\int_{0}^{\infty} e^{-s t} \int_{0}^{t} f(\tau) g(t-\tau) d \tau d t
$$

This iterated integral equals a double integral over the region shown in Figure 8.6.1. Reversing the order of integration yields

$$
\begin{equation*}
\mathcal{L}(f * g)=\int_{0}^{\infty} f(\tau) \int_{\tau}^{\infty} e^{-s t} g(t-\tau) d t d \tau \tag{8.6.6}
\end{equation*}
$$

However, the substitution $x=t-\tau$ shows that

$$
\begin{aligned}
\int_{\tau}^{\infty} e^{-s t} g(t-\tau) d t & =\int_{0}^{\infty} e^{-s(x+\tau)} g(x) d x \\
& =e^{-s \tau} \int_{0}^{\infty} e^{-s x} g(x) d x=e^{-s \tau} G(s)
\end{aligned}
$$

Substituting this into (8.6.6) and noting that $G(s)$ is independent of $\tau$ yields

$$
\begin{aligned}
\mathcal{L}(f * g) & =\int_{0}^{\infty} e^{-s \tau} f(\tau) G(s) d \tau \\
& =G(s) \int_{0}^{\infty} e^{-s t} f(\tau) d \tau=F(s) G(s)
\end{aligned}
$$

Example 8.6.1 Let

$$
f(t)=e^{a t} \quad \text { and } \quad g(t)=e^{b t} \quad(a \neq b)
$$

Verify that $\mathcal{L}(f * g)=\mathcal{L}(f) \mathcal{L}(g)$, as implied by the convolution theorem.

Solution We first compute

$$
\begin{aligned}
(f * g)(t) & =\int_{0}^{t} e^{a \tau} e^{b(t-\tau)} d \tau=e^{b t} \int_{0}^{t} e^{(a-b) \tau} d \tau \\
& =\left.e^{b t} \frac{e^{(a-b) \tau}}{a-b}\right|_{0} ^{t}=\frac{e^{b t}\left[e^{(a-b) t}-1\right]}{a-b} \\
& =\frac{e^{a t}-e^{b t}}{a-b}
\end{aligned}
$$

Since

$$
e^{a t} \leftrightarrow \frac{1}{s-a} \quad \text { and } \quad e^{b t} \leftrightarrow \frac{1}{s-b},
$$

it follows that

$$
\begin{aligned}
\mathcal{L}(f * g) & =\frac{1}{a-b}\left[\frac{1}{s-a}-\frac{1}{s-b}\right] \\
& =\frac{1}{(s-a)(s-b)} \\
& =\mathcal{L}\left(e^{a t}\right) \mathcal{L}\left(e^{b t}\right)=\mathcal{L}(f) \mathcal{L}(g)
\end{aligned}
$$

A Formula for the Solution of an Initial Value Problem
The convolution theorem provides a formula for the solution of an initial value problem for a linear constant coefficient second order equation with an unspecified. The next three examples illustrate this.

Example 8.6.2 Find a formula for the solution of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} . \tag{8.6.7}
\end{equation*}
$$

Solution Taking Laplace transforms in (8.6.7) yields

$$
\left(s^{2}-2 s+1\right) Y(s)=F(s)+\left(k_{1}+k_{0} s\right)-2 k_{0} .
$$

Therefore

$$
\begin{aligned}
Y(s) & =\frac{1}{(s-1)^{2}} F(s)+\frac{k_{1}+k_{0} s-2 k_{0}}{(s-1)^{2}} \\
& =\frac{1}{(s-1)^{2}} F(s)+\frac{k_{0}}{s-1}+\frac{k_{1}-k_{0}}{(s-1)^{2}}
\end{aligned}
$$

From the table of Laplace transforms,

$$
\mathcal{L}^{-1}\left(\frac{k_{0}}{s-1}+\frac{k_{1}-k_{0}}{(s-1)^{2}}\right)=e^{t}\left(k_{0}+\left(k_{1}-k_{0}\right) t\right) .
$$

Since

$$
\frac{1}{(s-1)^{2}} \leftrightarrow t e^{t} \quad \text { and } \quad F(s) \leftrightarrow f(t)
$$

the convolution theorem implies that

$$
\mathcal{L}^{-1}\left(\frac{1}{(s-1)^{2}} F(s)\right)=\int_{0}^{t} \tau e^{\tau} f(t-\tau) d \tau
$$

Therefore the solution of (8.6.7) is

$$
y(t)=e^{t}\left(k_{0}+\left(k_{1}-k_{0}\right) t\right)+\int_{0}^{t} \tau e^{\tau} f(t-\tau) d \tau
$$

Example 8.6.3 Find a formula for the solution of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+4 y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} . \tag{8.6.8}
\end{equation*}
$$

Solution Taking Laplace transforms in (8.6.8) yields

$$
\left(s^{2}+4\right) Y(s)=F(s)+k_{1}+k_{0} s
$$

Therefore

$$
Y(s)=\frac{1}{\left(s^{2}+4\right)} F(s)+\frac{k_{1}+k_{0} s}{s^{2}+4} .
$$

From the table of Laplace transforms,

$$
\mathcal{L}^{-1}\left(\frac{k_{1}+k_{0} s}{s^{2}+4}\right)=k_{0} \cos 2 t+\frac{k_{1}}{2} \sin 2 t .
$$

Since

$$
\frac{1}{\left(s^{2}+4\right)} \leftrightarrow \frac{1}{2} \sin 2 t \quad \text { and } \quad F(s) \leftrightarrow f(t)
$$

the convolution theorem implies that

$$
\mathcal{L}^{-1}\left(\frac{1}{\left(s^{2}+4\right)} F(s)\right)=\frac{1}{2} \int_{0}^{t} f(t-\tau) \sin 2 \tau d \tau
$$

Therefore the solution of (8.6.8) is

$$
y(t)=k_{0} \cos 2 t+\frac{k_{1}}{2} \sin 2 t+\frac{1}{2} \int_{0}^{t} f(t-\tau) \sin 2 \tau d \tau
$$

Example 8.6.4 Find a formula for the solution of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+2 y^{\prime}+2 y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} . \tag{8.6.9}
\end{equation*}
$$

Solution Taking Laplace transforms in (8.6.9) yields

$$
\left(s^{2}+2 s+2\right) Y(s)=F(s)+k_{1}+k_{0} s+2 k_{0} .
$$

Therefore

$$
\begin{aligned}
Y(s) & =\frac{1}{(s+1)^{2}+1} F(s)+\frac{k_{1}+k_{0} s+2 k_{0}}{(s+1)^{2}+1} \\
& =\frac{1}{(s+1)^{2}+1} F(s)+\frac{\left(k_{1}+k_{0}\right)+k_{0}(s+1)}{(s+1)^{2}+1}
\end{aligned}
$$

From the table of Laplace transforms,

$$
\mathcal{L}^{-1}\left(\frac{\left(k_{1}+k_{0}\right)+k_{0}(s+1)}{(s+1)^{2}+1}\right)=e^{-t}\left(\left(k_{1}+k_{0}\right) \sin t+k_{0} \cos t\right)
$$

Since

$$
\frac{1}{(s+1)^{2}+1} \leftrightarrow e^{-t} \sin t \quad \text { and } \quad F(s) \leftrightarrow f(t)
$$

the convolution theorem implies that

$$
\mathcal{L}^{-1}\left(\frac{1}{(s+1)^{2}+1} F(s)\right)=\int_{0}^{t} f(t-\tau) e^{-\tau} \sin \tau d \tau
$$

Therefore the solution of (8.6.9) is

$$
\begin{equation*}
y(t)=e^{-t}\left(\left(k_{1}+k_{0}\right) \sin t+k_{0} \cos t\right)+\int_{0}^{t} f(t-\tau) e^{-\tau} \sin \tau d \tau \tag{8.6.10}
\end{equation*}
$$

Evaluating Convolution Integrals
We'll say that an integral of the form $\int_{0}^{t} u(\tau) v(t-\tau) d \tau$ is a convolution integral. The convolution theorem provides a convenient way to evaluate convolution integrals.

Example 8.6.5 Evaluate the convolution integral

$$
h(t)=\int_{0}^{t}(t-\tau)^{5} \tau^{7} d \tau
$$

Solution We could evaluate this integral by expanding $(t-\tau)^{5}$ in powers of $\tau$ and then integrating. However, the convolution theorem provides an easier way. The integral is the convolution of $f(t)=t^{5}$ and $g(t)=t^{7}$. Since

$$
t^{5} \leftrightarrow \frac{5!}{s^{6}} \quad \text { and } \quad t^{7} \leftrightarrow \frac{7!}{s^{8}}
$$

the convolution theorem implies that

$$
h(t) \leftrightarrow \frac{5!7!}{s^{14}}=\frac{5!7!}{13!} \frac{13!}{s^{14}}
$$

where we have written the second equality because

$$
\frac{13!}{s^{14}} \leftrightarrow t^{13}
$$

Hence,

$$
h(t)=\frac{5!7!}{13!} t^{13}
$$

Example 8.6.6 Use the convolution theorem and a partial fraction expansion to evaluate the convolution integral

$$
h(t)=\int_{0}^{t} \sin a(t-\tau) \cos b \tau d \tau \quad(|a| \neq|b|)
$$

Solution Since

$$
\sin a t \leftrightarrow \frac{a}{s^{2}+a^{2}} \quad \text { and } \quad \cos b t \leftrightarrow \frac{s}{s^{2}+b^{2}},
$$

the convolution theorem implies that

$$
H(s)=\frac{a}{s^{2}+a^{2}} \frac{s}{s^{2}+b^{2}} .
$$

Expanding this in a partial fraction expansion yields

$$
H(s)=\frac{a}{b^{2}-a^{2}}\left[\frac{s}{s^{2}+a^{2}}-\frac{s}{s^{2}+b^{2}}\right]
$$

Therefore

$$
h(t)=\frac{a}{b^{2}-a^{2}}(\cos a t-\cos b t) .
$$

Volterra Integral Equations
An equation of the form

$$
\begin{equation*}
y(t)=f(t)+\int_{0}^{t} k(t-\tau) y(\tau) d \tau \tag{8.6.11}
\end{equation*}
$$

is a Volterra integral equation. Here $f$ and $k$ are given functions and $y$ is unknown. Since the integral on the right is a convolution integral, the convolution theorem provides a convenient formula for solving (8.6.11). Taking Laplace transforms in (8.6.11) yields

$$
Y(s)=F(s)+K(s) Y(s),
$$

and solving this for $Y(s)$ yields

$$
Y(s)=\frac{F(s)}{1-K(s)}
$$

We then obtain the solution of (8.6.11) as $y=\mathcal{L}^{-1}(Y)$.

Example 8.6.7 Solve the integral equation

$$
\begin{equation*}
y(t)=1+2 \int_{0}^{t} e^{-2(t-\tau)} y(\tau) d \tau \tag{8.6.12}
\end{equation*}
$$

Solution Taking Laplace transforms in (8.6.12) yields

$$
Y(s)=\frac{1}{s}+\frac{2}{s+2} Y(s)
$$

and solving this for $Y(s)$ yields

$$
Y(s)=\frac{1}{s}+\frac{2}{s^{2}}
$$

Hence,

$$
y(t)=1+2 t
$$

Transfer Functions
The next theorem presents a formula for the solution of the general initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

where we assume for simplicity that $f$ is continuous on $[0, \infty)$ and that $\mathcal{L}(f)$ exists. In Exercises $11-14$ it's shown that the formula is valid under much weaker conditions on $f$.

Theorem 8.6.3 Suppose $f$ is continuous on $[0, \infty)$ and has a Laplace transform. Then the solution of the initial value problem

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} \tag{8.6.13}
\end{equation*}
$$

is

$$
\begin{equation*}
y(t)=k_{0} y_{1}(t)+k_{1} y_{2}(t)+\int_{0}^{t} w(\tau) f(t-\tau) d \tau \tag{8.6.14}
\end{equation*}
$$

where $y_{1}$ and $y_{2}$ satisfy

$$
\begin{equation*}
a y_{1}^{\prime \prime}+b y_{1}^{\prime}+c y_{1}=0, \quad y_{1}(0)=1, \quad y_{1}^{\prime}(0)=0 \tag{8.6.15}
\end{equation*}
$$

and

$$
\begin{equation*}
a y_{2}^{\prime \prime}+b y_{2}^{\prime}+c y_{2}=0, \quad y_{2}(0)=0, \quad y_{2}^{\prime}(0)=1 \tag{8.6.16}
\end{equation*}
$$

and

$$
\begin{equation*}
w(t)=\frac{1}{a} y_{2}(t) . \tag{8.6.17}
\end{equation*}
$$

Proof Taking Laplace transforms in (8.6.13) yields

$$
p(s) Y(s)=F(s)+a\left(k_{1}+k_{0} s\right)+b k_{0}
$$

where

$$
p(s)=a s^{2}+b s+c
$$

Hence,

$$
\begin{equation*}
Y(s)=W(s) F(s)+V(s) \tag{8.6.18}
\end{equation*}
$$

with

$$
\begin{equation*}
W(s)=\frac{1}{p(s)} \tag{8.6.19}
\end{equation*}
$$

and

$$
\begin{equation*}
V(s)=\frac{a\left(k_{1}+k_{0} s\right)+b k_{0}}{p(s)} \tag{8.6.20}
\end{equation*}
$$

Taking Laplace transforms in (8.6.15) and (8.6.16) shows that

$$
p(s) Y_{1}(s)=a s+b \quad \text { and } \quad p(s) Y_{2}(s)=a
$$

Therefore

$$
Y_{1}(s)=\frac{a s+b}{p(s)}
$$

and

$$
\begin{equation*}
Y_{2}(s)=\frac{a}{p(s)} \tag{8.6.21}
\end{equation*}
$$

Hence, (8.6.20) can be rewritten as

$$
V(s)=k_{0} Y_{1}(s)+k_{1} Y_{2}(s)
$$

Substituting this into (8.6.18) yields

$$
Y(s)=k_{0} Y_{1}(s)+k_{1} Y_{2}(s)+\frac{1}{a} Y_{2}(s) F(s) .
$$

Taking inverse transforms and invoking the convolution theorem yields (8.6.14). Finally, (8.6.19) and (8.6.21) imply (8.6.17).

It is useful to note from (8.6.14) that $y$ is of the form

$$
y=v+h
$$

where

$$
v(t)=k_{0} y_{1}(t)+k_{1} y_{2}(t)
$$

depends on the initial conditions and is independent of the forcing function, while

$$
h(t)=\int_{0}^{t} w(\tau) f(t-\tau) d \tau
$$

depends on the forcing function and is independent of the initial conditions. If the zeros of the characteristic polynomial

$$
p(s)=a s^{2}+b s+c
$$

of the complementary equation have negative real parts, then $y_{1}$ and $y_{2}$ both approach zero as $t \rightarrow \infty$, so $\lim _{t \rightarrow \infty} v(t)=0$ for any choice of initial conditions. Moreover, the value of $h(t)$ is essentially independent of the values of $f(t-\tau)$ for large $\tau$, since $\lim _{\tau \rightarrow \infty} w(\tau)=0$. In this case we say that $v$ and $h$ are transient and steady state components, respectively, of the solution $y$ of (8.6.13). These definitions apply to the initial value problem of Example 8.6.4, where the zeros of

$$
p(s)=s^{2}+2 s+2=(s+1)^{2}+1
$$

are $-1 \pm i$. From (8.6.10), we see that the solution of the general initial value problem of Example 8.6.4 is $y=v+h$, where

$$
v(t)=e^{-t}\left(\left(k_{1}+k_{0}\right) \sin t+k_{0} \cos t\right)
$$

is the transient component of the solution and

$$
h(t)=\int_{0}^{t} f(t-\tau) e^{-\tau} \sin \tau d \tau
$$

is the steady state component. The definitions don't apply to the initial value problems considered in Examples 8.6.2 and 8.6.3, since the zeros of the characteristic polynomials in these two examples don't have negative real parts.

In physical applications where the input $f$ and the output $y$ of a device are related by (8.6.13), the zeros of the characteristic polynomial usually do have negative real parts. Then $W=\mathcal{L}(w)$ is called the transfer function of the device. Since

$$
H(s)=W(s) F(s)
$$

we see that

$$
W(s)=\frac{H(s)}{F(s)}
$$

is the ratio of the transform of the steady state output to the transform of the input.

Because of the form of

$$
h(t)=\int_{0}^{t} w(\tau) f(t-\tau) d \tau
$$

$w$ is sometimes called the weighting function of the device, since it assigns weights to past values of the input $f$. It is also called the impulse response of the device, for reasons discussed in the next section.

Formula (8.6.14) is given in more detail in Exercises 8-10 for the three possible cases where the zeros of $p(s)$ are real and distinct, real and repeated, or complex conjugates, respectively.

### 8.6 Exercises

1. Express the inverse transform as an integral.
(a) $\frac{1}{s^{2}\left(s^{2}+4\right)}$
(b) $\frac{s}{(s+2)\left(s^{2}+9\right)}$
(c) $\frac{s}{\left(s^{2}+4\right)\left(s^{2}+9\right)}$
(d) $\frac{s}{\left(s^{2}+1\right)^{2}}$
(e) $\frac{1}{s(s-a)}$
(f) $\frac{1}{(s+1)\left(s^{2}+2 s+2\right)}$
(g) $\frac{1}{(s+1)^{2}\left(s^{2}+4 s+5\right)}$
(h) $\frac{1}{(s-1)^{3}(s+2)^{2}}$
(i) $\frac{s-1}{s^{2}\left(s^{2}-2 s+2\right)}$
(j) $\frac{s(s+3)}{\left(s^{2}+4\right)\left(s^{2}+6 s+10\right)}$
(k) $\frac{1}{(s-3)^{5} s^{6}}$
(l) $\frac{1}{(s-1)^{3}\left(s^{2}+4\right)}$
(m) $\frac{1}{s^{2}(s-2)^{3}}$
(n) $\frac{1}{s^{7}(s-2)^{6}}$
2. Find the Laplace transform.
(a) $\int_{0}^{t} \sin a \tau \cos b(t-\tau) d \tau$
(b) $\int_{0}^{t} e^{\tau} \sin a(t-\tau) d \tau$
(c) $\int_{0}^{t} \sinh a \tau \cosh a(t-\tau) d \tau$
(d) $\int_{0}^{t} \tau(t-\tau) \sin \omega \tau \cos \omega(t-\tau) d \tau$
(e) $e^{t} \int_{0}^{t} \sin \omega \tau \cos \omega(t-\tau) d \tau$
(f) $e^{t} \int_{0}^{t} \tau^{2}(t-\tau) e^{\tau} d \tau$
(g) $e^{-t} \int_{0}^{t} e^{-\tau} \tau \cos \omega(t-\tau) d \tau$
(h) $e^{t} \int_{0}^{t} e^{2 \tau} \sinh (t-\tau) d \tau$
(i) $\int_{0}^{t} \tau e^{2 \tau} \sin 2(t-\tau) d \tau$
(j) $\int_{0}^{t}(t-\tau)^{3} e^{\tau} d \tau$
(k) $\int_{0}^{t} \tau^{6} e^{-(t-\tau)} \sin 3(t-\tau) d \tau$
(l) $\int_{0}^{t} \tau^{2}(t-\tau)^{3} d \tau$
(m) $\int_{0}^{t}(t-\tau)^{7} e^{-\tau} \sin 2 \tau d \tau$
(n) $\int_{0}^{t}(t-\tau)^{4} \sin 2 \tau d \tau$
3. Find a formula for the solution of the initial value problem.
(a) $y^{\prime \prime}+3 y^{\prime}+y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$
(b) $y^{\prime \prime}+4 y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$
(c) $y^{\prime \prime}+2 y^{\prime}+y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0$
(d) $y^{\prime \prime}+k^{2} y=f(t), \quad y(0)=1, \quad y^{\prime}(0)=-1$
(e) $y^{\prime \prime}+6 y^{\prime}+9 y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=-2$
(f) $y^{\prime \prime}-4 y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=3$
(g) $y^{\prime \prime}-5 y^{\prime}+6 y=f(t), \quad y(0)=1, \quad y^{\prime}(0)=3$
(h) $y^{\prime \prime}+\omega^{2} y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}$
4. Solve the integral equation.
(a) $y(t)=t-\int_{0}^{t}(t-\tau) y(\tau) d \tau$
(b) $y(t)=\sin t-2 \int_{0}^{t} \cos (t-\tau) y(\tau) d \tau$
$\begin{array}{ll}\text { (c) } y(t)=1+2 \int_{0}^{t} y(\tau) \cos (t-\tau) d \tau & \text { (d) } y(t)=t+\int_{0}^{t} y(\tau) e^{-(t-\tau)} d \tau\end{array}$
(e) $y^{\prime}(t)=t+\int_{0}^{t} y(\tau) \cos (t-\tau) d \tau, y(0)=4$
(f) $y(t)=\cos t-\sin t+\int_{0}^{t} y(\tau) \sin (t-\tau) d \tau$
5. Use the convolution theorem to evaluate the integral.
(a) $\int_{0}^{t}(t-\tau)^{7} \tau^{8} d \tau$
(b) $\int_{0}^{t}(t-\tau)^{13} \tau^{7} d \tau$
(c) $\int_{0}^{t}(t-\tau)^{6} \tau^{7} d \tau$
(d) $\int_{0}^{t} e^{-\tau} \sin (t-\tau) d \tau$
(e) $\int_{0}^{t} \sin \tau \cos 2(t-\tau) d \tau$
6. Show that

$$
\int_{0}^{t} f(t-\tau) g(\tau) d \tau=\int_{0}^{t} f(\tau) g(t-\tau) d \tau
$$

by introducing the new variable of integration $x=t-\tau$ in the first integral.
7. Use the convolution theorem to show that if $f(t) \leftrightarrow F(s)$ then

$$
\int_{0}^{t} f(\tau) d \tau \leftrightarrow \frac{F(s)}{s}
$$

8. Show that if $p(s)=a s^{2}+b s+c$ has distinct real zeros $r_{1}$ and $r_{2}$ then the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

is

$$
\begin{aligned}
y(t)= & k_{0} \frac{r_{2} e^{r_{1} t}-r_{1} e^{r_{2} t}}{r_{2}-r_{1}}+k_{1} \frac{e^{r_{2} t}-e^{r_{1} t}}{r_{2}-r_{1}} \\
& +\frac{1}{a\left(r_{2}-r_{1}\right)} \int_{0}^{t}\left(e^{r_{2} \tau}-e^{r_{1} \tau}\right) f(t-\tau) d \tau
\end{aligned}
$$

9. Show that if $p(s)=a s^{2}+b s+c$ has a repeated real zero $r_{1}$ then the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

is

$$
y(t)=k_{0}\left(1-r_{1} t\right) e^{r_{1} t}+k_{1} t e^{r_{1} t}+\frac{1}{a} \int_{0}^{t} \tau e^{r_{1} \tau} f(t-\tau) d \tau
$$

10. Show that if $p(s)=a s^{2}+b s+c$ has complex conjugate zeros $\lambda \pm i \omega$ then the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

is

$$
\begin{aligned}
y(t)= & e^{\lambda t}\left[k_{0}\left(\cos \omega t-\frac{\lambda}{\omega} \sin \omega t\right)+\frac{k_{1}}{\omega} \sin \omega t\right] \\
& +\frac{1}{a \omega} \int_{0}^{t} e^{\lambda t} f(t-\tau) \sin \omega \tau d \tau
\end{aligned}
$$

11. Let

$$
w=\mathcal{L}^{-1}\left(\frac{1}{a s^{2}+b s+c}\right)
$$

where $a, b$, and $c$ are constants and $a \neq 0$.
(a) Show that $w$ is the solution of

$$
a w^{\prime \prime}+b w^{\prime}+c w=0, \quad w(0)=0, \quad w^{\prime}(0)=\frac{1}{a}
$$

(b) Let $f$ be continuous on $[0, \infty)$ and define

$$
h(t)=\int_{0}^{t} w(t-\tau) f(\tau) d \tau
$$

Use Leibniz's rule for differentiating an integral with respect to a parameter to show that $h$ is the solution of

$$
a h^{\prime \prime}+b h^{\prime}+c h=f, \quad h(0)=0, \quad h^{\prime}(0)=0 .
$$

(c) Show that the function $y$ in Eqn. (8.6.14) is the solution of Eqn. (8.6.13) provided that $f$ is continuous on $[0, \infty)$; thus, it's not necessary to assume that $f$ has a Laplace transform.
12. Consider the initial value problem

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0, \tag{A}
\end{equation*}
$$

where $a, b$, and $c$ are constants, $a \neq 0$, and

$$
f(t)=\left\{\begin{array}{cc}
f_{0}(t), \quad 0 \leq t<t_{1} \\
f_{1}(t), & t \geq t_{1}
\end{array}\right.
$$

Assume that $f_{0}$ is continuous and of exponential order on $[0, \infty)$ and $f_{1}$ is continuous and of exponential order on $\left[t_{1}, \infty\right)$. Let

$$
p(s)=a s^{2}+b s+c .
$$

(a) Show that the Laplace transform of the solution of (A) is

$$
Y(s)=\frac{F_{0}(s)+e^{-s t_{1}} G(s)}{p(s)}
$$

where $g(t)=f_{1}\left(t+t_{1}\right)-f_{0}\left(t+t_{1}\right)$.
(b) Let $w$ be as in Exercise 11. Use Theorem 8.4.2 and the convolution theorem to show that the solution of (A) is

$$
y(t)=\int_{0}^{t} w(t-\tau) f_{0}(\tau) d \tau+u\left(t-t_{1}\right) \int_{0}^{t-t_{1}} w\left(t-t_{1}-\tau\right) g(\tau) d \tau
$$

for $t>0$.
(c) Henceforth, assume only that $f_{0}$ is continuous on $[0, \infty)$ and $f_{1}$ is continuous on $\left[t_{1}, \infty\right)$. Use Exercise 11 (a) and (b) to show that

$$
y^{\prime}(t)=\int_{0}^{t} w^{\prime}(t-\tau) f_{0}(\tau) d \tau+u\left(t-t_{1}\right) \int_{0}^{t-t_{1}} w^{\prime}\left(t-t_{1}-\tau\right) g(\tau) d \tau
$$

for $t>0$, and

$$
y^{\prime \prime}(t)=\frac{f(t)}{a}+\int_{0}^{t} w^{\prime \prime}(t-\tau) f_{0}(\tau) d \tau+u\left(t-t_{1}\right) \int_{0}^{t-t_{1}} w^{\prime \prime}\left(t-t_{1}-\tau\right) g(\tau) d \tau
$$

for $0<t<t_{1}$ and $t>t_{1}$. Also, show $y$ satisfies the differential equation in (A) on $\left(0, t_{1}\right)$ and $\left(t_{1}, \infty\right)$.
(d) Show that $y$ and $y^{\prime}$ are continuous on $[0, \infty)$.
13. Suppose

$$
f(t)=\left\{\begin{array}{cl}
f_{0}(t), & 0 \leq t<t_{1} \\
f_{1}(t), & t_{1} \leq t<t_{2} \\
& \vdots \\
f_{k-1}(t), & t_{k-1} \leq t<t_{k} \\
f_{k}(t), & t \geq t_{k}
\end{array}\right.
$$

where $f_{m}$ is continuous on $\left[t_{m}, \infty\right)$ for $m=0, \ldots, k$ (let $t_{0}=0$ ), and define

$$
g_{m}(t)=f_{m}\left(t+t_{m}\right)-f_{m-1}\left(t+t_{m}\right), m=1, \ldots, k
$$

Extend the results of Exercise 12 to show that the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0
$$

is

$$
y(t)=\int_{0}^{t} w(t-\tau) f_{0}(\tau) d \tau+\sum_{m=1}^{k} u\left(t-t_{m}\right) \int_{0}^{t-t_{m}} w\left(t-t_{m}-\tau\right) g_{m}(\tau) d \tau
$$

14. Let $\left\{t_{m}\right\}_{m=0}^{\infty}$ be a sequence of points such that $t_{0}=0, t_{m+1}>t_{m}$, and $\lim _{m \rightarrow \infty} t_{m}=\infty$. For each nonegative integer $m$ let $f_{m}$ be continuous on $\left[t_{m}, \infty\right)$, and let $f$ be defined on $[0, \infty)$ by

$$
f(t)=f_{m}(t), \quad t_{m} \leq t<t_{m+1} \quad m=0,1,2 \ldots
$$

Let

$$
g_{m}(t)=f_{m}\left(t+t_{m}\right)-f_{m-1}\left(t+t_{m}\right), \quad m=1, \ldots, k
$$

Extend the results of Exercise 13 to show that the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=0, \quad y^{\prime}(0)=0
$$

is

$$
y(t)=\int_{0}^{t} w(t-\tau) f_{0}(\tau) d \tau+\sum_{m=1}^{\infty} u\left(t-t_{m}\right) \int_{0}^{t-t_{m}} w\left(t-t_{m}-\tau\right) g_{m}(\tau) d \tau
$$

Hint: See Exercise30.

### 8.7 CONSTANT COEFFICIENT EQUATIONS WITH IMPULSES

So far in this chapter, we've considered initial value problems for the constant coefficient equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t)
$$

where $f$ is continuous or piecewise continuous on $[0, \infty)$. In this section we consider initial value problems where $f$ represents a force that's very large for a short time and zero otherwise. We say that such forces are impulsive. Impulsive forces occur, for example, when two objects collide. Since it isn't feasible to represent such forces as continuous or piecewise continuous functions, we must construct a different mathematical model to deal with them.

If $f$ is an integrable function and $f(t)=0$ for $t$ outside of the interval $\left[t_{0}, t_{0}+h\right]$, then $\int_{t_{0}}^{t_{0}+h} f(t) d t$ is called the total impulse of $f$. We're interested in the idealized situation where $h$ is so small that the total impulse can be assumed to be applied instantaneously at $t=t_{0}$. We say in this case that $f$ is an impulse function. In particular, we denote by $\delta\left(t-t_{0}\right)$ the impulse function with total impulse equal to one, applied at $t=t_{0}$. (The impulse function $\delta(t)$ obtained by setting $t_{0}=0$ is the Dirac $\delta$ function.) It must be understood, however, that $\delta\left(t-t_{0}\right)$ isn't a function in the standard sense, since our "definition" implies that $\delta\left(t-t_{0}\right)=0$ if $t \neq t_{0}$, while

$$
\int_{t_{0}}^{t_{0}} \delta\left(t-t_{0}\right) d t=1
$$



Figure 8.7.1 $y=f_{h}(t)$

From calculus we know that no function can have these properties; nevertheless, there's a branch of mathematics known as the theory of distributions where the definition can be made rigorous. Since the theory of distributions is beyond the scope of this book, we'll take an intuitive approach to impulse functions.

Our first task is to define what we mean by the solution of the initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=\delta\left(t-t_{0}\right), \quad y(0)=0, \quad y^{\prime}(0)=0
$$

where $t_{0}$ is a fixed nonnegative number. The next theorem will motivate our definition.
Theorem 8.7.1 Suppose $t_{0} \geq 0$. For each positive number $h$, let $y_{h}$ be the solution of the initial value problem

$$
\begin{equation*}
a y_{h}^{\prime \prime}+b y_{h}^{\prime}+c y_{h}=f_{h}(t), \quad y_{h}(0)=0, \quad y_{h}^{\prime}(0)=0 \tag{8.7.1}
\end{equation*}
$$

where

$$
f_{h}(t)=\left\{\begin{array}{cl}
0, & 0 \leq t<t_{0}  \tag{8.7.2}\\
1 / h, & t_{0} \leq t<t_{0}+h \\
0, & t \geq t_{0}+h
\end{array}\right.
$$

so $f_{h}$ has unit total impulse equal to the area of the shaded rectangle in Figure 8.7.1. Then

$$
\begin{equation*}
\lim _{h \rightarrow 0+} y_{h}(t)=u\left(t-t_{0}\right) w\left(t-t_{0}\right) \tag{8.7.3}
\end{equation*}
$$

where

$$
w=\mathcal{L}^{-1}\left(\frac{1}{a s^{2}+b s+c}\right)
$$

Proof Taking Laplace transforms in (8.7.1) yields

$$
\left(a s^{2}+b s+c\right) Y_{h}(s)=F_{h}(s)
$$

so

$$
Y_{h}(s)=\frac{F_{h}(s)}{a s^{2}+b s+c}
$$

The convolution theorem implies that

$$
y_{h}(t)=\int_{0}^{t} w(t-\tau) f_{h}(\tau) d \tau
$$

Therefore, (8.7.2) implies that

$$
y_{h}(t)=\left\{\begin{array}{cl}
0, & 0 \leq t<t_{0}  \tag{8.7.4}\\
\frac{1}{h} \int_{t_{0}}^{t} w(t-\tau) d \tau, & t_{0} \leq t \leq t_{0}+h \\
\frac{1}{h} \int_{t_{0}}^{t_{0}+h} w(t-\tau) d \tau, & t>t_{0}+h
\end{array}\right.
$$

Since $y_{h}(t)=0$ for all $h$ if $0 \leq t \leq t_{0}$, it follows that

$$
\begin{equation*}
\lim _{h \rightarrow 0+} y_{h}(t)=0 \quad \text { if } \quad 0 \leq t \leq t_{0} \tag{8.7.5}
\end{equation*}
$$

We'll now show that

$$
\begin{equation*}
\lim _{h \rightarrow 0+} y_{h}(t)=w\left(t-t_{0}\right) \quad \text { if } \quad t>t_{0} . \tag{8.7.6}
\end{equation*}
$$

Suppose $t$ is fixed and $t>t_{0}$. From (8.7.4),

$$
\begin{equation*}
y_{h}(t)=\frac{1}{h} \int_{t_{0}}^{t_{0}+h} w(t-\tau) d \tau \quad \text { if } \quad h<t-t_{0} \tag{8.7.7}
\end{equation*}
$$

Since

$$
\begin{equation*}
\frac{1}{h} \int_{t_{0}}^{t_{0}+h} d \tau=1 \tag{8.7.8}
\end{equation*}
$$

we can write

$$
w\left(t-t_{0}\right)=\frac{1}{h} w\left(t-t_{0}\right) \int_{t_{0}}^{t_{0}+h} d \tau=\frac{1}{h} \int_{t_{0}}^{t_{0}+h} w\left(t-t_{0}\right) d \tau
$$

From this and (8.7.7),

$$
y_{h}(t)-w\left(t-t_{0}\right)=\frac{1}{h} \int_{t_{0}}^{t_{0}+h}\left(w(t-\tau)-w\left(t-t_{0}\right)\right) d \tau
$$

Therefore

$$
\begin{equation*}
\left|y_{h}(t)-w\left(t-t_{0}\right)\right| \leq \frac{1}{h} \int_{t_{0}}^{t_{0}+h}\left|w(t-\tau)-w\left(t-t_{0}\right)\right| d \tau \tag{8.7.9}
\end{equation*}
$$

Now let $M_{h}$ be the maximum value of $\left|w(t-\tau)-w\left(t-t_{0}\right)\right|$ as $\tau$ varies over the interval $\left[t_{0}, t_{0}+h\right]$. (Remember that $t$ and $t_{0}$ are fixed.) Then (8.7.8) and (8.7.9) imply that

$$
\begin{equation*}
\left|y_{h}(t)-w\left(t-t_{0}\right)\right| \leq \frac{1}{h} M_{h} \int_{t_{0}}^{t_{0}+h} d \tau=M_{h} \tag{8.7.10}
\end{equation*}
$$

But $\lim _{h \rightarrow 0+} M_{h}=0$, since $w$ is continuous. Therefore (8.7.10) implies (8.7.6). This and (8.7.5) imply (8.7.3).

Theorem 8.7.1 motivates the next definition.
Definition 8.7.2 If $t_{0}>0$, then the solution of the initial value problem

$$
\begin{equation*}
a y^{\prime \prime}+b y^{\prime}+c y=\delta\left(t-t_{0}\right), \quad y(0)=0, \quad y^{\prime}(0)=0 \tag{8.7.11}
\end{equation*}
$$

is defined to be

$$
y=u\left(t-t_{0}\right) w\left(t-t_{0}\right),
$$

where

$$
w=\mathcal{L}^{-1}\left(\frac{1}{a s^{2}+b s+c}\right)
$$

In physical applications where the input $f$ and the output $y$ of a device are related by the differential equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t)
$$

$w$ is called the impulse response of the device. Note that $w$ is the solution of the initial value problem

$$
\begin{equation*}
a w^{\prime \prime}+b w^{\prime}+c w=0, \quad w(0)=0, \quad w^{\prime}(0)=1 / a \tag{8.7.12}
\end{equation*}
$$

as can be seen by using the Laplace transform to solve this problem. (Verify.) On the other hand, we can solve (8.7.12) by the methods of Section 5.2 and show that $w$ is defined on $(-\infty, \infty)$ by

$$
\begin{equation*}
w=\frac{e^{r_{2} t}-e^{r_{1} t}}{a\left(r_{2}-r_{1}\right)}, \quad w=\frac{1}{a} t e^{r_{1} t}, \quad \text { or } \quad w=\frac{1}{a \omega} e^{\lambda t} \sin \omega t, \tag{8.7.13}
\end{equation*}
$$

depending upon whether the polynomial $p(r)=a r^{2}+b r+c$ has distinct real zeros $r_{1}$ and $r_{2}$, a repeated zero $r_{1}$, or complex conjugate zeros $\lambda \pm i \omega$. (In most physical applications, the zeros of the characteristic polynomial have negative real parts, so $\lim _{t \rightarrow \infty} w(t)=0$.) This means that $y=u\left(t-t_{0}\right) w\left(t-t_{0}\right)$ is defined on $(-\infty, \infty)$ and has the following properties:

$$
\begin{gathered}
y(t)=0, \quad t<t_{0} \\
a y^{\prime \prime}+b y^{\prime}+c y=0 \quad \text { on } \quad\left(-\infty, t_{0}\right) \quad \text { and } \quad\left(t_{0}, \infty\right)
\end{gathered}
$$

and

$$
\begin{equation*}
y_{-}^{\prime}\left(t_{0}\right)=0, \quad y_{+}^{\prime}\left(t_{0}\right)=1 / a \tag{8.7.14}
\end{equation*}
$$

(remember that $y_{-}^{\prime}\left(t_{0}\right)$ and $y_{+}^{\prime}\left(t_{0}\right)$ are derivatives from the right and left, respectively) and $y^{\prime}\left(t_{0}\right)$ does not exist. Thus, even though we defined $y=u\left(t-t_{0}\right) w\left(t-t_{0}\right)$ to be the solution of (8.7.11), this function doesn't satisfy the differential equation in (8.7.11) at $t_{0}$, since it isn't differentiable there; in fact (8.7.14) indicates that an impulse causes a jump discontinuity in velocity. (To see that this is reasonable, think of what happens when you hit a ball with a bat.) This means that the initial value problem (8.7.11) doesn't make sense if $t_{0}=0$, since $y^{\prime}(0)$ doesn't exist in this case. However $y=u(t) w(t)$ can be defined to be the solution of the modified initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=\delta(t), \quad y(0)=0, \quad y_{-}^{\prime}(0)=0
$$

where the condition on the derivative at $t=0$ has been replaced by a condition on the derivative from the left.

Figure 8.7.2 illustrates Theorem 8.7.1 for the case where the impulse response $w$ is the first expression in (8.7.13) and $r_{1}$ and $r_{2}$ are distinct and both negative. The solid curve in the figure is the graph of $w$. The dashed curves are solutions of (8.7.1) for various values of $h$. As $h$ decreases the graph of $y_{h}$ moves to the left toward the graph of $w$.

Example 8.7.1 Find the solution of the initial value problem

$$
\begin{equation*}
y^{\prime \prime}-2 y^{\prime}+y=\delta\left(t-t_{0}\right), \quad y(0)=0, \quad y^{\prime}(0)=0 \tag{8.7.15}
\end{equation*}
$$

where $t_{0}>0$. Then interpret the solution for the case where $t_{0}=0$.

Solution Here

$$
w=\mathcal{L}^{-1}\left(\frac{1}{s^{2}-2 s+1}\right)=\mathcal{L}^{-1}\left(\frac{1}{(s-1)^{2}}\right)=t e^{-t}
$$

so Definition 8.7.2 yields

$$
y=u\left(t-t_{0}\right)\left(t-t_{0}\right) e^{-\left(t-t_{0}\right)}
$$

as the solution of (8.7.15) if $t_{0}>0$. If $t_{0}=0$, then (8.7.15) doesn't have a solution; however, $y=$ $u(t) t e^{-t}$ (which we would usually write simply as $y=t e^{-t}$ ) is the solution of the modified initial value problem

$$
y^{\prime \prime}-2 y^{\prime}+y=\delta(t), \quad y(0)=0, \quad y_{-}^{\prime}(0)=0
$$

The graph of $y=u\left(t-t_{0}\right)\left(t-t_{0}\right) e^{-\left(t-t_{0}\right)}$ is shown in Figure 8.7.3
Definition 8.7.2 and the principle of superposition motivate the next definition.


Figure 8.7.2 An illustration of Theorem 8.7.1


Figure 8.7.3 $y=u\left(t-t_{0}\right)\left(t-t_{0}\right) e^{-\left(t-t_{0}\right)}$

Definition 8.7.3 Suppose $\alpha$ is a nonzero constant and $f$ is piecewise continuous on $[0, \infty)$. If $t_{0}>0$, then the solution of the initial value problem

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t)+\alpha \delta\left(t-t_{0}\right), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1}
$$

is defined to be

$$
y(t)=\hat{y}(t)+\alpha u\left(t-t_{0}\right) w\left(t-t_{0}\right),
$$

where $\hat{y}$ is the solution of

$$
a y^{\prime \prime}+b y^{\prime}+c y=f(t), \quad y(0)=k_{0}, \quad y^{\prime}(0)=k_{1} .
$$

This definition also applies if $t_{0}=0$, provided that the initial condition $y^{\prime}(0)=k_{1}$ is replaced by $y_{-}^{\prime}(0)=k_{1}$.

Example 8.7.2 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+6 y^{\prime}+5 y=3 e^{-2 t}+2 \delta(t-1), \quad y(0)=-3, \quad y^{\prime}(0)=2 . \tag{8.7.16}
\end{equation*}
$$

Solution We leave it to you to show that the solution of

$$
y^{\prime \prime}+6 y^{\prime}+5 y=3 e^{-2 t}, \quad y(0)=-3, y^{\prime}(0)=2
$$

is

$$
\hat{y}=-e^{-2 t}+\frac{1}{2} e^{-5 t}-\frac{5}{2} e^{-t}
$$

Since

$$
\begin{aligned}
w(t) & =\mathcal{L}^{-1}\left(\frac{1}{s^{2}+6 s+5}\right)=\mathcal{L}^{-1}\left(\frac{1}{(s+1)(s+5)}\right) \\
& =\frac{1}{4} \mathcal{L}^{-1}\left(\frac{1}{s+1}-\frac{1}{s+5}\right)=\frac{e^{-t}-e^{-5 t}}{4}
\end{aligned}
$$

the solution of (8.7.16) is

$$
\begin{equation*}
y=-e^{-2 t}+\frac{1}{2} e^{-5 t}-\frac{5}{2} e^{-t}+u(t-1) \frac{e^{-(t-1)}-e^{-5(t-1)}}{2} \tag{8.7.17}
\end{equation*}
$$

(Figure 8.7.4)


Figure 8.7.4 Graph of (8.7.17)


Figure 8.7.5 Graph of (8.7.19)

Definition 8.7.3 can be extended in the obvious way to cover the case where the forcing function contains more than one impulse.

Example 8.7.3 Solve the initial value problem

$$
\begin{equation*}
y^{\prime \prime}+y=1+2 \delta(t-\pi)-3 \delta(t-2 \pi), \quad y(0)=-1, y^{\prime}(0)=2 \tag{8.7.18}
\end{equation*}
$$

Solution We leave it to you to show that

$$
\hat{y}=1-2 \cos t+2 \sin t
$$

is the solution of

$$
y^{\prime \prime}+y=1, \quad y(0)=-1, \quad y^{\prime}(0)=2
$$

Since

$$
w=\mathcal{L}^{-1}\left(\frac{1}{s^{2}+1}\right)=\sin t
$$

the solution of (8.7.18) is

$$
\begin{aligned}
y & =1-2 \cos t+2 \sin t+2 u(t-\pi) \sin (t-\pi)-3 u(t-2 \pi) \sin (t-2 \pi) \\
& =1-2 \cos t+2 \sin t-2 u(t-\pi) \sin t-3 u(t-2 \pi) \sin t
\end{aligned}
$$

or

$$
y=\left\{\begin{array}{cl}
1-2 \cos t+2 \sin t, & 0 \leq t<\pi  \tag{8.7.19}\\
1-2 \cos t, & \pi \leq t<2 \pi \\
1-2 \cos t-3 \sin t, & t \geq 2 \pi
\end{array}\right.
$$

(Figure 8.7.5).

### 8.7 Exercises

In Exercises 1-20 solve the initial value problem. Where indicated by $\mathrm{C} / \mathrm{G}$, graph the solution.

1. $y^{\prime \prime}+3 y^{\prime}+2 y=6 e^{2 t}+2 \delta(t-1), \quad y(0)=2, \quad y^{\prime}(0)=-6$
2. $\mathrm{C} / \mathrm{G} \quad y^{\prime \prime}+y^{\prime}-2 y=-10 e^{-t}+5 \delta(t-1), \quad y(0)=7, \quad y^{\prime}(0)=-9$
3. $y^{\prime \prime}-4 y=2 e^{-t}+5 \delta(t-1), \quad y(0)=-1, \quad y^{\prime}(0)=2$
4. $\mathrm{C} / \mathrm{G} \quad y^{\prime \prime}+y=\sin 3 t+2 \delta(t-\pi / 2), \quad y(0)=1, \quad y^{\prime}(0)=-1$
5. $y^{\prime \prime}+4 y=4+\delta(t-3 \pi), \quad y(0)=0, \quad y^{\prime}(0)=1$
6. $\quad y^{\prime \prime}-y=8+2 \delta(t-2), \quad y(0)=-1, \quad y^{\prime}(0)=1$
7. $y^{\prime \prime}+y^{\prime}=e^{t}+3 \delta(t-6), \quad y(0)=-1, \quad y^{\prime}(0)=4$
8. $y^{\prime \prime}+4 y=8 e^{2 t}+\delta(t-\pi / 2), \quad y(0)=8, \quad y^{\prime}(0)=0$
9. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+3 y^{\prime}+2 y=1+\delta(t-1), \quad y(0)=1, \quad y^{\prime}(0)=-1$
10. $y^{\prime \prime}+2 y^{\prime}+y=e^{t}+2 \delta(t-2), \quad y(0)=-1, \quad y^{\prime}(0)=2$
11. $\mathrm{C} / \mathrm{G} y^{\prime \prime}+4 y=\sin t+\delta(t-\pi / 2), \quad y(0)=0, \quad y^{\prime}(0)=2$
12. $y^{\prime \prime}+2 y^{\prime}+2 y=\delta(t-\pi)-3 \delta(t-2 \pi), \quad y(0)=-1, \quad y^{\prime}(0)=2$
13. $y^{\prime \prime}+4 y^{\prime}+13 y=\delta(t-\pi / 6)+2 \delta(t-\pi / 3), \quad y(0)=1, \quad y^{\prime}(0)=2$
14. $2 y^{\prime \prime}-3 y^{\prime}-2 y=1+\delta(t-2), \quad y(0)=-1, \quad y^{\prime}(0)=2$
15. $4 y^{\prime \prime}-4 y^{\prime}+5 y=4 \sin t-4 \cos t+\delta(t-\pi / 2)-\delta(t-\pi), \quad y(0)=1, \quad y^{\prime}(0)=1$
16. $y^{\prime \prime}+y=\cos 2 t+2 \delta(t-\pi / 2)-3 \delta(t-\pi), \quad y(0)=0, \quad y^{\prime}(0)=-1$
17. $\mathrm{C} / \mathrm{G} y^{\prime \prime}-y=4 e^{-t}-5 \delta(t-1)+3 \delta(t-2), \quad y(0)=0, \quad y^{\prime}(0)=0$
18. $y^{\prime \prime}+2 y^{\prime}+y=e^{t}-\delta(t-1)+2 \delta(t-2), \quad y(0)=0, \quad y^{\prime}(0)=-1$
19. $y^{\prime \prime}+y=f(t)+\delta(t-2 \pi), \quad y(0)=0, \quad y^{\prime}(0)=1$, and $f(t)=\left\{\begin{array}{cl}\sin 2 t, & 0 \leq t<\pi, \\ 0, & t \geq \pi .\end{array}\right.$
20. $y^{\prime \prime}+4 y=f(t)+\delta(t-\pi)-3 \delta(t-3 \pi / 2), \quad y(0)=1, \quad y^{\prime}(0)=-1$, and $f(t)= \begin{cases}1, & 0 \leq t<\pi / 2, \\ 2, & t \geq \pi / 2\end{cases}$
21. $\quad y^{\prime \prime}+y=\delta(t), \quad y(0)=1, \quad y_{-}^{\prime}(0)=-2$
22. $\quad y^{\prime \prime}-4 y=3 \delta(t), \quad y(0)=-1, \quad y_{-}^{\prime}(0)=7$
23. $y^{\prime \prime}+3 y^{\prime}+2 y=-5 \delta(t), \quad y(0)=0, \quad y_{-}^{\prime}(0)=0$
24. $y^{\prime \prime}+4 y^{\prime}+4 y=-\delta(t), \quad y(0)=1, \quad y_{-}^{\prime}(0)=5$
25. $4 y^{\prime \prime}+4 y^{\prime}+y=3 \delta(t), \quad y(0)=1, \quad y_{-}^{\prime}(0)=-6$

In Exercises 26-28, solve the initial value problem

$$
a y_{h}^{\prime \prime}+b y_{h}^{\prime}+c y_{h}=\left\{\begin{array}{cl}
0, & 0 \leq t<t_{0} \\
1 / h, & t_{0} \leq t<t_{0}+h, \\
0, & t \geq t_{0}+h
\end{array} \quad y_{h}(0)=0, \quad y_{h}^{\prime}(0)=0,\right.
$$

where $t_{0}>0$ and $h>0$. Then find

$$
w=\mathcal{L}^{-1}\left(\frac{1}{a s^{2}+b s+c}\right)
$$

and verify Theorem 8.7.1 by graphing $w$ and $y_{h}$ on the same axes, for small positive values of $h$.
26. L $y^{\prime \prime}+2 y^{\prime}+2 y=f_{h}(t), \quad y(0)=0, \quad y^{\prime}(0)=0$
27. L $y^{\prime \prime}+2 y^{\prime}+y=f_{h}(t), \quad y(0)=0, \quad y^{\prime}(0)=0$
28. $\mathrm{L} y^{\prime \prime}+3 y^{\prime}+2 y=f_{h}(t), \quad y(0)=0, \quad y^{\prime}(0)=0$
29. Recall from Section 6.2 that the displacement of an object of mass $m$ in a spring-mass system in free damped oscillation is

$$
m y^{\prime \prime}+c y^{\prime}+k y=0, \quad y(0)=y_{0}, \quad y^{\prime}(0)=v_{0},
$$

and that $y$ can be written as

$$
y=R e^{-c t / 2 m} \cos \left(\omega_{1} t-\phi\right)
$$

if the motion is underdamped. Suppose $y(\tau)=0$. Find the impulse that would have to be applied to the object at $t=\tau$ to put it in equilibrium.
30. Solve the initial value problem. Find a formula that does not involve step functions and represents $y$ on each subinterval of $[0, \infty)$ on which the forcing function is zero.
(a) $y^{\prime \prime}-y=\sum_{k=1}^{\infty} \delta(t-k), \quad y(0)=0, \quad y^{\prime}(0)=1$
(b) $y^{\prime \prime}+y=\sum_{k=1}^{\infty} \delta(t-2 k \pi), \quad y(0)=0, \quad y^{\prime}(0)=1$
(c) $y^{\prime \prime}-3 y^{\prime}+2 y=\sum_{k=1}^{\infty} \delta(t-k), \quad y(0)=0, \quad y^{\prime}(0)=1$
(d) $y^{\prime \prime}+y=\sum_{k=1}^{\infty} \delta(t-k \pi), \quad y(0)=0, \quad y^{\prime}(0)=0$

### 8.8 A BRIEF TABLE OF LAPLACE TRANSFORMS

| $f(t)$ | $F(s)$ |  |
| :---: | :---: | :---: |
| 1 | $\frac{1}{s}$ | $(s>0)$ |
| $\begin{aligned} & t^{n} \\ & (n=\text { integer }>0) \end{aligned}$ | $\frac{n!}{s^{n+1}}$ | $(s>0)$ |
| $t^{p}, p>-1$ | $\frac{\Gamma(p+1)}{s^{(p+1)}}$ | $(s>0)$ |
| $e^{a t}$ | $\frac{1}{s-a}$ | $(s>a)$ |
| $t^{n} e^{a t}$ | $\frac{n!}{(s-a)^{n+1}}$ | $(s>0)$ |
| $\begin{aligned} & (n=\text { integer }>0) \\ & \cos \omega t \end{aligned}$ | $\frac{s}{s^{2}+\omega^{2}}$ | $(s>0)$ |
| $\sin \omega t$ | $\frac{\omega}{s^{2}+\omega^{2}}$ | $(s>0)$ |
| $e^{\lambda t} \cos \omega t$ | $\frac{s-\lambda}{(s-\lambda)^{2}+\omega^{2}}$ | $(s>\lambda)$ |
| $e^{\lambda t} \sin \omega t$ | $\frac{\omega}{(s-\lambda)^{2}+\omega^{2}}$ | $(s>\lambda)$ |
| $\cosh b t$ | $\frac{s}{s^{2}-b^{2}}$ | $(s>\|b\|)$ |
| $\sinh b t$ | $\frac{b}{s^{2}-b^{2}}$ | $(s>\|b\|)$ |
| $t \cos \omega t$ | $\frac{s^{2}-\omega^{2}}{\left(s^{2}+\omega^{2}\right)^{2}}$ | $(s>0)$ |


| $t \sin \omega t$ | $\frac{2 \omega s}{\left(s^{2}+\omega^{2}\right)^{2}}$ | $(s>0)$ |
| :---: | :---: | :---: |
| $\sin \omega t-\omega t \cos \omega t$ | $\frac{2 \omega^{3}}{\left(s^{2}+\omega^{2}\right)^{2}}$ | $(s>0)$ |
| $\omega t-\sin \omega t$ | $\frac{\omega^{3}}{s^{2}\left(s^{2}+\omega^{2}\right)^{2}}$ | $(s>0)$ |
| $\frac{1}{t} \sin \omega t$ | $\arctan \left(\frac{\omega}{s}\right)$ | $(s>0)$ |
| $e^{a t} f(t)$ | $F(s-a)$ |  |
| $t^{k} f(t)$ | $(-1)^{k} F^{(k)}(s)$ |  |
| $f(\omega t)$ | $\frac{1}{\omega} F\left(\frac{s}{\omega}\right), \quad \omega>0$ |  |
| $u(t-\tau)$ | $\frac{e^{-\tau s}}{s}$ | $(s>0)$ |
| $u(t-\tau) f(t-\tau)(\tau>0)$ | $e^{-\tau s} F(s)$ |  |
| $\int_{o}^{t} f(\tau) g(t-\tau) d \tau$ | $F(s) \cdot G(s)$ |  |
| $\delta(t-a)$ | $e^{-a s}$ | $(s>0)$ |

## CHAPTER 10 Linear Systems of Differential Equations

IN THIS CHAPTER we consider systems of differential equations involving more than one unknown function. Such systems arise in many physical applications.

SECTION 10.1 presents examples of physical situations that lead to systems of differential equations.
SECTION 10.2 discusses linear systems of differential equations.
SECTION 10.3 deals with the basic theory of homogeneous linear systems.
SECTIONS 10.4, 10.5, AND 10.6 present the theory of constant coefficient homogeneous systems.
SECTION 10.7 presents the method of variation of parameters for nonhomogeneous linear systems.

### 10.1 INTRODUCTION TO SYSTEMS OF DIFFERENTIAL EQUATIONS

Many physical situations are modelled by systems of $n$ differential equations in $n$ unknown functions, where $n \geq 2$. The next three examples illustrate physical problems that lead to systems of differential equations. In these examples and throughout this chapter we'll denote the independent variable by $t$.

Example 10.1.1 Tanks $T_{1}$ and $T_{2}$ contain 100 gallons and 300 gallons of salt solutions, respectively. Salt solutions are simultaneously added to both tanks from external sources, pumped from each tank to the other, and drained from both tanks (Figure 10.1.1). A solution with 1 pound of salt per gallon is pumped into $T_{1}$ from an external source at $5 \mathrm{gal} / \mathrm{min}$, and a solution with 2 pounds of salt per gallon is pumped into $T_{2}$ from an external source at $4 \mathrm{gal} / \mathrm{min}$. The solution from $T_{1}$ is pumped into $T_{2}$ at $2 \mathrm{gal} / \mathrm{min}$, and the solution from $T_{2}$ is pumped into $T_{1}$ at $3 \mathrm{gal} / \mathrm{min}$. $T_{1}$ is drained at $6 \mathrm{gal} / \mathrm{min}$ and $T_{2}$ is drained at $3 \mathrm{gal} / \mathrm{min}$. Let $Q_{1}(t)$ and $Q_{2}(t)$ be the number of pounds of salt in $T_{1}$ and $T_{2}$, respectively, at time $t>0$. Derive a system of differential equations for $Q_{1}$ and $Q_{2}$. Assume that both mixtures are well stirred.


Figure 10.1.1

Solution As in Section 4.2, let rate in and rate out denote the rates ( $\mathrm{lb} / \mathrm{min}$ ) at which salt enters and leaves a tank; thus,

$$
\begin{aligned}
Q_{1}^{\prime} & =(\text { rate in })_{1}-(\text { rate out })_{1} \\
Q_{2}^{\prime} & =(\text { rate in })_{2}-(\text { rate out })_{2}
\end{aligned}
$$

Note that the volumes of the solutions in $T_{1}$ and $T_{2}$ remain constant at 100 gallons and 300 gallons, respectively.
$T_{1}$ receives salt from the external source at the rate of

$$
(1 \mathrm{lb} / \mathrm{gal}) \times(5 \mathrm{gal} / \mathrm{min})=5 \mathrm{lb} / \mathrm{min},
$$

and from $T_{2}$ at the rate of

$$
\left(\mathrm{lb} / \mathrm{gal} \text { in } T_{2}\right) \times(3 \mathrm{gal} / \mathrm{min})=\frac{1}{300} Q_{2} \times 3=\frac{1}{100} Q_{2} \mathrm{lb} / \mathrm{min} .
$$

Therefore

$$
\begin{equation*}
(\text { rate in })_{1}=5+\frac{1}{100} Q_{2} \tag{10.1.1}
\end{equation*}
$$

Solution leaves $T_{1}$ at the rate of $8 \mathrm{gal} / \mathrm{min}$, since $6 \mathrm{gal} / \mathrm{min}$ are drained and $2 \mathrm{gal} / \mathrm{min}$ are pumped to $T_{2}$; hence,

$$
\begin{equation*}
(\text { rate out })_{1}=\left(\mathrm{lb} / \mathrm{gal} \text { in } \mathrm{T}_{1}\right) \times(8 \mathrm{gal} / \mathrm{min})=\frac{1}{100} Q_{1} \times 8=\frac{2}{25} Q_{1} \tag{10.1.2}
\end{equation*}
$$

Eqns. (10.1.1) and (10.1.2) imply that

$$
\begin{equation*}
Q_{1}^{\prime}=5+\frac{1}{100} Q_{2}-\frac{2}{25} Q_{1} . \tag{10.1.3}
\end{equation*}
$$

$T_{2}$ receives salt from the external source at the rate of

$$
(2 \mathrm{lb} / \mathrm{gal}) \times(4 \mathrm{gal} / \mathrm{min})=8 \mathrm{lb} / \mathrm{min},
$$

and from $T_{1}$ at the rate of

$$
\left(\mathrm{lb} / \mathrm{gal} \text { in } T_{1}\right) \times(2 \mathrm{gal} / \mathrm{min})=\frac{1}{100} Q_{1} \times 2=\frac{1}{50} Q_{1} \mathrm{lb} / \mathrm{min} .
$$

Therefore

$$
\begin{equation*}
(\text { rate in })_{2}=8+\frac{1}{50} Q_{1} \tag{10.1.4}
\end{equation*}
$$

Solution leaves $T_{2}$ at the rate of $6 \mathrm{gal} / \mathrm{min}$, since $3 \mathrm{gal} / \mathrm{min}$ are drained and $3 \mathrm{gal} / \mathrm{min}$ are pumped to $T_{1}$; hence,

$$
\begin{equation*}
(\text { rate out })_{2}=\left(\mathrm{lb} / \mathrm{gal} \text { in } \mathrm{T}_{2}\right) \times(6 \mathrm{gal} / \mathrm{min})=\frac{1}{300} Q_{2} \times 6=\frac{1}{50} Q_{2} . \tag{10.1.5}
\end{equation*}
$$

Eqns. (10.1.4) and (10.1.5) imply that

$$
\begin{equation*}
Q_{2}^{\prime}=8+\frac{1}{50} Q_{1}-\frac{1}{50} Q_{2} \tag{10.1.6}
\end{equation*}
$$

We say that (10.1.3) and (10.1.6) form a system of two first order equations in two unknowns, and write them together as

$$
\begin{aligned}
Q_{1}^{\prime} & =5-\frac{2}{25} Q_{1}+\frac{1}{100} Q_{2} \\
Q_{2}^{\prime} & =8+\frac{1}{50} Q_{1}-\frac{1}{50} Q_{2}
\end{aligned}
$$

Example 10.1.2 A mass $m_{1}$ is suspended from a rigid support on a spring $S_{1}$ and a second mass $m_{2}$ is suspended from the first on a spring $S_{2}$ (Figure 10.1.2). The springs obey Hooke's law, with spring constants $k_{1}$ and $k_{2}$. Internal friction causes the springs to exert damping forces proportional to the rates of change of their lengths, with damping constants $c_{1}$ and $c_{2}$. Let $y_{1}=y_{1}(t)$ and $y_{2}=y_{2}(t)$ be the displacements of the two masses from their equilibrium positions at time $t$, measured positive upward. Derive a system of differential equations for $y_{1}$ and $y_{2}$, assuming that the masses of the springs are negligible and that vertical external forces $F_{1}$ and $F_{2}$ also act on the objects.

Solution In equilibrium, $S_{1}$ supports both $m_{1}$ and $m_{2}$ and $S_{2}$ supports only $m_{2}$. Therefore, if $\Delta \ell_{1}$ and $\Delta \ell_{2}$ are the elongations of the springs in equilibrium then

$$
\begin{equation*}
\left(m_{1}+m_{2}\right) g=k_{1} \Delta \ell_{1} \quad \text { and } \quad m_{2} g=k_{2} \Delta \ell_{2} . \tag{10.1.7}
\end{equation*}
$$

Let $H_{1}$ be the Hooke's law force acting on $m_{1}$, and let $D_{1}$ be the damping force on $m_{1}$. Similarly, let $H_{2}$ and $D_{2}$ be the Hooke's law and damping forces acting on $m_{2}$. According to Newton's second law of motion,

$$
\begin{align*}
& m_{1} y_{1}^{\prime \prime}=-m_{1} g+H_{1}+D_{1}+F_{1}, \\
& m_{2} y_{2}^{\prime \prime}=-m_{2} g+H_{2}+D_{2}+F_{2} \tag{10.1.8}
\end{align*}
$$

When the displacements are $y_{1}$ and $y_{2}$, the change in length of $S_{1}$ is $-y_{1}+\Delta \ell_{1}$ and the change in length of $S_{2}$ is $-y_{2}+y_{1}+\Delta \ell_{2}$. Both springs exert Hooke's law forces on $m_{1}$, while only $S_{2}$ exerts a Hooke's law force on $m_{2}$. These forces are in directions that tend to restore the springs to their natural lengths. Therefore

$$
\begin{equation*}
H_{1}=k_{1}\left(-y_{1}+\Delta \ell_{1}\right)-k_{2}\left(-y_{2}+y_{1}+\Delta \ell_{2}\right) \quad \text { and } \quad H_{2}=k_{2}\left(-y_{2}+y_{1}+\Delta \ell_{2}\right) \tag{10.1.9}
\end{equation*}
$$



Figure 10.1.2

When the velocities are $y_{1}^{\prime}$ and $y_{2}^{\prime}, S_{1}$ and $S_{2}$ are changing length at the rates $-y_{1}^{\prime}$ and $-y_{2}^{\prime}+y_{1}^{\prime}$, respectively. Both springs exert damping forces on $m_{1}$, while only $S_{2}$ exerts a damping force on $m_{2}$. Since the force due to damping exerted by a spring is proportional to the rate of change of length of the spring and in a direction that opposes the change, it follows that

$$
\begin{equation*}
D_{1}=-c_{1} y_{1}^{\prime}+c_{2}\left(y_{2}^{\prime}-y_{1}^{\prime}\right) \quad \text { and } \quad D_{2}=-c_{2}\left(y_{2}^{\prime}-y_{1}^{\prime}\right) \tag{10.1.10}
\end{equation*}
$$

From (10.1.8), (10.1.9), and (10.1.10),

$$
\begin{align*}
m_{1} y_{1}^{\prime \prime}= & -m_{1} g+k_{1}\left(-y_{1}+\Delta \ell_{1}\right)-k_{2}\left(-y_{2}+y_{1}+\Delta \ell_{2}\right) \\
& -c_{1} y_{1}^{\prime}+c_{2}\left(y_{2}^{\prime}-y_{1}^{\prime}\right)+F_{1} \\
= & -\left(m_{1} g-k_{1} \Delta \ell_{1}+k_{2} \Delta \ell_{2}\right)-k_{1} y_{1}+k_{2}\left(y_{2}-y_{1}\right)  \tag{10.1.11}\\
& -c_{1} y_{1}^{\prime}+c_{2}\left(y_{2}^{\prime}-y_{1}^{\prime}\right)+F_{1}
\end{align*}
$$

and

$$
\begin{align*}
m_{2} y_{2}^{\prime \prime} & =-m_{2} g+k_{2}\left(-y_{2}+y_{1}+\Delta \ell_{2}\right)-c_{2}\left(y_{2}^{\prime}-y_{1}^{\prime}\right)+F_{2}  \tag{10.1.12}\\
& =-\left(m_{2} g-k_{2} \Delta \ell_{2}\right)-k_{2}\left(y_{2}-y_{1}\right)-c_{2}\left(y_{2}^{\prime}-y_{1}^{\prime}\right)+F_{2}
\end{align*}
$$

From (10.1.7),

$$
m_{1} g-k_{1} \Delta \ell_{1}+k_{2} \Delta \ell_{2}=-m_{2} g+k_{2} \Delta \ell_{2}=0
$$

Therefore we can rewrite (10.1.11) and (10.1.12) as

$$
\begin{aligned}
& m_{1} y_{1}^{\prime \prime}=-\left(c_{1}+c_{2}\right) y_{1}^{\prime}+c_{2} y_{2}^{\prime}-\left(k_{1}+k_{2}\right) y_{1}+k_{2} y_{2}+F_{1} \\
& m_{2} y_{2}^{\prime \prime}=c_{2} y_{1}^{\prime}-c_{2} y_{2}^{\prime}+k_{2} y_{1}-k_{2} y_{2}+F_{2} .
\end{aligned}
$$

Example 10.1.3 Let $\mathbf{X}=\mathbf{X}(t)=x(t) \mathbf{i}+y(t) \mathbf{j}+z(t) \mathbf{k}$ be the position vector at time $t$ of an object with mass $m$, relative to a rectangular coordinate system with origin at Earth's center (Figure 10.1.3). According to Newton's law of gravitation, Earth's gravitational force $\mathbf{F}=\mathbf{F}(x, y, z)$ on the object is inversely proportional to the square of the distance of the object from Earth's center, and directed toward the center; thus,

$$
\begin{equation*}
\mathbf{F}=\frac{K}{\|\mathbf{X}\|^{2}}\left(-\frac{\mathbf{X}}{\|\mathbf{X}\|}\right)=-K \frac{x \mathbf{i}+y \mathbf{j}+z \mathbf{k}}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}} \tag{10.1.13}
\end{equation*}
$$

where $K$ is a constant. To determine $K$, we observe that the magnitude of $\mathbf{F}$ is

$$
\|\mathbf{F}\|=K \frac{\|\mathbf{X}\|}{\|\mathbf{X}\|^{3}}=\frac{K}{\|\mathbf{X}\|^{2}}=\frac{K}{\left(x^{2}+y^{2}+z^{2}\right)}
$$



Figure 10.1.3

Let $R$ be Earth's radius. Since $\|\mathbf{F}\|=m g$ when the object is at Earth's surface,

$$
m g=\frac{K}{R^{2}}, \quad \text { so } \quad K=m g R^{2}
$$

Therefore we can rewrite (10.1.13) as

$$
\mathbf{F}=-m g R^{2} \frac{x \mathbf{i}+y \mathbf{j}+z \mathbf{k}}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}
$$

Now suppose $\mathbf{F}$ is the only force acting on the object. According to Newton's second law of motion, $\mathbf{F}=m \mathbf{X}^{\prime \prime}$; that is,

$$
m\left(x^{\prime \prime} \mathbf{i}+y^{\prime \prime} \mathbf{j}+z^{\prime \prime} \mathbf{k}\right)=-m g R^{2} \frac{x \mathbf{i}+y \mathbf{j}+z \mathbf{k}}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}
$$

Cancelling the common factor $m$ and equating components on the two sides of this equation yields the system

$$
\begin{align*}
x^{\prime \prime} & =-\frac{g R^{2} x}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}} \\
y^{\prime \prime} & =-\frac{g R^{2} y}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}  \tag{10.1.14}\\
z^{\prime \prime} & =-\frac{g R^{2} z}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}} .
\end{align*}
$$

Rewriting Higher Order Systems as First Order Systems
A system of the form

$$
\begin{align*}
y_{1}^{\prime} & =g_{1}\left(t, y_{1}, y_{2}, \ldots, y_{n}\right) \\
y_{2}^{\prime} & =g_{2}\left(t, y_{1}, y_{2}, \ldots, y_{n}\right) \\
& \vdots  \tag{10.1.15}\\
y_{n}^{\prime} & =g_{n}\left(t, y_{1}, y_{2}, \ldots, y_{n}\right)
\end{align*}
$$

is called a first order system, since the only derivatives occurring in it are first derivatives. The derivative of each of the unknowns may depend upon the independent variable and all the unknowns, but not on the derivatives of other unknowns. When we wish to emphasize the number of unknown functions in (10.1.15) we will say that (10.1.15) is an $n \times n$ system.

Systems involving higher order derivatives can often be reformulated as first order systems by introducing additional unknowns. The next two examples illustrate this.

Example 10.1.4 Rewrite the system

$$
\begin{align*}
& m_{1} y_{1}^{\prime \prime}=-\left(c_{1}+c_{2}\right) y_{1}^{\prime}+c_{2} y_{2}^{\prime}-\left(k_{1}+k_{2}\right) y_{1}+k_{2} y_{2}+F_{1}  \tag{10.1.16}\\
& m_{2} y_{2}^{\prime \prime}=c_{2} y_{1}^{\prime}-c_{2} y_{2}^{\prime}+k_{2} y_{1}-k_{2} y_{2}+F_{2} .
\end{align*}
$$

derived in Example 10.1.2 as a system of first order equations.

Solution If we define $v_{1}=y_{1}^{\prime}$ and $v_{2}=y_{2}^{\prime}$, then $v_{1}^{\prime}=y_{1}^{\prime \prime}$ and $v_{2}^{\prime}=y_{2}^{\prime \prime}$, so (10.1.16) becomes

$$
\begin{aligned}
& m_{1} v_{1}^{\prime}=-\left(c_{1}+c_{2}\right) v_{1}+c_{2} v_{2}-\left(k_{1}+k_{2}\right) y_{1}+k_{2} y_{2}+F_{1} \\
& m_{2} v_{2}^{\prime}=c_{2} v_{1}-c_{2} v_{2}+k_{2} y_{1}-k_{2} y_{2}+F_{2} .
\end{aligned}
$$

Therefore $\left\{y_{1}, y_{2}, v_{1}, v_{2}\right\}$ satisfies the $4 \times 4$ first order system

$$
\begin{align*}
y_{1}^{\prime} & =v_{1} \\
y_{2}^{\prime} & =v_{2} \\
v_{1}^{\prime} & =\frac{1}{m_{1}}\left[-\left(c_{1}+c_{2}\right) v_{1}+c_{2} v_{2}-\left(k_{1}+k_{2}\right) y_{1}+k_{2} y_{2}+F_{1}\right]  \tag{10.1.17}\\
v_{2}^{\prime} & =\frac{1}{m_{2}}\left[c_{2} v_{1}-c_{2} v_{2}+k_{2} y_{1}-k_{2} y_{2}+F_{2}\right] .
\end{align*}
$$

REMARK: The difference in form between (10.1.15) and (10.1.17), due to the way in which the unknowns are denoted in the two systems, isn't important; (10.1.17) is a first order system, in that each equation in (10.1.17) expresses the first derivative of one of the unknown functions in a way that does not involve derivatives of any of the other unknowns.

Example 10.1.5 Rewrite the system

$$
\begin{aligned}
x^{\prime \prime} & =f\left(t, x, x^{\prime}, y, y^{\prime}, y^{\prime \prime}\right) \\
y^{\prime \prime \prime} & =g\left(t, x, x^{\prime}, y, y^{\prime} y^{\prime \prime}\right)
\end{aligned}
$$

as a first order system.

Solution We regard $x, x^{\prime}, y, y^{\prime}$, and $y^{\prime \prime}$ as unknown functions, and rename them

$$
x=x_{1}, x^{\prime}=x_{2}, \quad y=y_{1}, \quad y^{\prime}=y_{2}, \quad y^{\prime \prime}=y_{3} .
$$

These unknowns satisfy the system

$$
\begin{aligned}
x_{1}^{\prime} & =x_{2} \\
x_{2}^{\prime} & =f\left(t, x_{1}, x_{2}, y_{1}, y_{2}, y_{3}\right) \\
y_{1}^{\prime} & =y_{2} \\
y_{2}^{\prime} & =y_{3} \\
y_{3}^{\prime} & =g\left(t, x_{1}, x_{2}, y_{1}, y_{2}, y_{3}\right) .
\end{aligned}
$$

Rewriting Scalar Differential Equations as Systems
In this chapter we'll refer to differential equations involving only one unknown function as scalar differential equations. Scalar differential equations can be rewritten as systems of first order equations by the method illustrated in the next two examples.

Example 10.1.6 Rewrite the equation

$$
\begin{equation*}
y^{(4)}+4 y^{\prime \prime \prime}+6 y^{\prime \prime}+4 y^{\prime}+y=0 \tag{10.1.18}
\end{equation*}
$$

as a $4 \times 4$ first order system.

Solution We regard $y, y^{\prime}, y^{\prime \prime}$, and $y^{\prime \prime \prime}$ as unknowns and rename them

$$
y=y_{1}, \quad y^{\prime}=y_{2}, \quad y^{\prime \prime}=y_{3}, \quad \text { and } \quad y^{\prime \prime \prime}=y_{4} .
$$

Then $y^{(4)}=y_{4}^{\prime}$, so (10.1.18) can be written as

$$
y_{4}^{\prime}+4 y_{4}+6 y_{3}+4 y_{2}+y_{1}=0 .
$$

Therefore $\left\{y_{1}, y_{2}, y_{3}, y_{4}\right\}$ satisfies the system

$$
\begin{aligned}
y_{1}^{\prime} & =y_{2} \\
y_{2}^{\prime} & =y_{3} \\
y_{3}^{\prime} & =y_{4} \\
y_{4}^{\prime} & =-4 y_{4}-6 y_{3}-4 y_{2}-y_{1}
\end{aligned}
$$

Example 10.1.7 Rewrite

$$
x^{\prime \prime \prime}=f\left(t, x, x^{\prime}, x^{\prime \prime}\right)
$$

as a system of first order equations.

Solution We regard $x, x^{\prime}$, and $x^{\prime \prime}$ as unknowns and rename them

$$
x=y_{1}, \quad x^{\prime}=y_{2}, \quad \text { and } \quad x^{\prime \prime}=y_{3} .
$$

Then

$$
y_{1}^{\prime}=x^{\prime}=y_{2}, \quad y_{2}^{\prime}=x^{\prime \prime}=y_{3}, \quad \text { and } \quad y_{3}^{\prime}=x^{\prime \prime \prime} .
$$

Therefore $\left\{y_{1}, y_{2}, y_{3}\right\}$ satisfies the first order system

$$
\begin{aligned}
y_{1}^{\prime} & =y_{2} \\
y_{2}^{\prime} & =y_{3} \\
y_{3}^{\prime} & =f\left(t, y_{1}, y_{2}, y_{3}\right)
\end{aligned}
$$

Since systems of differential equations involving higher derivatives can be rewritten as first order systems by the method used in Examples 10.1.5-10.1.7, we'll consider only first order systems.

Numerical Solution of Systems
The numerical methods that we studied in Chapter 3 can be extended to systems, and most differential equation software packages include programs to solve systems of equations. We won't go into detail on numerical methods for systems; however, for illustrative purposes we'll describe the Runge-Kutta method for the numerical solution of the initial value problem

$$
\begin{array}{ll}
y_{1}^{\prime}=g_{1}\left(t, y_{1}, y_{2}\right), & y_{1}\left(t_{0}\right)=y_{10} \\
y_{2}^{\prime}=g_{2}\left(t, y_{1}, y_{2}\right), & y_{2}\left(t_{0}\right)=y_{20}
\end{array}
$$

at equally spaced points $t_{0}, t_{1}, \ldots, t_{n}=b$ in an interval $\left[t_{0}, b\right]$. Thus,

$$
t_{i}=t_{0}+i h, \quad i=0,1, \ldots, n
$$

where

$$
h=\frac{b-t_{0}}{n} .
$$

We'll denote the approximate values of $y_{1}$ and $y_{2}$ at these points by $y_{10}, y_{11}, \ldots, y_{1 n}$ and $y_{20}, y_{21}, \ldots, y_{2 n}$.

The Runge-Kutta method computes these approximate values as follows: given $y_{1 i}$ and $y_{2 i}$, compute

$$
\begin{aligned}
I_{1 i} & =g_{1}\left(t_{i}, y_{1 i}, y_{2 i}\right) \\
J_{1 i} & =g_{2}\left(t_{i}, y_{1 i}, y_{2 i}\right), \\
I_{2 i} & =g_{1}\left(t_{i}+\frac{h}{2}, y_{1 i}+\frac{h}{2} I_{1 i}, y_{2 i}+\frac{h}{2} J_{1 i}\right), \\
J_{2 i} & =g_{2}\left(t_{i}+\frac{h}{2}, y_{1 i}+\frac{h}{2} I_{1 i}, y_{2 i}+\frac{h}{2} J_{1 i}\right), \\
I_{3 i} & =g_{1}\left(t_{i}+\frac{h}{2}, y_{1 i}+\frac{h}{2} I_{2 i}, y_{2 i}+\frac{h}{2} J_{2 i}\right), \\
J_{3 i} & =g_{2}\left(t_{i}+\frac{h}{2}, y_{1 i}+\frac{h}{2} I_{2 i}, y_{2 i}+\frac{h}{2} J_{2 i}\right), \\
I_{4 i} & =g_{1}\left(t_{i}+h, y_{1 i}+h I_{3 i}, y_{2 i}+h J_{3 i}\right), \\
J_{4 i} & =g_{2}\left(t_{i}+h, y_{1 i}+h I_{3 i}, y_{2 i}+h J_{3 i}\right),
\end{aligned}
$$

and

$$
\begin{aligned}
y_{1, i+1} & =y_{1 i}+\frac{h}{6}\left(I_{1 i}+2 I_{2 i}+2 I_{3 i}+I_{4 i}\right), \\
y_{2, i+1} & =y_{2 i}+\frac{h}{6}\left(J_{1 i}+2 J_{2 i}+2 J_{3 i}+J_{4 i}\right)
\end{aligned}
$$

for $i=0, \ldots, n-1$. Under appropriate conditions on $g_{1}$ and $g_{2}$, it can be shown that the global truncation error for the Runge-Kutta method is $O\left(h^{4}\right)$, as in the scalar case considered in Section 3.3.

### 10.1 Exercises

1. Tanks $T_{1}$ and $T_{2}$ contain 50 gallons and 100 gallons of salt solutions, respectively. A solution with 2 pounds of salt per gallon is pumped into $T_{1}$ from an external source at $1 \mathrm{gal} / \mathrm{min}$, and a solution with 3 pounds of salt per gallon is pumped into $T_{2}$ from an external source at $2 \mathrm{gal} / \mathrm{min}$. The solution from $T_{1}$ is pumped into $T_{2}$ at $3 \mathrm{gal} / \mathrm{min}$, and the solution from $T_{2}$ is pumped into $T_{1}$ at $4 \mathrm{gal} / \mathrm{min}$. $T_{1}$ is drained at $2 \mathrm{gal} / \mathrm{min}$ and $T_{2}$ is drained at $1 \mathrm{gal} / \mathrm{min}$. Let $Q_{1}(t)$ and $Q_{2}(t)$ be the number of pounds of salt in $T_{1}$ and $T_{2}$, respectively, at time $t>0$. Derive a system of differential equations for $Q_{1}$ and $Q_{2}$. Assume that both mixtures are well stirred.
2. Two 500 gallon tanks $T_{1}$ and $T_{2}$ initially contain 100 gallons each of salt solution. A solution with 2 pounds of salt per gallon is pumped into $T_{1}$ from an external source at $6 \mathrm{gal} / \mathrm{min}$, and a solution with 1 pound of salt per gallon is pumped into $T_{2}$ from an external source at $5 \mathrm{gal} / \mathrm{min}$. The solution from $T_{1}$ is pumped into $T_{2}$ at $2 \mathrm{gal} / \mathrm{min}$, and the solution from $T_{2}$ is pumped into $T_{1}$ at $1 \mathrm{gal} / \mathrm{min}$. Both tanks are drained at $3 \mathrm{gal} / \mathrm{min}$. Let $Q_{1}(t)$ and $Q_{2}(t)$ be the number of pounds of salt in $T_{1}$ and $T_{2}$, respectively, at time $t>0$. Derive a system of differential equations for $Q_{1}$ and $Q_{2}$ that's valid until a tank is about to overflow. Assume that both mixtures are well stirred.
3. A mass $m_{1}$ is suspended from a rigid support on a spring $S_{1}$ with spring constant $k_{1}$ and damping constant $c_{1}$. A second mass $m_{2}$ is suspended from the first on a spring $S_{2}$ with spring constant $k_{2}$ and damping constant $c_{2}$, and a third mass $m_{3}$ is suspended from the second on a spring $S_{3}$ with spring constant $k_{3}$ and damping constant $c_{3}$. Let $y_{1}=y_{1}(t), y_{2}=y_{2}(t)$, and $y_{3}=y_{3}(t)$ be the displacements of the three masses from their equilibrium positions at time $t$, measured positive upward. Derive a system of differential equations for $y_{1}, y_{2}$ and $y_{3}$, assuming that the masses of the springs are negligible and that vertical external forces $F_{1}, F_{2}$, and $F_{3}$ also act on the masses.
4. Let $\mathbf{X}=x \mathbf{i}+y \mathbf{j}+z \mathbf{k}$ be the position vector of an object with mass $m$, expressed in terms of a rectangular coordinate system with origin at Earth's center (Figure 10.1.3). Derive a system of differential equations for $x, y$, and $z$, assuming that the object moves under Earth's gravitational force (given by Newton's law of gravitation, as in Example 10.1.3 ) and a resistive force proportional to the speed of the object. Let $\alpha$ be the constant of proportionality.
5. Rewrite the given system as a first order system.
(a) $x^{\prime \prime \prime}=f\left(t, x, y, y^{\prime}\right)$ $y^{\prime \prime}=g\left(t, y, y^{\prime}\right)$

$$
u^{\prime}=f\left(t, u, v, v^{\prime}, w^{\prime}\right)
$$

(b) $v^{\prime \prime}=g\left(t, u, v, v^{\prime}, w\right)$
$w^{\prime \prime}=h\left(t, u, v, v^{\prime}, w, w^{\prime}\right)$
(c) $y^{\prime \prime \prime}=f\left(t, y, y^{\prime}, y^{\prime \prime}\right)$
(d) $y^{(4)}=f(t, y)$
(e) $\begin{aligned} x^{\prime \prime} & =f(t, x, y) \\ y^{\prime \prime} & =g(t, x, y)\end{aligned}$
6. Rewrite the system (10.1.14) of differential equations derived in Example 10.1.3 as a first order system.
7. Formulate a version of Euler's method (Section 3.1) for the numerical solution of the initial value problem

$$
\begin{array}{ll}
y_{1}^{\prime}=g_{1}\left(t, y_{1}, y_{2}\right), & y_{1}\left(t_{0}\right)=y_{10}, \\
y_{2}^{\prime}=g_{2}\left(t, y_{1}, y_{2}\right), & y_{2}\left(t_{0}\right)=y_{20},
\end{array}
$$

on an interval $\left[t_{0}, b\right]$.
8. Formulate a version of the improved Euler method (Section 3.2) for the numerical solution of the initial value problem

$$
\begin{array}{ll}
y_{1}^{\prime}=g_{1}\left(t, y_{1}, y_{2}\right), & y_{1}\left(t_{0}\right)=y_{10} \\
y_{2}^{\prime}=g_{2}\left(t, y_{1}, y_{2}\right), & y_{2}\left(t_{0}\right)=y_{20}
\end{array}
$$

on an interval $\left[t_{0}, b\right]$.

### 10.2 LINEAR SYSTEMS OF DIFFERENTIAL EQUATIONS

A first order system of differential equations that can be written in the form

$$
\begin{align*}
y_{1}^{\prime} & =a_{11}(t) y_{1}+a_{12}(t) y_{2}+\cdots+a_{1 n}(t) y_{n}+f_{1}(t) \\
y_{2}^{\prime} & =a_{21}(t) y_{1}+a_{22}(t) y_{2}+\cdots+a_{2 n}(t) y_{n}+f_{2}(t) \\
& \vdots  \tag{10.2.1}\\
y_{n}^{\prime} & =a_{n 1}(t) y_{1}+a_{n 2}(t) y_{2}+\cdots+a_{n n}(t) y_{n}+f_{n}(t)
\end{align*}
$$

is called a linear system.
The linear system (10.2.1) can be written in matrix form as

$$
\left[\begin{array}{c}
y_{1}^{\prime} \\
y_{2}^{\prime} \\
\vdots \\
y_{n}^{\prime}
\end{array}\right]=\left[\begin{array}{cccc}
a_{11}(t) & a_{12}(t) & \cdots & a_{1 n}(t) \\
a_{21}(t) & a_{22}(t) & \cdots & a_{2 n}(t) \\
\vdots & \vdots & \ddots & \vdots \\
a_{n 1}(t) & a_{n 2}(t) & \cdots & a_{n n}(t)
\end{array}\right]\left[\begin{array}{c}
y_{1} \\
y_{2} \\
\vdots \\
y_{n}
\end{array}\right]+\left[\begin{array}{c}
f_{1}(t) \\
f_{2}(t) \\
\vdots \\
f_{n}(t)
\end{array}\right]
$$

or more briefly as

$$
\begin{equation*}
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t) \tag{10.2.2}
\end{equation*}
$$

where

$$
\mathbf{y}=\left[\begin{array}{c}
y_{1} \\
y_{2} \\
\vdots \\
y_{n}
\end{array}\right], \quad A(t)=\left[\begin{array}{cccc}
a_{11}(t) & a_{12}(t) & \cdots & a_{1 n}(t) \\
a_{21}(t) & a_{22}(t) & \cdots & a_{2 n}(t) \\
\vdots & \vdots & \ddots & \vdots \\
a_{n 1}(t) & a_{n 2}(t) & \cdots & a_{n n}(t)
\end{array}\right], \quad \text { and } \quad \mathbf{f}(t)=\left[\begin{array}{c}
f_{1}(t) \\
f_{2}(t) \\
\vdots \\
f_{n}(t)
\end{array}\right]
$$

We call $A$ the coefficient matrix of (10.2.2) and $\mathbf{f}$ the forcing function. We'll say that $A$ and $\mathbf{f}$ are continuous if their entries are continuous. If $\mathbf{f}=\mathbf{0}$, then (10.2.2) is homogeneous; otherwise, (10.2.2) is nonhomogeneous.

An initial value problem for (10.2.2) consists of finding a solution of (10.2.2) that equals a given constant vector

$$
\mathbf{k}=\left[\begin{array}{c}
k_{1} \\
k_{2} \\
\vdots \\
k_{n}
\end{array}\right]
$$

at some initial point $t_{0}$. We write this initial value problem as

$$
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t), \quad \mathbf{y}\left(t_{0}\right)=\mathbf{k}
$$

The next theorem gives sufficient conditions for the existence of solutions of initial value problems for (10.2.2). We omit the proof.

Theorem 10.2.1 Suppose the coefficient matrix $A$ and the forcing function $\mathbf{f}$ are continuous on $(a, b)$, let $t_{0}$ be in $(a, b)$, and let $\mathbf{k}$ be an arbitrary constant $n$-vector. Then the initial value problem

$$
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t), \quad \mathbf{y}\left(t_{0}\right)=\mathbf{k}
$$

has a unique solution on $(a, b)$.

## Example 10.2.1

(a) Write the system

$$
\begin{align*}
& y_{1}^{\prime}=y_{1}+2 y_{2}+2 e^{4 t} \\
& y_{2}^{\prime}=2 y_{1}+y_{2}+e^{4 t} \tag{10.2.3}
\end{align*}
$$

in matrix form and conclude from Theorem 10.2.1 that every initial value problem for (10.2.3) has a unique solution on $(-\infty, \infty)$.
(b) Verify that

$$
\mathbf{y}=\frac{1}{5}\left[\begin{array}{l}
8  \tag{10.2.4}\\
7
\end{array}\right] e^{4 t}+c_{1}\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{3 t}+c_{2}\left[\begin{array}{l}
1 \\
-1
\end{array}\right] e^{-t}
$$

is a solution of (10.2.3) for all values of the constants $c_{1}$ and $c_{2}$.
(c) Find the solution of the initial value problem

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
1 & 2  \tag{10.2.5}\\
2 & 1
\end{array}\right] \mathbf{y}+\left[\begin{array}{l}
2 \\
1
\end{array}\right] e^{4 t}, \quad \mathbf{y}(0)=\frac{1}{5}\left[\begin{array}{l}
3 \\
22
\end{array}\right]
$$

$\underline{\text { Solution(a) }}$ The system (10.2.3) can be written in matrix form as

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
1 & 2 \\
2 & 1
\end{array}\right] \mathbf{y}+\left[\begin{array}{l}
2 \\
1
\end{array}\right] e^{4 t} .
$$

An initial value problem for (10.2.3) can be written as

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
1 & 2 \\
2 & 1
\end{array}\right] \mathbf{y}+\left[\begin{array}{l}
2 \\
1
\end{array}\right] e^{4 t}, \quad y\left(t_{0}\right)=\left[\begin{array}{l}
k_{1} \\
k_{2}
\end{array}\right]
$$

Since the coefficient matrix and the forcing function are both continuous on $(-\infty, \infty)$, Theorem 10.2.1 implies that this problem has a unique solution on $(-\infty, \infty)$.
$\underline{\text { SOLUTION(b) If } y \text { is given by (10.2.4), then }}$

$$
\begin{aligned}
A \mathbf{y}+\mathbf{f}= & \frac{1}{5}\left[\begin{array}{ll}
1 & 2 \\
2 & 1
\end{array}\right]\left[\begin{array}{l}
8 \\
7
\end{array}\right] e^{4 t}+c_{1}\left[\begin{array}{ll}
1 & 2 \\
2 & 1
\end{array}\right]\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{3 t} \\
& +c_{2}\left[\begin{array}{ll}
1 & 2 \\
2 & 1
\end{array}\right]\left[\begin{array}{l}
1 \\
-1
\end{array}\right] e^{-t}+\left[\begin{array}{l}
2 \\
1
\end{array}\right] e^{4 t} \\
= & \frac{1}{5}\left[\begin{array}{l}
22 \\
23
\end{array}\right] e^{4 t}+c_{1}\left[\begin{array}{l}
3 \\
3
\end{array}\right] e^{3 t}+c_{2}\left[\begin{array}{l}
-1 \\
1
\end{array}\right] e^{-t}+\left[\begin{array}{l}
2 \\
1
\end{array}\right] e^{4 t} \\
= & \frac{1}{5}\left[\begin{array}{l}
32 \\
28
\end{array}\right] e^{4 t}+3 c_{1}\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{3 t}-c_{2}\left[\begin{array}{l}
1 \\
-1
\end{array}\right] e^{-t}=\mathbf{y}^{\prime} .
\end{aligned}
$$

$\underline{\text { SOLUTION(c) }}$ We must choose $c_{1}$ and $c_{2}$ in (10.2.4) so that

$$
\frac{1}{5}\left[\begin{array}{l}
8 \\
7
\end{array}\right]+c_{1}\left[\begin{array}{l}
1 \\
1
\end{array}\right]+c_{2}\left[\begin{array}{l}
1 \\
-1
\end{array}\right]=\frac{1}{5}\left[\begin{array}{l}
3 \\
22
\end{array}\right]
$$

which is equivalent to

$$
\left[\begin{array}{rr}
1 & 1 \\
1 & -1
\end{array}\right]\left[\begin{array}{l}
c_{1} \\
c_{2}
\end{array}\right]=\left[\begin{array}{l}
-1 \\
3
\end{array}\right]
$$

Solving this system yields $c_{1}=1, c_{2}=-2$, so

$$
\mathbf{y}=\frac{1}{5}\left[\begin{array}{l}
8 \\
7
\end{array}\right] e^{4 t}+\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{3 t}-2\left[\begin{array}{l}
1 \\
-1
\end{array}\right] e^{-t}
$$

is the solution of (10.2.5).
REMARK: The theory of $n \times n$ linear systems of differential equations is analogous to the theory of the scalar $n$-th order equation

$$
\begin{equation*}
P_{0}(t) y^{(n)}+P_{1}(t) y^{(n-1)}+\cdots+P_{n}(t) y=F(t) \tag{10.2.6}
\end{equation*}
$$

as developed in Sections 9.1. For example, by rewriting (10.2.6) as an equivalent linear system it can be shown that Theorem 10.2.1 implies Theorem ?? (Exercise 12).

### 10.2 Exercises

1. Rewrite the system in matrix form and verify that the given vector function satisfies the system for any choice of the constants $c_{1}$ and $c_{2}$.
(a) $\begin{aligned} & y_{1}^{\prime}=2 y_{1}+4 y_{2} \\ & y_{2}^{\prime}=4 y_{1}+2 y_{2} ;\end{aligned} \quad \mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{6 t}+c_{2}\left[\begin{array}{l}1 \\ -1\end{array}\right] e^{-2 t}$
(b) $\begin{aligned} & y_{1}^{\prime}=-2 y_{1}-2 y_{2} \\ & y_{2}^{\prime}=-5 y_{1}+y_{2} ;\end{aligned} \quad \mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{-4 t}+c_{2}\left[\begin{array}{l}-2 \\ 5\end{array}\right] e^{3 t}$
(c) $\begin{aligned} & y_{1}^{\prime}=-4 y_{1}-10 y_{2} \\ & y_{2}^{\prime}=3 y_{1}+7 y_{2} ;\end{aligned} \quad \mathbf{y}=c_{1}\left[\begin{array}{l}-5 \\ 3\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{l}2 \\ -1\end{array}\right] e^{t}$
(d) $\begin{aligned} & y_{1}^{\prime}=2 y_{1}+y_{2} \\ & y_{2}^{\prime}=y_{1}+2 y_{2} ;\end{aligned} \quad \mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{3 t}+c_{2}\left[\begin{array}{l}1 \\ -1\end{array}\right] e^{t}$
2. Rewrite the system in matrix form and verify that the given vector function satisfies the system for any choice of the constants $c_{1}, c_{2}$, and $c_{3}$.
(a) $\begin{array}{rrr}y_{1}^{\prime} & = & -y_{1}+2 y_{2}+3 y_{3} \\ y_{2}^{\prime} & = & y_{2}+6 y_{3} \\ y_{3}^{\prime} & = & -2 y_{3} ;\end{array}$
$\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1 \\ 0\end{array}\right] e^{t}+c_{2}\left[\begin{array}{l}1 \\ 0 \\ 0\end{array}\right] e^{-t}+c_{3}\left[\begin{array}{r}1 \\ -2 \\ 1\end{array}\right] e^{-2 t}$
(b) $\begin{array}{rlr}y_{1}^{\prime} & = & 2 y_{2}+2 y_{3} \\ y_{2}^{\prime} & =2 y_{1} & +2 y_{3}\end{array}$
$y_{3}^{\prime}=2 y_{1}+2 y_{2}$;
$\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ 0 \\ 1\end{array}\right] e^{-2 t}+c_{2}\left[\begin{array}{r}0 \\ -1 \\ 1\end{array}\right] e^{-2 t}+c_{3}\left[\begin{array}{l}1 \\ 1 \\ 1\end{array}\right] e^{4 t}$
(c) $\begin{aligned} y_{1}^{\prime} & =-y_{1}+2 y_{2}+2 y_{3} \\ y_{2}^{\prime} & =2 y_{1}-y_{2}+2 y_{3} \\ y_{3}^{\prime} & =2 y_{1}+2 y_{2}-y_{3}\end{aligned}$.
$y_{3}^{\prime}=2 y_{1}+2 y_{2}-y_{3}$;
$\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ 0 \\ 1\end{array}\right] e^{-3 t}+c_{2}\left[\begin{array}{r}0 \\ -1 \\ 1\end{array}\right] e^{-3 t}+c_{3}\left[\begin{array}{l}1 \\ 1 \\ 1\end{array}\right] e^{3 t}$
(d) $\begin{aligned} y_{1}^{\prime} & =3 y_{1}-y_{2}-y_{3} \\ y_{2}^{\prime} & =-2 y_{1}+3 y_{2}+2 y_{3} \\ y_{3}^{\prime} & =4 y_{1}-y_{2}-2 y_{3} ;\end{aligned}$

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{3 t}+c_{3}\left[\begin{array}{r}
1 \\
-3 \\
7
\end{array}\right] e^{-t}
$$

3. Rewrite the initial value problem in matrix form and verify that the given vector function is a solution.

$$
\begin{array}{ll}
\text { (a) } \quad y_{1}^{\prime}=y_{1}+y_{2} \quad y_{1}(0)=1 \\
y_{2}^{\prime}=-2 y_{1}+4 y_{2}, \quad y_{2}(0)=0 ; & \mathbf{y}=2\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{2 t}-\left[\begin{array}{l}
1 \\
2
\end{array}\right] e^{3 t} \\
\text { (b) } \quad y_{1}^{\prime}=5 y_{1}+3 y_{2} \quad y_{1}(0)=12 \\
y_{2}^{\prime}=-y_{1}+y_{2}, \quad y_{2}(0)=-6 ; \quad \mathbf{y}=3\left[\begin{array}{l}
1 \\
-1
\end{array}\right] e^{2 t}+3\left[\begin{array}{l}
3 \\
-1
\end{array}\right] e^{4 t}
\end{array}
$$

4. Rewrite the initial value problem in matrix form and verify that the given vector function is a solution.
(a) $\begin{array}{rlll}y_{1}^{\prime} & =6 y_{1}+4 y_{2}+4 y_{3} & y_{1}(0) & =3 \\ y_{2}^{\prime} & =-7 y_{1}-2 y_{2}-y_{3},, & y_{2}(0) & =-6 \\ y_{3}^{\prime} & =7 y_{1}+4 y_{2}+3 y_{3} & y_{3}(0) & =4\end{array}$

$$
\mathbf{y}=\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{6 t}+2\left[\begin{array}{r}
1 \\
-2 \\
1
\end{array}\right] e^{2 t}+\left[\begin{array}{r}
0 \\
-1 \\
1
\end{array}\right] e^{-t}
$$

$$
y_{1}^{\prime}=8 y_{1}+7 y_{2}+7 y_{3} \quad y_{1}(0)=2
$$

(b) $\begin{array}{rll}y_{2}^{\prime}=-5 y_{1}-6 y_{2}-9 y_{3}, & y_{2}(0)=-4 \\ y_{3}^{\prime} & =5 y_{1}+7 y_{2}+10 y_{3}, & y_{3}(0)=\end{array}$

$$
\mathbf{y}=\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{8 t}+\left[\begin{array}{r}
0 \\
-1 \\
1
\end{array}\right] e^{3 t}+\left[\begin{array}{r}
1 \\
-2 \\
1
\end{array}\right] e^{t}
$$

5. Rewrite the system in matrix form and verify that the given vector function satisfies the system for any choice of the constants $c_{1}$ and $c_{2}$.
(a) $\begin{aligned} & y_{1}^{\prime}=-3 y_{1}+2 y_{2}+3-2 t \\ & y_{2}^{\prime}=-5 y_{1}+3 y_{2}+6-3 t\end{aligned}$

$$
\mathbf{y}=c_{1}\left[\begin{array}{c}
2 \cos t \\
3 \cos t-\sin t
\end{array}\right]+c_{2}\left[\begin{array}{c}
2 \sin t \\
3 \sin t+\cos t
\end{array}\right]+\left[\begin{array}{l}
1 \\
t
\end{array}\right]
$$

(b) $\begin{aligned} y_{1}^{\prime} & =3 y_{1}+y_{2}-5 e^{t} \\ y_{2}^{\prime} & =-y_{1}+y_{2}+e^{t}\end{aligned}$

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
-1 \\
1
\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{c}
1+t \\
-t
\end{array}\right] e^{2 t}+\left[\begin{array}{l}
1 \\
3
\end{array}\right] e^{t}
$$

(c) $\begin{aligned} & y_{1}^{\prime}=-y_{1}-4 y_{2}+4 e^{t}+8 t e^{t} \\ & y_{2}^{\prime}=-y_{1}-y_{2}+e^{3 t}+(4 t+2) e^{t}\end{aligned}$

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
2 \\
1
\end{array}\right] e^{-3 t}+c_{2}\left[\begin{array}{l}
-2 \\
1
\end{array}\right] e^{t}+\left[\begin{array}{c}
e^{3 t} \\
2 t e^{t}
\end{array}\right]
$$

(d) $\begin{aligned} & y_{1}^{\prime}=-6 y_{1}-3 y_{2}+14 e^{2 t}+12 e^{t} \\ & y_{2}^{\prime}=y_{1}-2 y_{2}+7 e^{2 t}-12 e^{t}\end{aligned}$

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
-3 \\
1
\end{array}\right] e^{-5 t}+c_{2}\left[\begin{array}{l}
-1 \\
1
\end{array}\right] e^{-3 t}+\left[\begin{array}{c}
e^{2 t}+3 e^{t} \\
2 e^{2 t}-3 e^{t}
\end{array}\right]
$$

6. Convert the linear scalar equation

$$
\begin{equation*}
P_{0}(t) y^{(n)}+P_{1}(t) y^{(n-1)}+\cdots+P_{n}(t) y(t)=F(t) \tag{A}
\end{equation*}
$$

into an equivalent $n \times n$ system

$$
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t),
$$

and show that $A$ and $\mathbf{f}$ are continuous on an interval $(a, b)$ if and only if (A) is normal on $(a, b)$.
7. A matrix function

$$
Q(t)=\left[\begin{array}{cccc}
q_{11}(t) & q_{12}(t) & \cdots & q_{1 s}(t) \\
q_{21}(t) & q_{22}(t) & \cdots & q_{2 s}(t) \\
\vdots & \vdots & \ddots & \vdots \\
q_{r 1}(t) & q_{r 2}(t) & \cdots & q_{r s}(t)
\end{array}\right]
$$

is said to be differentiable if its entries $\left\{q_{i j}\right\}$ are differentiable. Then the derivative $Q^{\prime}$ is defined by

$$
Q^{\prime}(t)=\left[\begin{array}{cccc}
q_{11}^{\prime}(t) & q_{12}^{\prime}(t) & \cdots & q_{1 s}^{\prime}(t) \\
q_{21}^{\prime}(t) & q_{22}^{\prime}(t) & \cdots & q_{2 s}^{\prime}(t) \\
\vdots & \vdots & \ddots & \vdots \\
q_{r 1}^{\prime}(t) & q_{r 2}^{\prime}(t) & \cdots & q_{r s}^{\prime}(t)
\end{array}\right] .
$$

(a) Prove: If $P$ and $Q$ are differentiable matrices such that $P+Q$ is defined and if $c_{1}$ and $c_{2}$ are constants, then

$$
\left(c_{1} P+c_{2} Q\right)^{\prime}=c_{1} P^{\prime}+c_{2} Q^{\prime} .
$$

(b) Prove: If $P$ and $Q$ are differentiable matrices such that $P Q$ is defined, then

$$
(P Q)^{\prime}=P^{\prime} Q+P Q^{\prime}
$$

8. Verify that $Y^{\prime}=A Y$.
(a) $Y=\left[\begin{array}{rr}e^{6 t} & e^{-2 t} \\ e^{6 t} & -e^{-2 t}\end{array}\right], \quad A=\left[\begin{array}{ll}2 & 4 \\ 4 & 2\end{array}\right]$
(b) $Y=\left[\begin{array}{rr}e^{-4 t} & -2 e^{3 t} \\ e^{-4 t} & 5 e^{3 t}\end{array}\right], \quad A=\left[\begin{array}{rr}-2 & -2 \\ -5 & 1\end{array}\right]$
(c) $Y=\left[\begin{array}{rr}-5 e^{2 t} & 2 e^{t} \\ 3 e^{2 t} & -e^{t}\end{array}\right], \quad A=\left[\begin{array}{rr}-4 & -10 \\ 3 & 7\end{array}\right]$
(d) $Y=\left[\begin{array}{rr}e^{3 t} & e^{t} \\ e^{3 t} & -e^{t}\end{array}\right], \quad A=\left[\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right]$
(e) $Y=\left[\begin{array}{rrr}e^{t} & e^{-t} & e^{-2 t} \\ e^{t} & 0 & -2 e^{-2 t} \\ 0 & 0 & e^{-2 t}\end{array}\right], \quad A=\left[\begin{array}{rrr}-1 & 2 & 3 \\ 0 & 1 & 6 \\ 0 & 0 & -2\end{array}\right]$
(f) $Y=\left[\begin{array}{ccc}-e^{-2 t} & -e^{-2 t} & e^{4 t} \\ 0 & e^{-2 t} & e^{4 t} \\ e^{-2 t} & 0 & e^{4 t}\end{array}\right], \quad A=\left[\begin{array}{lll}0 & 2 & 2 \\ 2 & 0 & 2 \\ 2 & 2 & 0\end{array}\right]$
(g) $Y=\left[\begin{array}{ccc}e^{3 t} & e^{-3 t} & 0 \\ e^{3 t} & 0 & -e^{-3 t} \\ e^{3 t} & e^{-3 t} & e^{-3 t}\end{array}\right], \quad A=\left[\begin{array}{ccc}-9 & 6 & 6 \\ -6 & 3 & 6 \\ -6 & 6 & 3\end{array}\right]$
(h) $Y=\left[\begin{array}{crr}e^{2 t} & e^{3 t} & e^{-t} \\ 0 & -e^{3 t} & -3 e^{-t} \\ e^{2 t} & e^{3 t} & 7 e^{-t}\end{array}\right], \quad A=\left[\begin{array}{rrr}3 & -1 & -1 \\ -2 & 3 & 2 \\ 4 & -1 & -2\end{array}\right]$
9. Suppose

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
y_{11} \\
y_{21}
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=\left[\begin{array}{l}
y_{12} \\
y_{22}
\end{array}\right]
$$

are solutions of the homogeneous system

$$
\begin{equation*}
\mathbf{y}^{\prime}=A(t) \mathbf{y} \tag{A}
\end{equation*}
$$

and define

$$
Y=\left[\begin{array}{ll}
y_{11} & y_{12} \\
y_{21} & y_{22}
\end{array}\right]
$$

(a) Show that $Y^{\prime}=A Y$.
(b) Show that if $\mathbf{c}$ is a constant vector then $\mathbf{y}=Y \mathbf{c}$ is a solution of (A).
(c) State generalizations of (a) and (b) for $n \times n$ systems.
10. Suppose $Y$ is a differentiable square matrix.
(a) Find a formula for the derivative of $Y^{2}$.
(b) Find a formula for the derivative of $Y^{n}$, where $n$ is any positive integer.
(c) State how the results obtained in (a) and (b) are analogous to results from calculus concerning scalar functions.
11. It can be shown that if $Y$ is a differentiable and invertible square matrix function, then $Y^{-1}$ is differentiable.
(a) Show that $\left(Y^{-1}\right)^{\prime}=-Y^{-1} Y^{\prime} Y^{-1}$. (Hint: Differentiate the identity $Y^{-1} Y=I$.)
(b) Find the derivative of $Y^{-n}=\left(Y^{-1}\right)^{n}$, where $n$ is a positive integer.
(c) State how the results obtained in (a) and (b) are analogous to results from calculus concerning scalar functions.
12. Show that Theorem 10.2.1 implies Theorem ??. Hint: Write the scalar equation

$$
P_{0}(x) y^{(n)}+P_{1}(x) y^{(n-1)}+\cdots+P_{n}(x) y=F(x)
$$

as an $n \times n$ system of linear equations.
13. Suppose $\mathbf{y}$ is a solution of the $n \times n$ system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$, and that the $n \times n$ matrix $P$ is invertible and differentiable on $(a, b)$. Find a matrix $B$ such that the function $\mathbf{x}=P \mathbf{y}$ is a solution of $\mathbf{x}^{\prime}=B \mathbf{x}$ on $(a, b)$.

### 10.3 BASIC THEORY OF HOMOGENEOUS LINEAR SYSTEMS

In this section we consider homogeneous linear systems $\mathbf{y}^{\prime}=A(t) \mathbf{y}$, where $A=A(t)$ is a continuous $n \times n$ matrix function on an interval $(a, b)$. The theory of linear homogeneous systems has much in common with the theory of linear homogeneous scalar equations, which we considered in Sections 2.1, 5.1, and 9.1.

Whenever we refer to solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ we'll mean solutions on $(a, b)$. Since $\mathbf{y} \equiv \mathbf{0}$ is obviously a solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$, we call it the trivial solution. Any other solution is nontrivial.

If $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ are vector functions defined on an interval $(a, b)$ and $c_{1}, c_{2}, \ldots, c_{n}$ are constants, then

$$
\begin{equation*}
\mathbf{y}=c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}+\cdots+c_{n} \mathbf{y}_{n} \tag{10.3.1}
\end{equation*}
$$

is a linear combination of $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$. It's easy show that if $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ are solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$, then so is any linear combination of $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ (Exercise 1 ). We say that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is a fundamental set of solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$ on if every solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$ can be written as a linear combination of $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$, as in (10.3.1). In this case we say that (10.3.1) is the general solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$.

It can be shown that if $A$ is continuous on $(a, b)$ then $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ has infinitely many fundamental sets of solutions on $(a, b)$ (Exercises 15 and 16). The next definition will help to characterize fundamental sets of solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$.

We say that a set $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ of $n$-vector functions is linearly independent on $(a, b)$ if the only constants $c_{1}, c_{2}, \ldots, c_{n}$ such that

$$
\begin{equation*}
c_{1} \mathbf{y}_{1}(t)+c_{2} \mathbf{y}_{2}(t)+\cdots+c_{n} \mathbf{y}_{n}(t)=0, \quad a<t<b \tag{10.3.2}
\end{equation*}
$$

are $c_{1}=c_{2}=\cdots=c_{n}=0$. If (10.3.2) holds for some set of constants $c_{1}, c_{2}, \ldots, c_{n}$ that are not all zero, then $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is linearly dependent on $(a, b)$

The next theorem is analogous to Theorems 5.1.3 and ??.
Theorem 10.3.1 Suppose the $n \times n$ matrix $A=A(t)$ is continuous on $(a, b)$. Then a set $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ of $n$ solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$ is a fundamental set if and only if it's linearly independent on $(a, b)$.

Example 10.3.1 Show that the vector functions

$$
\mathbf{y}_{1}=\left[\begin{array}{c}
e^{t} \\
0 \\
e^{-t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
0 \\
e^{3 t} \\
1
\end{array}\right], \quad \text { and } \quad \mathbf{y}_{3}=\left[\begin{array}{c}
e^{2 t} \\
e^{3 t} \\
0
\end{array}\right]
$$

are linearly independent on every interval $(a, b)$.

## Solution Suppose

$$
c_{1}\left[\begin{array}{c}
e^{t} \\
0 \\
e^{-t}
\end{array}\right]+c_{2}\left[\begin{array}{c}
0 \\
e^{3 t} \\
1
\end{array}\right]+c_{3}\left[\begin{array}{c}
e^{2 t} \\
e^{3 t} \\
0
\end{array}\right]=\left[\begin{array}{l}
0 \\
0 \\
0
\end{array}\right], \quad a<t<b .
$$

We must show that $c_{1}=c_{2}=c_{3}=0$. Rewriting this equation in matrix form yields

$$
\left[\begin{array}{ccc}
e^{t} & 0 & e^{2 t} \\
0 & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right]\left[\begin{array}{l}
c_{1} \\
c_{2} \\
c_{3}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0 \\
0
\end{array}\right], \quad a<t<b
$$

Expanding the determinant of this system in cofactors of the entries of the first row yields

$$
\begin{aligned}
\left|\begin{array}{ccc}
e^{t} & 0 & e^{2 t} \\
0 & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right| & =e^{t}\left|\begin{array}{cc}
e^{3 t} & e^{3 t} \\
1 & 0
\end{array}\right|-0\left|\begin{array}{cc}
0 & e^{3 t} \\
e^{-t} & 0
\end{array}\right|+e^{2 t}\left|\begin{array}{cc}
0 & e^{3 t} \\
e^{-t} & 1
\end{array}\right| \\
& =e^{t}\left(-e^{3 t}\right)+e^{2 t}\left(-e^{2 t}\right)=-2 e^{4 t}
\end{aligned}
$$

Since this determinant is never zero, $c_{1}=c_{2}=c_{3}=0$.
We can use the method in Example 10.3.1 to test $n$ solutions $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ of any $n \times n$ system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ for linear independence on an interval $(a, b)$ on which $A$ is continuous. To explain this (and for other purposes later), it's useful to write a linear combination of $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ in a different way. We first write the vector functions in terms of their components as

$$
\mathbf{y}_{1}=\left[\begin{array}{c}
y_{11} \\
y_{21} \\
\vdots \\
y_{n 1}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
y_{12} \\
y_{22} \\
\vdots \\
y_{n 2}
\end{array}\right], \ldots, \quad \mathbf{y}_{n}=\left[\begin{array}{c}
y_{1 n} \\
y_{2 n} \\
\vdots \\
y_{n n}
\end{array}\right]
$$

If

$$
\mathbf{y}=c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}+\cdots+c_{n} \mathbf{y}_{n}
$$

then

$$
\begin{aligned}
\mathbf{y} & =c_{1}\left[\begin{array}{c}
y_{11} \\
y_{21} \\
\vdots \\
y_{n 1}
\end{array}\right]+c_{2}\left[\begin{array}{c}
y_{12} \\
y_{22} \\
\vdots \\
y_{n 2}
\end{array}\right]+\cdots+c_{n}\left[\begin{array}{c}
y_{1 n} \\
y_{2 n} \\
\vdots \\
y_{n n}
\end{array}\right] \\
& =\left[\begin{array}{cccc}
y_{11} & y_{12} & \cdots & y_{1 n} \\
y_{21} & y_{22} & \cdots & y_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
y_{n 1} & y_{n 2} & \cdots & y_{n n}
\end{array}\right]\left[\begin{array}{c}
c_{1} \\
c_{2} \\
\vdots \\
c_{n}
\end{array}\right]
\end{aligned}
$$

This shows that

$$
\begin{equation*}
c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}+\cdots+c_{n} \mathbf{y}_{n}=Y \mathbf{c} \tag{10.3.3}
\end{equation*}
$$

where

$$
\mathbf{c}=\left[\begin{array}{c}
c_{1} \\
c_{2} \\
\vdots \\
c_{n}
\end{array}\right]
$$

and

$$
Y=\left[\begin{array}{llll}
\mathbf{y}_{1} & \mathbf{y}_{2} & \cdots & \mathbf{y}_{n}
\end{array}\right]=\left[\begin{array}{cccc}
y_{11} & y_{12} & \cdots & y_{1 n}  \tag{10.3.4}\\
y_{21} & y_{22} & \cdots & y_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
y_{n 1} & y_{n 2} & \cdots & y_{n n}
\end{array}\right]
$$

that is, the columns of $Y$ are the vector functions $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$.

For reference below, note that

$$
\begin{aligned}
Y^{\prime} & =\left[\begin{array}{llll}
\mathbf{y}_{1}^{\prime} & \mathbf{y}_{2}^{\prime} & \cdots & \mathbf{y}_{n}^{\prime}
\end{array}\right] \\
& =\left[\begin{array}{lllll}
A \mathbf{y}_{1} A \mathbf{y}_{2} & \cdots & A \mathbf{y}_{n}
\end{array}\right] \\
& =A\left[\begin{array}{llll}
\mathbf{y}_{1} & \mathbf{y}_{2} & \cdots & \mathbf{y}_{n}
\end{array}\right]=A Y
\end{aligned}
$$

that is, $Y$ satisfies the matrix differential equation

$$
Y^{\prime}=A Y
$$

The determinant of $Y$,

$$
W \xlongequal{ }\left|\begin{array}{cccc}
y_{11} & y_{12} & \cdots & y_{1 n}  \tag{10.3.5}\\
y_{21} & y_{22} & \cdots & y_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
y_{n 1} & y_{n 2} & \cdots & y_{n n}
\end{array}\right|
$$

is called the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$. It can be shown (Exercises 2 and 3 ) that this definition is analogous to definitions of the Wronskian of scalar functions given in Sections 5.1 and 9.1. The next theorem is analogous to Theorems 5.1.4 and ??. The proof is sketched in Exercise 4 for $n=2$ and in Exercise 5 for general $n$.

Theorem 10.3.2 [Abel's Formula] Suppose the $n \times n$ matrix $A=A(t)$ is continuous on $(a, b)$, let $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ be solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$, and let $t_{0}$ be in $(a, b)$. Then the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is given by

$$
\begin{equation*}
W(t)=W\left(t_{0}\right) \exp \left(\int_{t_{0}}^{t}\left[a_{11}(s)+a_{22}(s)+\cdots+a_{n n}(s)\right] d s\right), \quad a<t<b \tag{10.3.6}
\end{equation*}
$$

Therefore, either $W$ has no zeros in $(a, b)$ or $W \equiv 0$ on $(a, b)$.
REMARK: The sum of the diagonal entries of a square matrix $A$ is called the trace of $A$, denoted by $\operatorname{tr}(A)$. Thus, for an $n \times n$ matrix $A$,

$$
\operatorname{tr}(A)=a_{11}+a_{22}+\cdots+a_{n n}
$$

and (10.3.6) can be written as

$$
W(t)=W\left(t_{0}\right) \exp \left(\int_{t_{0}}^{t} \operatorname{tr}(A(s)) d s\right), \quad a<t<b .
$$

The next theorem is analogous to Theorems 5.1.6 and ??.
Theorem 10.3.3 Suppose the $n \times n$ matrix $A=A(t)$ is continuous on $(a, b)$ and let $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ be solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$. Then the following statements are equivalent; that is, they are either all true or all false:
(a) The general solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$ is $\mathbf{y}=c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}+\cdots+c_{n} \mathbf{y}_{n}$, where $c_{1}, c_{2}, \ldots$, $c_{n}$ are arbitrary constants.
(b) $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is a fundamental set of solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$.
(c) $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is linearly independent on $(a, b)$.
(d) The Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is nonzero at some point in $(a, b)$.
(e) The Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is nonzero at all points in $(a, b)$.

We say that $Y$ in (10.3.4) is a fundamental matrix for $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ if any (and therefore all) of the statements (a)-(e) of Theorem 10.3.2 are true for the columns of $Y$. In this case, (10.3.3) implies that the general solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ can be written as $\mathbf{y}=Y \mathbf{c}$, where $\mathbf{c}$ is an arbitrary constant $n$-vector.

Example 10.3.2 The vector functions

$$
\mathbf{y}_{1}=\left[\begin{array}{c}
-e^{2 t} \\
2 e^{2 t}
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=\left[\begin{array}{r}
-e^{-t} \\
e^{-t}
\end{array}\right]
$$

are solutions of the constant coefficient system

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rr}
-4 & -3  \tag{10.3.7}\\
6 & 5
\end{array}\right] \mathbf{y}
$$

on $(-\infty, \infty)$. (Verify.)
(a) Compute the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}\right\}$ directly from the definition (10.3.5)
(b) Verify Abel's formula (10.3.6) for the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}\right\}$.
(c) Find the general solution of (10.3.7).
(d) Solve the initial value problem

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rr}
-4 & -3  \tag{10.3.8}\\
6 & 5
\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}
4 \\
-5
\end{array}\right] .
$$

$\underline{\text { SOLUTION(a) }}$ From (10.3.5)

$$
W(t)=\left|\begin{array}{cc}
-e^{2 t} & -e^{-t}  \tag{10.3.9}\\
2 e^{2 t} & e^{-t}
\end{array}\right|=e^{2 t} e^{-t}\left[\begin{array}{rr}
-1 & -1 \\
2 & 1
\end{array}\right]=e^{t} .
$$

$\underline{\text { SOLUTION(b) }}$ Here

$$
A=\left[\begin{array}{rr}
-4 & -3 \\
6 & 5
\end{array}\right]
$$

so $\operatorname{tr}(A)=-4+5=1$. If $t_{0}$ is an arbitrary real number then (10.3.6) implies that

$$
W(t)=W\left(t_{0}\right) \exp \left(\int_{t_{0}}^{t} 1 d s\right)=\left|\begin{array}{cc}
-e^{2 t_{0}} & -e^{-t_{0}} \\
2 e^{2 t_{0}} & e^{-t_{0}}
\end{array}\right| e^{\left(t-t_{0}\right)}=e^{t_{0}} e^{t-t_{0}}=e^{t}
$$

which is consistent with (10.3.9).
 of (10.3.7) and

$$
Y=\left[\begin{array}{cc}
-e^{2 t} & -e^{-t} \\
2 e^{2 t} & e^{-t}
\end{array}\right]
$$

is a fundamental matrix for (10.3.7). Therefore the general solution of (10.3.7) is

$$
\mathbf{y}=c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}=c_{1}\left[\begin{array}{c}
-e^{2 t}  \tag{10.3.10}\\
2 e^{2 t}
\end{array}\right]+c_{2}\left[\begin{array}{l}
-e^{-t} \\
e^{-t}
\end{array}\right]=\left[\begin{array}{cc}
-e^{2 t} & -e^{-t} \\
2 e^{2 t} & e^{-t}
\end{array}\right]\left[\begin{array}{l}
c_{1} \\
c_{2}
\end{array}\right] .
$$

$\underline{\text { SOLUTION(d) }}$ Setting $t=0$ in (10.3.10) and imposing the initial condition in (10.3.8) yields

$$
c_{1}\left[\begin{array}{r}
-1 \\
2
\end{array}\right]+c_{2}\left[\begin{array}{r}
-1 \\
1
\end{array}\right]=\left[\begin{array}{r}
4 \\
-5
\end{array}\right] .
$$

Thus,

$$
\begin{array}{rlr}
-c_{1}-c_{2} & = & 4 \\
2 c_{1}+c_{2} & = & -5 .
\end{array}
$$

The solution of this system is $c_{1}=-1, c_{2}=-3$. Substituting these values into (10.3.10) yields

$$
\mathbf{y}=-\left[\begin{array}{c}
-e^{2 t} \\
2 e^{2 t}
\end{array}\right]-3\left[\begin{array}{c}
-e^{-t} \\
e^{-t}
\end{array}\right]=\left[\begin{array}{c}
e^{2 t}+3 e^{-t} \\
-2 e^{2 t}-3 e^{-t}
\end{array}\right]
$$

as the solution of (10.3.8).

### 10.3 Exercises

1. Prove: If $\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ are solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$, then any linear combination of $\mathbf{y}_{1}$, $\mathbf{y}_{2}, \ldots, \mathbf{y}_{n}$ is also a solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$.
2. In Section 5.1 the Wronskian of two solutions $y_{1}$ and $y_{2}$ of the scalar second order equation

$$
\begin{equation*}
P_{0}(x) y^{\prime \prime}+P_{1}(x) y^{\prime}+P_{2}(x) y=0 \tag{A}
\end{equation*}
$$

was defined to be

$$
W=\left|\begin{array}{ll}
y_{1} & y_{2} \\
y_{1}^{\prime} & y_{2}^{\prime}
\end{array}\right|
$$

(a) Rewrite (A) as a system of first order equations and show that $W$ is the Wronskian (as defined in this section) of two solutions of this system.
(b) Apply Eqn. (10.3.6) to the system derived in (a), and show that

$$
W(x)=W\left(x_{0}\right) \exp \left\{-\int_{x_{0}}^{x} \frac{P_{1}(s)}{P_{0}(s)} d s\right\}
$$

which is the form of Abel's formula given in Theorem 9.1.3.
3. In Section 9.1 the Wronskian of $n$ solutions $y_{1}, y_{2}, \ldots, y_{n}$ of the $n$-th order equation

$$
\begin{equation*}
P_{0}(x) y^{(n)}+P_{1}(x) y^{(n-1)}+\cdots+P_{n}(x) y=0 \tag{A}
\end{equation*}
$$

was defined to be

$$
W=\left|\begin{array}{cccc}
y_{1} & y_{2} & \cdots & y_{n} \\
y_{1}^{\prime} & y_{2}^{\prime} & \cdots & y_{n}^{\prime} \\
\vdots & \vdots & \ddots & \vdots \\
y_{1}^{(n-1)} & y_{2}^{(n-1)} & \cdots & y_{n}^{(n-1)}
\end{array}\right| .
$$

(a) Rewrite (A) as a system of first order equations and show that $W$ is the Wronskian (as defined in this section) of $n$ solutions of this system.
(b) Apply Eqn. (10.3.6) to the system derived in (a), and show that

$$
W(x)=W\left(x_{0}\right) \exp \left\{-\int_{x_{0}}^{x} \frac{P_{1}(s)}{P_{0}(s)} d s\right\},
$$

which is the form of Abel's formula given in Theorem 9.1.3.
4. Suppose

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
y_{11} \\
y_{21}
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=\left[\begin{array}{l}
y_{12} \\
y_{22}
\end{array}\right]
$$

are solutions of the $2 \times 2$ system $\mathbf{y}^{\prime}=A \mathbf{y}$ on $(a, b)$, and let

$$
Y=\left[\begin{array}{ll}
y_{11} & y_{12} \\
y_{21} & y_{22}
\end{array}\right] \quad \text { and } \quad W=\left|\begin{array}{ll}
y_{11} & y_{12} \\
y_{21} & y_{22}
\end{array}\right|
$$

thus, $W$ is the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}\right\}$.
(a) Deduce from the definition of determinant that

$$
W^{\prime}=\left|\begin{array}{ll}
y_{11}^{\prime} & y_{12}^{\prime} \\
y_{21} & y_{22}
\end{array}\right|+\left|\begin{array}{ll}
y_{11} & y_{12} \\
y_{21}^{\prime} & y_{22}^{\prime}
\end{array}\right| .
$$

(b) Use the equation $Y^{\prime}=A(t) Y$ and the definition of matrix multiplication to show that

$$
\left[\begin{array}{ll}
y_{11}^{\prime} & y_{12}^{\prime}
\end{array}\right]=a_{11}\left[\begin{array}{ll}
y_{11} & y_{12}
\end{array}\right]+a_{12}\left[\begin{array}{ll}
y_{21} & y_{22}
\end{array}\right]
$$

and

$$
\left[\begin{array}{ll}
y_{21}^{\prime} & y_{22}^{\prime}
\end{array}\right]=a_{21}\left[\begin{array}{ll}
y_{11} & y_{12}
\end{array}\right]+a_{22}\left[\begin{array}{ll}
y_{21} & y_{22}
\end{array}\right]
$$

(c) Use properties of determinants to deduce from (a) and (a) that

$$
\left|\begin{array}{ll}
y_{11}^{\prime} & y_{12}^{\prime} \\
y_{21} & y_{22}
\end{array}\right|=a_{11} W \quad \text { and } \quad\left|\begin{array}{cc}
y_{11} & y_{12} \\
y_{21}^{\prime} & y_{22}^{\prime}
\end{array}\right|=a_{22} W .
$$

(d) Conclude from (c) that

$$
W^{\prime}=\left(a_{11}+a_{22}\right) W
$$

and use this to show that if $a<t_{0}<b$ then

$$
W(t)=W\left(t_{0}\right) \exp \left(\int_{t_{0}}^{t}\left[a_{11}(s)+a_{22}(s)\right] d s\right) \quad a<t<b
$$

5. Suppose the $n \times n$ matrix $A=A(t)$ is continuous on $(a, b)$. Let

$$
Y=\left[\begin{array}{cccc}
y_{11} & y_{12} & \cdots & y_{1 n} \\
y_{21} & y_{22} & \cdots & y_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
y_{n 1} & y_{n 2} & \cdots & y_{n n}
\end{array}\right]
$$

where the columns of $Y$ are solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$. Let

$$
r_{i}=\left[y_{i 1} y_{i 2} \ldots y_{i n}\right]
$$

be the $i$ th row of $Y$, and let $W$ be the determinant of $Y$.
(a) Deduce from the definition of determinant that

$$
W^{\prime}=W_{1}+W_{2}+\cdots+W_{n}
$$

where, for $1 \leq m \leq n$, the $i$ th row of $W_{m}$ is $r_{i}$ if $i \neq m$, and $r_{m}^{\prime}$ if $i=m$.
(b) Use the equation $Y^{\prime}=A Y$ and the definition of matrix multiplication to show that

$$
r_{m}^{\prime}=a_{m 1} r_{1}+a_{m 2} r_{2}+\cdots+a_{m n} r_{n} .
$$

(c) Use properties of determinants to deduce from (b) that

$$
\operatorname{det}\left(W_{m}\right)=a_{m m} W
$$

(d) Conclude from (a) and (c) that

$$
W^{\prime}=\left(a_{11}+a_{22}+\cdots+a_{n n}\right) W
$$

and use this to show that if $a<t_{0}<b$ then

$$
W(t)=W\left(t_{0}\right) \exp \left(\int_{t_{0}}^{t}\left[a_{11}(s)+a_{22}(s)+\cdots+a_{n n}(s)\right] d s\right), \quad a<t<b
$$

6. Suppose the $n \times n$ matrix $A$ is continuous on $(a, b)$ and $t_{0}$ is a point in $(a, b)$. Let $Y$ be a fundamental matrix for $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$.
(a) Show that $Y\left(t_{0}\right)$ is invertible.
(b) Show that if $\mathbf{k}$ is an arbitrary $n$-vector then the solution of the initial value problem

$$
\mathbf{y}^{\prime}=A(t) \mathbf{y}, \quad \mathbf{y}\left(t_{0}\right)=\mathbf{k}
$$

is

$$
\mathbf{y}=Y(t) Y^{-1}\left(t_{0}\right) \mathbf{k}
$$

7. Let

$$
A=\left[\begin{array}{ll}
2 & 4 \\
4 & 2
\end{array}\right], \quad \mathbf{y}_{1}=\left[\begin{array}{l}
e^{6 t} \\
e^{6 t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
e^{-2 t} \\
-e^{-2 t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{r}
-3 \\
9
\end{array}\right]
$$

(a) Verify that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}\right\}$ is a fundamental set of solutions for $\mathbf{y}^{\prime}=A \mathbf{y}$.
(b) Solve the initial value problem

$$
\begin{equation*}
\mathbf{y}^{\prime}=A \mathbf{y}, \quad \mathbf{y}(0)=\mathbf{k} . \tag{A}
\end{equation*}
$$

(c) Use the result of Exercise 6(b) to find a formula for the solution of (A) for an arbitrary initial vector $\mathbf{k}$.
8. Repeat Exercise 7 with

$$
A=\left[\begin{array}{rr}
-2 & -2 \\
-5 & 1
\end{array}\right], \quad \mathbf{y}_{1}=\left[\begin{array}{l}
e^{-4 t} \\
e^{-4 t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{r}
-2 e^{3 t} \\
5 e^{3 t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{c}
10 \\
-4
\end{array}\right]
$$

9. Repeat Exercise 7 with

$$
A=\left[\begin{array}{rr}
-4 & -10 \\
3 & 7
\end{array}\right], \quad \mathbf{y}_{1}=\left[\begin{array}{r}
-5 e^{2 t} \\
3 e^{2 t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
2 e^{t} \\
-e^{t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{r}
-19 \\
11
\end{array}\right] .
$$

10. Repeat Exercise 7 with

$$
A=\left[\begin{array}{ll}
2 & 1 \\
1 & 2
\end{array}\right], \quad \mathbf{y}_{1}=\left[\begin{array}{c}
e^{3 t} \\
e^{3 t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{r}
e^{t} \\
-e^{t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{l}
2 \\
8
\end{array}\right] .
$$

11. Let

$$
\begin{aligned}
A & =\left[\begin{array}{rrr}
3 & -1 & -1 \\
-2 & 3 & 2 \\
4 & -1 & -2
\end{array}\right], \\
\mathbf{y}_{1} & =\left[\begin{array}{c}
e^{2 t} \\
0 \\
e^{2 t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
e^{3 t} \\
-e^{3 t} \\
e^{3 t}
\end{array}\right], \quad \mathbf{y}_{3}=\left[\begin{array}{c}
e^{-t} \\
-3 e^{-t} \\
7 e^{-t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{c}
2 \\
-7 \\
20
\end{array}\right] .
\end{aligned}
$$

(a) Verify that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ is a fundamental set of solutions for $\mathbf{y}^{\prime}=A \mathbf{y}$.
(b) Solve the initial value problem

$$
\begin{equation*}
\mathbf{y}^{\prime}=A \mathbf{y}, \quad \mathbf{y}(0)=\mathbf{k} . \tag{A}
\end{equation*}
$$

(c) Use the result of Exercise 6(b) to find a formula for the solution of (A) for an arbitrary initial vector $\mathbf{k}$.
12. Repeat Exercise 11 with

$$
\begin{aligned}
A & =\left[\begin{array}{lll}
0 & 2 & 2 \\
2 & 0 & 2 \\
2 & 2 & 0
\end{array}\right], \\
\mathbf{y}_{1} & =\left[\begin{array}{c}
-e^{-2 t} \\
0 \\
e^{-2 t}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
-e^{-2 t} \\
e^{-2 t} \\
0
\end{array}\right], \quad \mathbf{y}_{3}=\left[\begin{array}{c}
e^{4 t} \\
e^{4 t} \\
e^{4 t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{r}
0 \\
-9 \\
12
\end{array}\right] .
\end{aligned}
$$

13. Repeat Exercise 11 with

$$
\begin{aligned}
A & =\left[\begin{array}{rrr}
-1 & 2 & 3 \\
0 & 1 & 6 \\
0 & 0 & -2
\end{array}\right], \\
\mathbf{y}_{1} & =\left[\begin{array}{c}
e^{t} \\
e^{t} \\
0
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
e^{-t} \\
0 \\
0
\end{array}\right], \quad \mathbf{y}_{3}=\left[\begin{array}{c}
e^{-2 t} \\
-2 e^{-2 t} \\
e^{-2 t}
\end{array}\right], \quad \mathbf{k}=\left[\begin{array}{c}
5 \\
5 \\
-1
\end{array}\right] .
\end{aligned}
$$

14. Suppose $Y$ and $Z$ are fundamental matrices for the $n \times n$ system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$. Then some of the four matrices $Y Z^{-1}, Y^{-1} Z, Z^{-1} Y, Z Y^{-1}$ are necessarily constant. Identify them and prove that they are constant.
15. Suppose the columns of an $n \times n$ matrix $Y$ are solutions of the $n \times n$ system $\mathbf{y}^{\prime}=A \mathbf{y}$ and $C$ is an $n \times n$ constant matrix.
(a) Show that the matrix $Z=Y C$ satisfies the differential equation $Z^{\prime}=A Z$.
(b) Show that $Z$ is a fundamental matrix for $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ if and only if $C$ is invertible and $Y$ is a fundamental matrix for $\mathbf{y}^{\prime}=A(t) \mathbf{y}$.
16. Suppose the $n \times n$ matrix $A=A(t)$ is continuous on $(a, b)$ and $t_{0}$ is in $(a, b)$. For $i=1,2, \ldots$, $n$, let $\mathbf{y}_{i}$ be the solution of the initial value problem $\mathbf{y}_{i}^{\prime}=A(t) \mathbf{y}_{i}, \mathbf{y}_{i}\left(t_{0}\right)=\mathbf{e}_{i}$, where

$$
\mathbf{e}_{1}=\left[\begin{array}{c}
1 \\
0 \\
\vdots \\
0
\end{array}\right], \quad \mathbf{e}_{2}=\left[\begin{array}{c}
0 \\
1 \\
\vdots \\
0
\end{array}\right], \quad \cdots \quad \mathbf{e}_{n}=\left[\begin{array}{c}
0 \\
0 \\
\vdots \\
1
\end{array}\right]
$$

that is, the $j$ th component of $\mathbf{e}_{i}$ is 1 if $j=i$, or 0 if $j \neq i$.
(a) Show that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is a fundamental set of solutions of $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$.
(b) Conclude from (a) and Exercise 15 that $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ has infinitely many fundamental sets of solutions on $(a, b)$.
17. Show that $Y$ is a fundamental matrix for the system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ if and only if $Y^{-1}$ is a fundamental matrix for $\mathbf{y}^{\prime}=-A^{T}(t) \mathbf{y}$, where $A^{T}$ denotes the transpose of $A$. Hint: See Exercise 11.
18. Let $Z$ be the fundamental matrix for the constant coefficient system $\mathbf{y}^{\prime}=A \mathbf{y}$ such that $Z(0)=I$.
(a) Show that $Z(t) Z(s)=Z(t+s)$ for all $s$ and $t$. Hint: For fixed s let $\Gamma_{1}(t)=Z(t) Z(s)$ and $\Gamma_{2}(t)=Z(t+s)$. Show that $\Gamma_{1}$ and $\Gamma_{2}$ are both solutions of the matrix initial value problem $\Gamma^{\prime}=A \Gamma, \quad \Gamma(0)=Z(s)$. Then conclude from Theorem 10.2.1 that $\Gamma_{1}=\Gamma_{2}$.
(b) Show that $(Z(t))^{-1}=Z(-t)$.
(c) The matrix $Z$ defined above is sometimes denoted by $e^{t A}$. Discuss the motivation for this notation.

### 10.4 CONSTANT COEFFICIENT HOMOGENEOUS SYSTEMS I

We'll now begin our study of the homogeneous system

$$
\begin{equation*}
\mathbf{y}^{\prime}=A \mathbf{y} \tag{10.4.1}
\end{equation*}
$$

where $A$ is an $n \times n$ constant matrix. Since $A$ is continuous on $(-\infty, \infty)$, Theorem 10.2.1 implies that all solutions of (10.4.1) are defined on $(-\infty, \infty)$. Therefore, when we speak of solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$, we'll mean solutions on $(-\infty, \infty)$.

In this section we assume that all the eigenvalues of $A$ are real and that $A$ has a set of $n$ linearly independent eigenvectors. In the next two sections we consider the cases where some of the eigenvalues of $A$ are complex, or where $A$ does not have $n$ linearly independent eigenvectors.

In Example 10.3.2 we showed that the vector functions

$$
\mathbf{y}_{1}=\left[\begin{array}{c}
-e^{2 t} \\
2 e^{2 t}
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=\left[\begin{array}{l}
-e^{-t} \\
e^{-t}
\end{array}\right]
$$

form a fundamental set of solutions of the system

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rr}
-4 & -3  \tag{10.4.2}\\
6 & 5
\end{array}\right] \mathbf{y}
$$

but we did not show how we obtained $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ in the first place. To see how these solutions can be obtained we write (10.4.2) as

$$
\begin{align*}
& y_{1}^{\prime}=-4 y_{1}-3 y_{2}  \tag{10.4.3}\\
& y_{2}^{\prime}=6 y_{1}+5 y_{2}
\end{align*}
$$

and look for solutions of the form

$$
\begin{equation*}
y_{1}=x_{1} e^{\lambda t} \quad \text { and } \quad y_{2}=x_{2} e^{\lambda t} \tag{10.4.4}
\end{equation*}
$$

where $x_{1}, x_{2}$, and $\lambda$ are constants to be determined. Differentiating (10.4.4) yields

$$
y_{1}^{\prime}=\lambda x_{1} e^{\lambda t} \quad \text { and } \quad y_{2}^{\prime}=\lambda x_{2} e^{\lambda t}
$$

Substituting this and (10.4.4) into (10.4.3) and canceling the common factor $e^{\lambda t}$ yields

$$
\begin{aligned}
-4 x_{1}-3 x_{2} & =\lambda x_{1} \\
6 x_{1}+5 x_{2} & =\lambda x_{2}
\end{aligned}
$$

For a given $\lambda$, this is a homogeneous algebraic system, since it can be rewritten as

$$
\begin{align*}
(-4-\lambda) x_{1}-3 x_{2} & =0  \tag{10.4.5}\\
6 x_{1}+(5-\lambda) x_{2} & =0 .
\end{align*}
$$

The trivial solution $x_{1}=x_{2}=0$ of this system isn't useful, since it corresponds to the trivial solution $y_{1} \equiv y_{2} \equiv 0$ of (10.4.3), which can't be part of a fundamental set of solutions of (10.4.2). Therefore we consider only those values of $\lambda$ for which (10.4.5) has nontrivial solutions. These are the values of $\lambda$ for which the determinant of (10.4.5) is zero; that is,

$$
\begin{aligned}
\left|\begin{array}{cc}
-4-\lambda & -3 \\
6 & 5-\lambda
\end{array}\right| & =(-4-\lambda)(5-\lambda)+18 \\
& =\lambda^{2}-\lambda-2 \\
& =(\lambda-2)(\lambda+1)=0
\end{aligned}
$$

which has the solutions $\lambda_{1}=2$ and $\lambda_{2}=-1$.
Taking $\lambda=2$ in (10.4.5) yields

$$
\begin{array}{r}
-6 x_{1}-3 x_{2}=0 \\
6 x_{1}+3 x_{2}=0
\end{array}
$$

which implies that $x_{1}=-x_{2} / 2$, where $x_{2}$ can be chosen arbitrarily. Choosing $x_{2}=2$ yields the solution $y_{1}=-e^{2 t}, y_{2}=2 e^{2 t}$ of (10.4.3). We can write this solution in vector form as

$$
\mathbf{y}_{1}=\left[\begin{array}{r}
-1  \tag{10.4.6}\\
2
\end{array}\right] e^{2 t} .
$$

Taking $\lambda=-1$ in (10.4.5) yields the system

$$
\begin{array}{r}
-3 x_{1}-3 x_{2}=0 \\
6 x_{1}+6 x_{2}=0
\end{array}
$$

so $x_{1}=-x_{2}$. Taking $x_{2}=1$ here yields the solution $y_{1}=-e^{-t}, y_{2}=e^{-t}$ of (10.4.3). We can write this solution in vector form as

$$
\mathbf{y}_{2}=\left[\begin{array}{r}
-1  \tag{10.4.7}\\
1
\end{array}\right] e^{-t}
$$

In (10.4.6) and (10.4.7) the constant coefficients in the arguments of the exponential functions are the eigenvalues of the coefficient matrix in (10.4.2), and the vector coefficients of the exponential functions are associated eigenvectors. This illustrates the next theorem.

Theorem 10.4.1 Suppose the $n \times n$ constant matrix $A$ has $n$ real eigenvalues $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{n}$ (which need not be distinct) with associated linearly independent eigenvectors $\mathbf{x}_{1}, \mathrm{x}_{2}, \ldots, \mathrm{x}_{n}$. Then the functions

$$
\mathbf{y}_{1}=\mathbf{x}_{1} e^{\lambda_{1} t}, \mathbf{y}_{2}=\mathbf{x}_{2} e^{\lambda_{2} t}, \ldots, \mathbf{y}_{n}=\mathbf{x}_{n} e^{\lambda_{n} t}
$$

form a fundamental set of solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$; that is, the general solution of this system is

$$
\mathbf{y}=c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}+c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}+\cdots+c_{n} \mathbf{x}_{n} e^{\lambda_{n} t}
$$

Proof Differentiating $\mathbf{y}_{i}=\mathbf{x}_{i} e^{\lambda_{i} t}$ and recalling that $A \mathbf{x}_{i}=\lambda_{i} \mathbf{x}_{i}$ yields

$$
\mathbf{y}_{i}^{\prime}=\lambda_{i} \mathbf{x}_{i} e^{\lambda_{i} t}=A \mathbf{x}_{i} e^{\lambda_{i} t}=A \mathbf{y}_{i}
$$

This shows that $\mathbf{y}_{i}$ is a solution of $\mathbf{y}^{\prime}=A \mathbf{y}$.
The Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is

$$
\left|\begin{array}{cccc}
x_{11} e^{\lambda_{1} t} & x_{12} e^{\lambda_{2} t} & \cdots & x_{1 n} e^{\lambda_{n} t} \\
x_{21} e^{\lambda_{1} t} & x_{22} e^{\lambda_{2} t} & \cdots & x_{2 n} e^{\lambda_{n} t} \\
\vdots & \vdots & \ddots & \vdots \\
x_{n 1} e^{\lambda_{1} t} & x_{n 2} e^{\lambda_{2} t} & \cdots & x_{n n} e^{\lambda_{n} t}
\end{array}\right|=e^{\lambda_{1} t} e^{\lambda_{2} t} \cdots e^{\lambda_{n} t}\left|\begin{array}{cccc}
x_{11} & x_{12} & \cdots & x_{1 n} \\
x_{21} & x_{22} & \cdots & x_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
x_{n 1} & x_{n 2} & \cdots & x_{n n}
\end{array}\right|
$$

Since the columns of the determinant on the right are $\mathbf{x}_{1}, \mathbf{x}_{2}, \ldots, \mathbf{x}_{n}$, which are assumed to be linearly independent, the determinant is nonzero. Therefore Theorem 10.3.3 implies that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ is a fundamental set of solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$.

## Example 10.4.1

(a) Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
2 & 4  \tag{10.4.8}\\
4 & 2
\end{array}\right] \mathbf{y} .
$$

(b) Solve the initial value problem

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
2 & 4  \tag{10.4.9}\\
4 & 2
\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}
5 \\
-1
\end{array}\right] .
$$

$\underline{\text { SOLUTION(a) }}$ The characteristic polynomial of the coefficient matrix $A$ in (10.4.8) is

$$
\begin{aligned}
\left|\begin{array}{cc}
2-\lambda & 4 \\
4 & 2-\lambda
\end{array}\right| & =(\lambda-2)^{2}-16 \\
& =(\lambda-2-4)(\lambda-2+4) \\
& =(\lambda-6)(\lambda+2)
\end{aligned}
$$

Hence, $\lambda_{1}=6$ and $\lambda_{2}=-2$ are eigenvalues of $A$. To obtain the eigenvectors, we must solve the system

$$
\left[\begin{array}{cc}
2-\lambda & 4  \tag{10.4.10}\\
4 & 2-\lambda
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right]
$$

with $\lambda=6$ and $\lambda=-2$. Setting $\lambda=6$ in (10.4.10) yields

$$
\left[\begin{array}{rr}
-4 & 4 \\
4 & -4
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right]
$$

which implies that $x_{1}=x_{2}$. Taking $x_{2}=1$ yields the eigenvector

$$
\mathbf{x}_{1}=\left[\begin{array}{l}
1 \\
1
\end{array}\right]
$$

so

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{6 t}
$$

is a solution of (10.4.8). Setting $\lambda=-2$ in (10.4.10) yields

$$
\left[\begin{array}{ll}
4 & 4 \\
4 & 4
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0
\end{array}\right]
$$

which implies that $x_{1}=-x_{2}$. Taking $x_{2}=1$ yields the eigenvector

$$
\mathbf{x}_{2}=\left[\begin{array}{r}
-1 \\
1
\end{array}\right]
$$

so

$$
\mathbf{y}_{2}=\left[\begin{array}{r}
-1 \\
1
\end{array}\right] e^{-2 t}
$$

is a solution of (10.4.8). From Theorem 10.4.1, the general solution of (10.4.8) is

$$
\mathbf{y}=c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}=c_{1}\left[\begin{array}{l}
1  \tag{10.4.11}\\
1
\end{array}\right] e^{6 t}+c_{2}\left[\begin{array}{r}
-1 \\
1
\end{array}\right] e^{-2 t}
$$

$\underline{\text { SOLUTION(b) }}$ To satisfy the initial condition in (10.4.9), we must choose $c_{1}$ and $c_{2}$ in (10.4.11) so that

$$
c_{1}\left[\begin{array}{l}
1 \\
1
\end{array}\right]+c_{2}\left[\begin{array}{r}
-1 \\
1
\end{array}\right]=\left[\begin{array}{r}
5 \\
-1
\end{array}\right] .
$$

This is equivalent to the system

$$
\begin{aligned}
c_{1}-c_{2} & =5 \\
c_{1}+c_{2} & =-1,
\end{aligned}
$$

so $c_{1}=2, c_{2}=-3$. Therefore the solution of (10.4.9) is

$$
\mathbf{y}=2\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{6 t}-3\left[\begin{array}{r}
-1 \\
1
\end{array}\right] e^{-2 t}
$$

or, in terms of components,

$$
y_{1}=2 e^{6 t}+3 e^{-2 t}, \quad y_{2}=2 e^{6 t}-3 e^{-2 t}
$$

## Example 10.4.2

(a) Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
3 & -1 & -1  \tag{10.4.12}\\
-2 & 3 & 2 \\
4 & -1 & -2
\end{array}\right] \mathbf{y} .
$$

(b) Solve the initial value problem

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
3 & -1 & -1  \tag{10.4.13}\\
-2 & 3 & 2 \\
4 & -1 & -2
\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}
2 \\
-1 \\
8
\end{array}\right]
$$

$\underline{\operatorname{Solution}(\mathbf{a})}$ The characteristic polynomial of the coefficient matrix $A$ in (10.4.12) is

$$
\left|\begin{array}{ccc}
3-\lambda & -1 & -1 \\
-2 & 3-\lambda & 2 \\
4 & -1 & -2-\lambda
\end{array}\right|=-(\lambda-2)(\lambda-3)(\lambda+1)
$$

Hence, the eigenvalues of $A$ are $\lambda_{1}=2, \lambda_{2}=3$, and $\lambda_{3}=-1$. To find the eigenvectors, we must solve the system

$$
\left[\begin{array}{ccc}
3-\lambda & -1 & -1  \tag{10.4.14}\\
-2 & 3-\lambda & 2 \\
4 & -1 & -2-\lambda
\end{array}\right]\left[\begin{array}{l}
x_{1} \\
x_{2} \\
x_{3}
\end{array}\right]=\left[\begin{array}{l}
0 \\
0 \\
0
\end{array}\right]
$$

with $\lambda=2,3,-1$. With $\lambda=2$, the augmented matrix of (10.4.14) is

$$
\left[\begin{array}{rrrcc}
1 & -1 & -1 & \vdots & 0 \\
-2 & 1 & 2 & \vdots & 0 \\
4 & -1 & -4 & \vdots & 0
\end{array}\right],
$$

which is row equivalent to

$$
\left[\begin{array}{rrrcc}
1 & 0 & -1 & \vdots & 0 \\
0 & 1 & 0 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=x_{3}$ and $x_{2}=0$. Taking $x_{3}=1$ yields

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{2 t}
$$

as a solution of (10.4.12). With $\lambda=3$, the augmented matrix of (10.4.14) is

$$
\left[\begin{array}{rrrcc}
0 & -1 & -1 & \vdots & 0 \\
-2 & 0 & 2 & \vdots & 0 \\
4 & -1 & -5 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrcc}
1 & 0 & -1 & \vdots & 0 \\
0 & 1 & 1 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=x_{3}$ and $x_{2}=-x_{3}$. Taking $x_{3}=1$ yields

$$
\mathbf{y}_{2}=\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{3 t}
$$

as a solution of (10.4.12). With $\lambda=-1$, the augmented matrix of (10.4.14) is

$$
\left[\begin{array}{rrrcc}
4 & -1 & -1 & \vdots & 0 \\
-2 & 4 & 2 & \vdots & 0 \\
4 & -1 & -1 & \vdots & 0
\end{array}\right],
$$

which is row equivalent to

$$
\left[\begin{array}{rrrcc}
1 & 0 & -\frac{1}{7} & \vdots & 0 \\
0 & 1 & \frac{3}{7} & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=x_{3} / 7$ and $x_{2}=-3 x_{3} / 7$. Taking $x_{3}=7$ yields

$$
\mathbf{y}_{3}=\left[\begin{array}{r}
1 \\
-3 \\
7
\end{array}\right] e^{-t}
$$

as a solution of (10.4.12). By Theorem 10.4.1, the general solution of (10.4.12) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{3 t}+c_{3}\left[\begin{array}{r}
1 \\
-3 \\
7
\end{array}\right] e^{-t}
$$

which can also be written as

$$
\mathbf{y}=\left[\begin{array}{ccc}
e^{2 t} & e^{3 t} & e^{-t}  \tag{10.4.15}\\
0 & -e^{3 t} & -3 e^{-t} \\
e^{2 t} & e^{3 t} & 7 e^{-t}
\end{array}\right]\left[\begin{array}{l}
c_{1} \\
c_{2} \\
c_{3}
\end{array}\right] .
$$

$\underline{\text { SOLUTION(b) }}$ To satisfy the initial condition in (10.4.13) we must choose $c_{1}, c_{2}, c_{3}$ in (10.4.15) so that

$$
\left[\begin{array}{rrr}
1 & 1 & 1 \\
0 & -1 & -3 \\
1 & 1 & 7
\end{array}\right]\left[\begin{array}{l}
c_{1} \\
c_{2} \\
c_{3}
\end{array}\right]=\left[\begin{array}{r}
2 \\
-1 \\
8
\end{array}\right] .
$$

Solving this system yields $c_{1}=3, c_{2}=-2, c_{3}=1$. Hence, the solution of (10.4.13) is

$$
\begin{aligned}
\mathbf{y} & =\left[\begin{array}{ccc}
e^{2 t} & e^{3 t} & e^{-t} \\
0 & -e^{3 t} & -3 e^{-t} \\
e^{2 t} & e^{3 t} & 7 e^{-t}
\end{array}\right]\left[\begin{array}{r}
3 \\
-2 \\
1
\end{array}\right] \\
& =3\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{2 t}-2\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{3 t}+\left[\begin{array}{r}
1 \\
-3 \\
7
\end{array}\right] e^{-t} .
\end{aligned}
$$

Example 10.4.3 Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
-3 & 2 & 2  \tag{10.4.16}\\
2 & -3 & 2 \\
2 & 2 & -3
\end{array}\right] \mathbf{y} .
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.4.16) is

$$
\left|\begin{array}{ccc}
-3-\lambda & 2 & 2 \\
2 & -3-\lambda & 2 \\
2 & 2 & -3-\lambda
\end{array}\right|=-(\lambda-1)(\lambda+5)^{2}
$$

Hence, $\lambda_{1}=1$ is an eigenvalue of multiplicity 1 , while $\lambda_{2}=-5$ is an eigenvalue of multiplicity 2 . Eigenvectors associated with $\lambda_{1}=1$ are solutions of the system with augmented matrix

$$
\left[\begin{array}{rrrcc}
-4 & 2 & 2 & \vdots & 0 \\
2 & -4 & 2 & \vdots & 0 \\
2 & 2 & -4 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrll}
1 & 0 & -1 & \vdots & 0 \\
0 & 1 & -1 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=x_{2}=x_{3}$, and we choose $x_{3}=1$ to obtain the solution

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1  \tag{10.4.17}\\
1 \\
1
\end{array}\right] e^{t}
$$

of (10.4.16). Eigenvectors associated with $\lambda_{2}=-5$ are solutions of the system with augmented matrix

$$
\left[\begin{array}{ccccc}
2 & 2 & 2 & \vdots & 0 \\
2 & 2 & 2 & \vdots & 0 \\
2 & 2 & 2 & \vdots & 0
\end{array}\right]
$$

Hence, the components of these eigenvectors need only satisfy the single condition

$$
x_{1}+x_{2}+x_{3}=0
$$

Since there's only one equation here, we can choose $x_{2}$ and $x_{3}$ arbitrarily. We obtain one eigenvector by choosing $x_{2}=0$ and $x_{3}=1$, and another by choosing $x_{2}=1$ and $x_{3}=0$. In both cases $x_{1}=-1$. Therefore

$$
\left[\begin{array}{r}
-1 \\
0 \\
1
\end{array}\right] \quad \text { and } \quad\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right]
$$

are linearly independent eigenvectors associated with $\lambda_{2}=-5$, and the corresponding solutions of (10.4.16) are

$$
\mathbf{y}_{2}=\left[\begin{array}{r}
-1 \\
0 \\
1
\end{array}\right] e^{-5 t} \quad \text { and } \quad \mathbf{y}_{3}=\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] e^{-5 t}
$$

Because of this and (10.4.17), Theorem 10.4.1 implies that the general solution of (10.4.16) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right] e^{t}+c_{2}\left[\begin{array}{r}
-1 \\
0 \\
1
\end{array}\right] e^{-5 t}+c_{3}\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] e^{-5 t} .
$$

Geometric Properties of Solutions when $n=2$
We'll now consider the geometric properties of solutions of a $2 \times 2$ constant coefficient system

$$
\left[\begin{array}{l}
y_{1}^{\prime}  \tag{10.4.18}\\
y_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]\left[\begin{array}{l}
y_{1} \\
y_{2}
\end{array}\right]
$$

It is convenient to think of a " $y_{1}-y_{2}$ plane," where a point is identified by rectangular coordinates $\left(y_{1}, y_{2}\right)$. If $\mathbf{y}=\left[\begin{array}{l}y_{1} \\ y_{2}\end{array}\right]$ is a non-constant solution of (10.4.18), then the point $\left(y_{1}(t), y_{2}(t)\right)$ moves along a curve $C$ in the $y_{1}-y_{2}$ plane as $t$ varies from $-\infty$ to $\infty$. We call $C$ the trajectory of $\mathbf{y}$. (We also say that $C$ is a trajectory of the system (10.4.18).) I's important to note that $C$ is the trajectory of infinitely many solutions of (10.4.18), since if $\tau$ is any real number, then $\mathbf{y}(t-\tau)$ is a solution of (10.4.18) (Exercise 28(b)), and $\left(y_{1}(t-\tau), y_{2}(t-\tau)\right)$ also moves along $C$ as $t$ varies from $-\infty$ to $\infty$. Moreover, Exercise 28(c) implies that distinct trajectories of (10.4.18) can't intersect, and that two solutions $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ of (10.4.18) have the same trajectory if and only if $\mathbf{y}_{2}(t)=\mathbf{y}_{1}(t-\tau)$ for some $\tau$.

From Exercise 28(a), a trajectory of a nontrivial solution of (10.4.18) can't contain ( 0,0 ), which we define to be the trajectory of the trivial solution $\mathbf{y} \equiv 0$. More generally, if $\mathbf{y}=\left[\begin{array}{l}k_{1} \\ k_{2}\end{array}\right] \neq \mathbf{0}$ is a constant solution of (10.4.18) (which could occur if zero is an eigenvalue of the matrix of (10.4.18)), we define the trajectory of $\mathbf{y}$ to be the single point $\left(k_{1}, k_{2}\right)$.

To be specific, this is the question: What do the trajectories look like, and how are they traversed? In this section we'll answer this question, assuming that the matrix

$$
A=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]
$$

of (10.4.18) has real eigenvalues $\lambda_{1}$ and $\lambda_{2}$ with associated linearly independent eigenvectors $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$. Then the general solution of $(10.4 .18)$ is

$$
\begin{equation*}
\mathbf{y}=c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}+c_{2} \mathbf{x}_{2} e^{\lambda_{2} t} \tag{10.4.19}
\end{equation*}
$$

We'll consider other situations in the next two sections.
We leave it to you (Exercise 35) to classify the trajectories of (10.4.18) if zero is an eigenvalue of $A$. We'll confine our attention here to the case where both eigenvalues are nonzero. In this case the simplest situation is where $\lambda_{1}=\lambda_{2} \neq 0$, so (10.4.19) becomes

$$
\mathbf{y}=\left(c_{1} \mathbf{x}_{1}+c_{2} \mathbf{x}_{2}\right) e^{\lambda_{1} t}
$$

Since $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$ are linearly independent, an arbitrary vector $\mathbf{x}$ can be written as $\mathbf{x}=c_{1} \mathbf{x}_{1}+c_{2} \mathbf{x}_{2}$. Therefore the general solution of (10.4.18) can be written as $\mathbf{y}=\mathbf{x} e^{\lambda_{1} t}$ where $\mathbf{x}$ is an arbitrary 2 vector, and the trajectories of nontrivial solutions of (10.4.18) are half-lines through (but not including) the origin. The direction of motion is away from the origin if $\lambda_{1}>0$ (Figure 10.4.1), toward it if $\lambda_{1}<0$ (Figure 10.4.2). (In these and the next figures an arrow through a point indicates the direction of motion along the trajectory through the point.)


Figure 10.4.1 Trajectories of a $2 \times 2$ system with a Figure 10.4.2 Trajectories of a $2 \times 2$ system with a repeated positive eigenvalue repeated negative eigenvalue

Now suppose $\lambda_{2}>\lambda_{1}$, and let $L_{1}$ and $L_{2}$ denote lines through the origin parallel to $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$, respectively. By a half-line of $L_{1}$ (or $L_{2}$ ), we mean either of the rays obtained by removing the origin from $L_{1}$ (or $L_{2}$ ).

Letting $c_{2}=0$ in (10.4.19) yields $\mathbf{y}=c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}$. If $c_{1} \neq 0$, the trajectory defined by this solution is a half-line of $L_{1}$. The direction of motion is away from the origin if $\lambda_{1}>0$, toward the origin if $\lambda_{1}<0$. Similarly, the trajectory of $\mathbf{y}=c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}$ with $c_{2} \neq 0$ is a half-line of $L_{2}$.

Henceforth, we assume that $c_{1}$ and $c_{2}$ in (10.4.19) are both nonzero. In this case, the trajectory of (10.4.19) can't intersect $L_{1}$ or $L_{2}$, since every point on these lines is on the trajectory of a solution for which either $c_{1}=0$ or $c_{2}=0$. (Remember: distinct trajectories can't intersect!). Therefore the trajectory of (10.4.19) must lie entirely in one of the four open sectors bounded by $L_{1}$ and $L_{2}$, but do not any point on $L_{1}$ or $L_{2}$. Since the initial point $\left(y_{1}(0), y_{2}(0)\right)$ defined by

$$
\mathbf{y}(0)=c_{1} \mathbf{x}_{1}+c_{2} \mathbf{x}_{2}
$$

is on the trajectory, we can determine which sector contains the trajectory from the signs of $c_{1}$ and $c_{2}$, as shown in Figure 10.4.3.
The direction of $\mathbf{y}(t)$ in (10.4.19) is the same as that of

$$
\begin{equation*}
e^{-\lambda_{2} t} \mathbf{y}(t)=c_{1} \mathbf{x}_{1} e^{-\left(\lambda_{2}-\lambda_{1}\right) t}+c_{2} \mathbf{x}_{2} \tag{10.4.20}
\end{equation*}
$$

and of

$$
\begin{equation*}
e^{-\lambda_{1} t} \mathbf{y}(t)=c_{1} \mathbf{x}_{1}+c_{2} \mathbf{x}_{2} e^{\left(\lambda_{2}-\lambda_{1}\right) t} \tag{10.4.21}
\end{equation*}
$$

Since the right side of (10.4.20) approaches $c_{2} \mathbf{x}_{2}$ as $t \rightarrow \infty$, the trajectory is asymptotically parallel to $L_{2}$ as $t \rightarrow \infty$. Since the right side of (10.4.21) approaches $c_{1} \mathbf{x}_{1}$ as $t \rightarrow-\infty$, the trajectory is asymptotically parallel to $L_{1}$ as $t \rightarrow-\infty$.

The shape and direction of traversal of the trajectory of (10.4.19) depend upon whether $\lambda_{1}$ and $\lambda_{2}$ are both positive, both negative, or of opposite signs. We'll now analyze these three cases.

Henceforth $\|\mathbf{u}\|$ denote the length of the vector $\mathbf{u}$.
Case 1: $\lambda_{2}>\lambda_{1}>0$
Figure 10.4.4 shows some typical trajectories. In this case, $\lim _{t \rightarrow-\infty}\|\mathbf{y}(t)\|=0$, so the trajectory is not only asymptotically parallel to $L_{1}$ as $t \rightarrow-\infty$, but is actually asymptotically tangent to $L_{1}$ at the origin. On the other hand, $\lim _{t \rightarrow \infty}\|\mathbf{y}(t)\|=\infty$ and

$$
\lim _{t \rightarrow \infty}\left\|\mathbf{y}(t)-c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}\right\|=\lim _{t \rightarrow \infty}\left\|c_{1} \mathbf{x}_{\mathbf{1}} e^{\lambda_{1} t}\right\|=\infty
$$

so, although the trajectory is asymptotically parallel to $L_{2}$ as $t \rightarrow \infty$, it's not asymptotically tangent to $L_{2}$. The direction of motion along each trajectory is away from the origin.

Case 2: $0>\lambda_{2}>\lambda_{1}$
Figure 10.4.5 shows some typical trajectories. In this case, $\lim _{t \rightarrow \infty}\|\mathbf{y}(t)\|=0$, so the trajectory is asymptotically tangent to $L_{2}$ at the origin as $t \rightarrow \infty$. On the other hand, $\lim _{t \rightarrow-\infty}\|\mathbf{y}(t)\|=\infty$ and

$$
\lim _{t \rightarrow-\infty}\left\|\mathbf{y}(t)-c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}\right\|=\lim _{t \rightarrow-\infty}\left\|c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}\right\|=\infty
$$



Figure 10.4.3 Four open sectors bounded by $L_{1}$ and $L_{2}$


Figure 10.4.4 Two positive eigenvalues; motion away from origin
so, although the trajectory is asymptotically parallel to $L_{1}$ as $t \rightarrow-\infty$, it's not asymptotically tangent to it. The direction of motion along each trajectory is toward the origin.


Figure 10.4.5 Two negative eigenvalues; motion toward the origin


Figure 10.4.6 Eigenvalues of different signs

Case 3: $\lambda_{2}>0>\lambda_{1}$
Figure 10.4.6 shows some typical trajectories. In this case,

$$
\lim _{t \rightarrow \infty}\|\mathbf{y}(t)\|=\infty \quad \text { and } \quad \lim _{t \rightarrow \infty}\left\|\mathbf{y}(t)-c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}\right\|=\lim _{t \rightarrow \infty}\left\|c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}\right\|=0
$$

so the trajectory is asymptotically tangent to $L_{2}$ as $t \rightarrow \infty$. Similarly,

$$
\lim _{t \rightarrow-\infty}\|\mathbf{y}(t)\|=\infty \quad \text { and } \quad \lim _{t \rightarrow-\infty}\left\|\mathbf{y}(t)-c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}\right\|=\lim _{t \rightarrow-\infty}\left\|c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}\right\|=0
$$

so the trajectory is asymptotically tangent to $L_{1}$ as $t \rightarrow-\infty$. The direction of motion is toward the origin on $L_{1}$ and away from the origin on $L_{2}$. The direction of motion along any other trajectory is away from $L_{1}$, toward $L_{2}$.

### 10.4 Exercises

In Exercises 1-15 find the general solution.

1. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}1 & 2 \\ 2 & 1\end{array}\right] \mathbf{y}$
2. $\mathbf{y}^{\prime}=\frac{1}{4}\left[\begin{array}{rr}-5 & 3 \\ 3 & -5\end{array}\right] \mathbf{y}$
3. $\mathbf{y}^{\prime}=\frac{1}{5}\left[\begin{array}{rr}-4 & 3 \\ -2 & -11\end{array}\right] \mathbf{y}$
4. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-1 & -4 \\ -1 & -1\end{array}\right] \mathbf{y}$
5. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}2 & -4 \\ -1 & -1\end{array}\right] \mathbf{y}$
6. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}4 & -3 \\ 2 & -1\end{array}\right] \mathbf{y}$
7. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-6 & -3 \\ 1 & -2\end{array}\right] \mathbf{y}$
8. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}1 & -1 & -2 \\ 1 & -2 & -3 \\ -4 & 1 & -1\end{array}\right] \mathbf{y}$
9. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-6 & -4 & -8 \\ -4 & 0 & -4 \\ -8 & -4 & -6\end{array}\right] \mathbf{y}$
10. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & 5 & 8 \\ 1 & -1 & -2 \\ -1 & -1 & -1\end{array}\right] \mathbf{y}$
11. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}1 & -1 & 2 \\ 12 & -4 & 10 \\ -6 & 1 & -7\end{array}\right] \mathbf{y}$
12. $\mathbf{y}^{\prime}=\left[\begin{array}{lll}4 & -1 & -4 \\ 4 & -3 & -2 \\ 1 & -1 & -1\end{array}\right] \mathbf{y}$
13. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-2 & 2 & -6 \\ 2 & 6 & 2 \\ -2 & -2 & 2\end{array}\right] \mathbf{y}$
14. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & 2 & -2 \\ -2 & 7 & -2 \\ -10 & 10 & -5\end{array}\right] \mathbf{y}$
15. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & 1 & -1 \\ 3 & 5 & 1 \\ -6 & 2 & 4\end{array}\right] \mathbf{y}$

In Exercises 16-27 solve the initial value problem.
16. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-7 & 4 \\ -6 & 7\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}2 \\ -4\end{array}\right]$
17. $\quad \mathbf{y}^{\prime}=\frac{1}{6}\left[\begin{array}{rr}7 & 2 \\ -2 & 2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}0 \\ -3\end{array}\right]$
18. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}21 & -12 \\ 24 & -15\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}5 \\ 3\end{array}\right]$
19. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-7 & 4 \\ -6 & 7\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}-1 \\ 7\end{array}\right]$
20. $\quad \mathbf{y}^{\prime}=\frac{1}{6}\left[\begin{array}{rrr}1 & 2 & 0 \\ 4 & -1 & 0 \\ 0 & 0 & 3\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}4 \\ 7 \\ 1\end{array}\right]$
21. $\quad \mathbf{y}^{\prime}=\frac{1}{3}\left[\begin{array}{rrr}2 & -2 & 3 \\ -4 & 4 & 3 \\ 2 & 1 & 0\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}1 \\ 1 \\ 5\end{array}\right]$
22. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}6 & -3 & -8 \\ 2 & 1 & -2 \\ 3 & -3 & -5\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}0 \\ -1 \\ -1\end{array}\right]$
23. $\quad \mathbf{y}^{\prime}=\frac{1}{3}\left[\begin{array}{rrr}2 & 4 & -7 \\ 1 & 5 & -5 \\ -4 & 4 & -1\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}4 \\ 1 \\ 3\end{array}\right]$
24. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & 0 & 1 \\ 11 & -2 & 7 \\ 1 & 0 & 3\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}2 \\ 7 \\ 6\end{array}\right]$
25. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-2 & -5 & -1 \\ -4 & -1 & 1 \\ 4 & 5 & 3\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}8 \\ -10 \\ -4\end{array}\right]$
26. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{lll}3 & -1 & 0 \\ 4 & -2 & 0 \\ 4 & -4 & 2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}7 \\ 10 \\ 2\end{array}\right]$
27. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-2 & 2 & 6 \\ 2 & 6 & 2 \\ -2 & -2 & 2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}6 \\ -10 \\ 7\end{array}\right]$
28. Let $A$ be an $n \times n$ constant matrix. Then Theorem 10.2.1 implies that the solutions of

$$
\begin{equation*}
\mathbf{y}^{\prime}=A \mathbf{y} \tag{A}
\end{equation*}
$$

are all defined on $(-\infty, \infty)$.
(a) Use Theorem 10.2.1 to show that the only solution of (A) that can ever equal the zero vector is $\mathbf{y} \equiv \mathbf{0}$.
(b) Suppose $\mathbf{y}_{1}$ is a solution of (A) and $\mathbf{y}_{2}$ is defined by $\mathbf{y}_{2}(t)=\mathbf{y}_{1}(t-\tau)$, where $\tau$ is an arbitrary real number. Show that $\mathbf{y}_{2}$ is also a solution of (A).
(c) Suppose $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ are solutions of (A) and there are real numbers $t_{1}$ and $t_{2}$ such that $\mathbf{y}_{1}\left(t_{1}\right)=\mathbf{y}_{2}\left(t_{2}\right)$. Show that $\mathbf{y}_{2}(t)=\mathbf{y}_{1}(t-\tau)$ for all $t$, where $\tau=t_{2}-t_{1}$. Hint: Show that $\mathbf{y}_{1}(t-\tau)$ and $\mathbf{y}_{2}(t)$ are solutions of the same initial value problem for $(\mathrm{A})$, and apply the uniqueness assertion of Theorem 10.2.1.

In Exercises 29-34 describe and graph trajectories of the given system.
29. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}1 & 1 \\ 1 & -1\end{array}\right] \mathbf{y}$
30. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-4 & 3 \\ -2 & -11\end{array}\right] \mathbf{y}$
31. $\quad \mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}9 & -3 \\ -1 & 11\end{array}\right] \mathbf{y}$
32. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}-1 & -10 \\ -5 & 4\end{array}\right] \mathbf{y}$
33. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{cc}5 & -4 \\ 1 & 10\end{array}\right] \mathbf{y}$
34. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}-7 & 1 \\ 3 & -5\end{array}\right] \mathbf{y}$
35. Suppose the eigenvalues of the $2 \times 2$ matrix $A$ are $\lambda=0$ and $\mu \neq 0$, with corresponding eigenvectors $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$. Let $L_{1}$ be the line through the origin parallel to $\mathbf{x}_{1}$.
(a) Show that every point on $L_{1}$ is the trajectory of a constant solution of $\mathbf{y}^{\prime}=A \mathbf{y}$.
(b) Show that the trajectories of nonconstant solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$ are half-lines parallel to $\mathbf{x}_{2}$ and on either side of $L_{1}$, and that the direction of motion along these trajectories is away from $L_{1}$ if $\mu>0$, or toward $L_{1}$ if $\mu<0$.

The matrices of the systems in Exercises 36-41 are singular. Describe and graph the trajectories of nonconstant solutions of the given systems.
36. $\mathrm{C} / \mathrm{G} \mathrm{y}^{\prime}=\left[\begin{array}{rr}-1 & 1 \\ 1 & -1\end{array}\right] \mathbf{y}$
37. $\mathrm{C} / \mathrm{G} \mathrm{y}^{\prime}=\left[\begin{array}{rr}-1 & -3 \\ 2 & 6\end{array}\right] \mathbf{y}$
38. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}1 & -3 \\ -1 & 3\end{array}\right] \mathbf{y}$
39. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}1 & -2 \\ -1 & 2\end{array}\right] \mathbf{y}$
40. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}-4 & -4 \\ 1 & 1\end{array}\right] \mathbf{y}$
41. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}3 & -1 \\ -3 & 1\end{array}\right] \mathbf{y}$
42. L Let $P=P(t)$ and $Q=Q(t)$ be the populations of two species at time $t$, and assume that each population would grow exponentially if the other didn't exist; that is, in the absence of competition,

$$
\begin{equation*}
P^{\prime}=a P \quad \text { and } \quad Q^{\prime}=b Q \tag{A}
\end{equation*}
$$

where $a$ and $b$ are positive constants. One way to model the effect of competition is to assume that the growth rate per individual of each population is reduced by an amount proportional to the
other population, so (A) is replaced by

$$
\begin{aligned}
P^{\prime} & =a P-\alpha Q \\
Q^{\prime} & =-\beta P+b Q
\end{aligned}
$$

where $\alpha$ and $\beta$ are positive constants. (Since negative population doesn't make sense, this system holds only while $P$ and $Q$ are both positive.) Now suppose $P(0)=P_{0}>0$ and $Q(0)=Q_{0}>0$.
(a) For several choices of $a, b, \alpha$, and $\beta$, verify experimentally (by graphing trajectories of (A) in the $P-Q$ plane) that there's a constant $\rho>0$ (depending upon $a, b, \alpha$, and $\beta$ ) with the following properties:
(i) If $Q_{0}>\rho P_{0}$, then $P$ decreases monotonically to zero in finite time, during which $Q$ remains positive.
(ii) If $Q_{0}<\rho P_{0}$, then $Q$ decreases monotonically to zero in finite time, during which $P$ remains positive.
(b) Conclude from (a) that exactly one of the species becomes extinct in finite time if $Q_{0} \neq \rho P_{0}$. Determine experimentally what happens if $Q_{0}=\rho P_{0}$.
(c) Confirm your experimental results and determine $\gamma$ by expressing the eigenvalues and associated eigenvectors of

$$
A=\left[\begin{array}{rr}
a & -\alpha \\
-\beta & b
\end{array}\right]
$$

in terms of $a, b, \alpha$, and $\beta$, and applying the geometric arguments developed at the end of this section.

### 10.5 CONSTANT COEFFICIENT HOMOGENEOUS SYSTEMS II

We saw in Section 10.4 that if an $n \times n$ constant matrix $A$ has $n$ real eigenvalues $\lambda_{1}, \lambda_{2}, \ldots, \lambda_{n}$ (which need not be distinct) with associated linearly independent eigenvectors $\mathbf{x}_{1}, \mathbf{x}_{2}, \ldots, \mathbf{x}_{n}$, then the general solution of $\mathbf{y}^{\prime}=A \mathbf{y}$ is

$$
\mathbf{y}=c_{1} \mathbf{x}_{1} e^{\lambda_{1} t}+c_{2} \mathbf{x}_{2} e^{\lambda_{2} t}+\cdots+c_{n} \mathbf{x}_{n} e^{\lambda_{n} t}
$$

In this section we consider the case where $A$ has $n$ real eigenvalues, but does not have $n$ linearly independent eigenvectors. It is shown in linear algebra that this occurs if and only if $A$ has at least one eigenvalue of multiplicity $r>1$ such that the associated eigenspace has dimension less than $r$. In this case $A$ is said to be defective. Since it's beyond the scope of this book to give a complete analysis of systems with defective coefficient matrices, we will restrict our attention to some commonly occurring special cases.

Example 10.5.1 Show that the system

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rr}
11 & -25  \tag{10.5.1}\\
4 & -9
\end{array}\right] \mathbf{y}
$$

does not have a fundamental set of solutions of the form $\left\{\mathbf{x}_{1} e^{\lambda_{1} t}, \mathbf{x}_{2} e^{\lambda_{2} t}\right\}$, where $\lambda_{1}$ and $\lambda_{2}$ are eigenvalues of the coefficient matrix $A$ of (10.5.1) and $\mathbf{x}_{1}$, and $\mathbf{x}_{2}$ are associated linearly independent eigenvectors.

Solution The characteristic polynomial of $A$ is

$$
\begin{aligned}
\left|\begin{array}{cc}
11-\lambda & -25 \\
4 & -9-\lambda
\end{array}\right| & =(\lambda-11)(\lambda+9)+100 \\
& =\lambda^{2}-2 \lambda+1=(\lambda-1)^{2}
\end{aligned}
$$

Hence, $\lambda=1$ is the only eigenvalue of $A$. The augmented matrix of the system $(A-I) \mathbf{x}=\mathbf{0}$ is

$$
\left[\begin{array}{rrcc}
10 & -25 & \vdots & 0 \\
4 & -10 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrr}
1 & -\frac{5}{2} & \vdots & 0 \\
0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=5 x_{2} / 2$ where $x_{2}$ is arbitrary. Therefore all eigenvectors of $A$ are scalar multiples of $\mathbf{x}_{1}=$ $\left[\begin{array}{l}5 \\ 2\end{array}\right]$, so $A$ does not have a set of two linearly independent eigenvectors.
From Example 10.5.1, we know that all scalar multiples of $\mathbf{y}_{1}=\left[\begin{array}{l}5 \\ 2\end{array}\right] e^{t}$ are solutions of (10.5.1); however, to find the general solution we must find a second solution $\mathbf{y}_{2}$ such that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}\right\}$ is linearly independent. Based on your recollection of the procedure for solving a constant coefficient scalar equation

$$
a y^{\prime \prime}+b y^{\prime}+c y=0
$$

in the case where the characteristic polynomial has a repeated root, you might expect to obtain a second solution of (10.5.1) by multiplying the first solution by $t$. However, this yields $\mathbf{y}_{2}=\left[\begin{array}{l}5 \\ 2\end{array}\right] t e^{t}$, which doesn't work, since

$$
\mathbf{y}_{2}^{\prime}=\left[\begin{array}{l}
5 \\
2
\end{array}\right]\left(t e^{t}+e^{t}\right), \quad \text { while } \quad\left[\begin{array}{rr}
11 & -25 \\
4 & -9
\end{array}\right] \mathbf{y}_{2}=\left[\begin{array}{l}
5 \\
2
\end{array}\right] t e^{t}
$$

The next theorem shows what to do in this situation.
Theorem 10.5.1 Suppose the $n \times n$ matrix $A$ has an eigenvalue $\lambda_{1}$ of multiplicity $\geq 2$ and the associated eigenspace has dimension 1 ; that is, all $\lambda_{1}$-eigenvectors of $A$ are scalar multiples of an eigenvector $\mathbf{x}$. Then there are infinitely many vectors $\mathbf{u}$ such that

$$
\begin{equation*}
\left(A-\lambda_{1} I\right) \mathbf{u}=\mathbf{x} \tag{10.5.2}
\end{equation*}
$$

Moreover, if $\mathbf{u}$ is any such vector then

$$
\begin{equation*}
\mathbf{y}_{1}=\mathbf{x} e^{\lambda_{1} t} \quad \text { and } \quad \mathbf{y}_{2}=\mathbf{u} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t} \tag{10.5.3}
\end{equation*}
$$

are linearly independent solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$.
A complete proof of this theorem is beyond the scope of this book. The difficulty is in proving that there's a vector $\mathbf{u}$ satisfying (10.5.2), since $\operatorname{det}\left(A-\lambda_{1} I\right)=0$. We'll take this without proof and verify the other assertions of the theorem.

We already know that $\mathbf{y}_{1}$ in (10.5.3) is a solution of $\mathbf{y}^{\prime}=A \mathbf{y}$. To see that $\mathbf{y}_{2}$ is also a solution, we compute

$$
\begin{aligned}
\mathbf{y}_{2}^{\prime}-A \mathbf{y}_{2} & =\lambda_{1} \mathbf{u} e^{\lambda_{1} t}+\mathbf{x} e^{\lambda_{1} t}+\lambda_{1} \mathbf{x} t e^{\lambda_{1} t}-A \mathbf{u} e^{\lambda_{1} t}-A \mathbf{x} t e^{\lambda_{1} t} \\
& =\left(\lambda_{1} \mathbf{u}+\mathbf{x}-A \mathbf{u}\right) e^{\lambda_{1} t}+\left(\lambda_{1} \mathbf{x}-A \mathbf{x}\right) t e^{\lambda_{1} t}
\end{aligned}
$$

Since $A \mathbf{x}=\lambda_{1} \mathbf{x}$, this can be written as

$$
\mathbf{y}_{2}^{\prime}-A \mathbf{y}_{2}=-\left(\left(A-\lambda_{1} I\right) \mathbf{u}-\mathbf{x}\right) e^{\lambda_{1} t}
$$

and now (10.5.2) implies that $\mathbf{y}_{2}^{\prime}=A \mathbf{y}_{2}$.
To see that $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ are linearly independent, suppose $c_{1}$ and $c_{2}$ are constants such that

$$
\begin{equation*}
c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}=c_{1} \mathbf{x} e^{\lambda_{1} t}+c_{2}\left(\mathbf{u} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t}\right)=\mathbf{0} . \tag{10.5.4}
\end{equation*}
$$

We must show that $c_{1}=c_{2}=0$. Multiplying (10.5.4) by $e^{-\lambda_{1} t}$ shows that

$$
\begin{equation*}
c_{1} \mathbf{x}+c_{2}(\mathbf{u}+\mathbf{x} t)=\mathbf{0} \tag{10.5.5}
\end{equation*}
$$

By differentiating this with respect to $t$, we see that $c_{2} \mathbf{x}=\mathbf{0}$, which implies $c_{2}=0$, because $\mathbf{x} \neq \mathbf{0}$. Substituting $c_{2}=0$ into (10.5.5) yields $c_{1} \mathbf{x}=\mathbf{0}$, which implies that $c_{1}=0$, again because $\mathbf{x} \neq \mathbf{0}$

Example 10.5.2 Use Theorem 10.5.1 to find the general solution of the system

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rr}
11 & -25  \tag{10.5.6}\\
4 & -9
\end{array}\right] \mathbf{y}
$$

considered in Example 10.5.1.

Solution In Example 10.5.1 we saw that $\lambda_{1}=1$ is an eigenvalue of multiplicity 2 of the coefficient matrix $A$ in (10.5.6), and that all of the eigenvectors of $A$ are multiples of

$$
\mathbf{x}=\left[\begin{array}{l}
5 \\
2
\end{array}\right]
$$

Therefore

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
5 \\
2
\end{array}\right] e^{t}
$$

is a solution of (10.5.6). From Theorem 10.5.1, a second solution is given by $\mathbf{y}_{2}=\mathbf{u} e^{t}+\mathbf{x} t e^{t}$, where $(A-I) \mathbf{u}=\mathbf{x}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrll}
10 & -25 & \vdots & 5 \\
4 & -10 & \vdots & 2
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrcc}
1 & -\frac{5}{2} & \vdots & \frac{1}{2} \\
0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore the components of u must satisfy

$$
u_{1}-\frac{5}{2} u_{2}=\frac{1}{2},
$$

where $u_{2}$ is arbitrary. We choose $u_{2}=0$, so that $u_{1}=1 / 2$ and

$$
\mathbf{u}=\left[\begin{array}{c}
\frac{1}{2} \\
0
\end{array}\right]
$$

Thus,

$$
\mathbf{y}_{2}=\left[\begin{array}{l}
1 \\
0
\end{array}\right] \frac{e^{t}}{2}+\left[\begin{array}{l}
5 \\
2
\end{array}\right] t e^{t}
$$

Since $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ are linearly independent by Theorem 10.5.1, they form a fundamental set of solutions of (10.5.6). Therefore the general solution of (10.5.6) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
5 \\
2
\end{array}\right] e^{t}+c_{2}\left(\left[\begin{array}{l}
1 \\
0
\end{array}\right] \frac{e^{t}}{2}+\left[\begin{array}{l}
5 \\
2
\end{array}\right] t e^{t}\right)
$$

Note that choosing the arbitrary constant $u_{2}$ to be nonzero is equivalent to adding a scalar multiple of $\mathbf{y}_{1}$ to the second solution $\mathbf{y}_{2}$ (Exercise 33).

Example 10.5.3 Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
3 & 4 & -10  \tag{10.5.7}\\
2 & 1 & -2 \\
2 & 2 & -5
\end{array}\right] \mathbf{y}
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.5.7) is

$$
\left|\begin{array}{ccc}
3-\lambda & 4 & -10 \\
2 & 1-\lambda & -2 \\
2 & 2 & -5-\lambda
\end{array}\right|=-(\lambda-1)(\lambda+1)^{2} .
$$

Hence, the eigenvalues are $\lambda_{1}=1$ with multiplicity 1 and $\lambda_{2}=-1$ with multiplicity 2 .
Eigenvectors associated with $\lambda_{1}=1$ must satisfy $(A-I) \mathbf{x}=\mathbf{0}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrrll}
2 & 4 & -10 & \vdots & 0 \\
2 & 0 & -2 & \vdots & 0 \\
2 & 2 & -6 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -1 & \vdots & 0 \\
0 & 1 & -2 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=x_{3}$ and $x_{2}=2 x_{3}$, where $x_{3}$ is arbitrary. Choosing $x_{3}=1$ yields the eigenvector

$$
\mathbf{x}_{1}=\left[\begin{array}{l}
1 \\
2 \\
1
\end{array}\right]
$$

Therefore

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1 \\
2 \\
1
\end{array}\right] e^{t}
$$

is a solution of (10.5.7).
Eigenvectors associated with $\lambda_{2}=-1$ satisfy $(A+I) \mathbf{x}=\mathbf{0}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrrrr}
4 & 4 & -10 & \vdots & 0 \\
2 & 2 & -2 & \vdots & 0 \\
2 & 2 & -4 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{ccccc}
1 & 1 & 0 & \vdots & 0 \\
0 & 0 & 1 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{3}=0$ and $x_{1}=-x_{2}$, where $x_{2}$ is arbitrary. Choosing $x_{2}=1$ yields the eigenvector

$$
\mathbf{x}_{2}=\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right]
$$

so

$$
\mathbf{y}_{2}=\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] e^{-t}
$$

is a solution of (10.5.7).
Since all the eigenvectors of $A$ associated with $\lambda_{2}=-1$ are multiples of $\mathbf{x}_{2}$, we must now use Theorem 10.5.1 to find a third solution of (10.5.7) in the form

$$
\mathbf{y}_{3}=\mathbf{u} e^{-t}+\left[\begin{array}{r}
-1  \tag{10.5.8}\\
1 \\
0
\end{array}\right] t e^{-t}
$$

where $\mathbf{u}$ is a solution of $(A+I) \mathbf{u}=\mathbf{x}_{2}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrrrr}
4 & 4 & -10 & \vdots & -1 \\
2 & 2 & -2 & \vdots & 1 \\
2 & 2 & -4 & \vdots & 0
\end{array}\right],
$$

which is row equivalent to

$$
\left[\begin{array}{ccccc}
1 & 1 & 0 & \vdots & 1 \\
0 & 0 & 1 & \vdots & \frac{1}{2} \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $u_{3}=1 / 2$ and $u_{1}=1-u_{2}$, where $u_{2}$ is arbitrary. Choosing $u_{2}=0$ yields

$$
\mathbf{u}=\left[\begin{array}{c}
1 \\
0 \\
\frac{1}{2}
\end{array}\right]
$$

and substituting this into (10.5.8) yields the solution

$$
\mathbf{y}_{3}=\left[\begin{array}{l}
2 \\
0 \\
1
\end{array}\right] \frac{e^{-t}}{2}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] t e^{-t}
$$

of (10.5.7).
Since the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ at $t=0$ is

$$
\left|\begin{array}{rrr}
1 & -1 & 1 \\
2 & 1 & 0 \\
1 & 0 & \frac{1}{2}
\end{array}\right|=\frac{1}{2}
$$

$\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ is a fundamental set of solutions of (10.5.7). Therefore the general solution of (10.5.7) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
1 \\
2 \\
1
\end{array}\right] e^{t}+c_{2}\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] e^{-t}+c_{3}\left(\left[\begin{array}{l}
2 \\
0 \\
1
\end{array}\right] \frac{e^{-t}}{2}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] t e^{-t}\right)
$$

Theorem 10.5.2 Suppose the $n \times n$ matrix $A$ has an eigenvalue $\lambda_{1}$ of multiplicity $\geq 3$ and the associated eigenspace is one-dimensional; that is, all eigenvectors associated with $\lambda_{1}$ are scalar multiples of the eigenvector $\mathbf{x}$. Then there are infinitely many vectors $\mathbf{u}$ such that

$$
\begin{equation*}
\left(A-\lambda_{1} I\right) \mathbf{u}=\mathbf{x} \tag{10.5.9}
\end{equation*}
$$

and, if $\mathbf{u}$ is any such vector, there are infinitely many vectors $\mathbf{v}$ such that

$$
\begin{equation*}
\left(A-\lambda_{1} I\right) \mathbf{v}=\mathbf{u} \tag{10.5.10}
\end{equation*}
$$

If $\mathbf{u}$ satisfies (10.5.9) and $\mathbf{v}$ satisfies (10.5.10), then

$$
\begin{aligned}
& \mathbf{y}_{1}=\mathbf{x} e^{\lambda_{1} t}, \\
& \mathbf{y}_{2}=\mathbf{u} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t}, \text { and } \\
& \mathbf{y}_{3}=\mathbf{v} e^{\lambda_{1} t}+\mathbf{u} t e^{\lambda_{1} t}+\mathbf{x} \frac{t^{2} e^{\lambda_{1} t}}{2}
\end{aligned}
$$

are linearly independent solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$.
Again, it's beyond the scope of this book to prove that there are vectors $\mathbf{u}$ and $\mathbf{v}$ that satisfy (10.5.9) and (10.5.10). Theorem 10.5.1 implies that $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ are solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$. We leave the rest of the proof to you (Exercise 34).

Example 10.5.4 Use Theorem 10.5.2 to find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
1 & 1 & 1  \tag{10.5.11}\\
1 & 3 & -1 \\
0 & 2 & 2
\end{array}\right] \mathbf{y} .
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.5.11) is

$$
\left|\begin{array}{ccc}
1-\lambda & 1 & 1 \\
1 & 3-\lambda & -1 \\
0 & 2 & 2-\lambda
\end{array}\right|=-(\lambda-2)^{3} .
$$

Hence, $\lambda_{1}=2$ is an eigenvalue of multiplicity 3 . The associated eigenvectors satisfy $(A-2 I) \mathbf{x}=\mathbf{0}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrrrr}
-1 & 1 & 1 & \vdots & 0 \\
1 & 1 & -1 & \vdots & 0 \\
0 & 2 & 0 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -1 & \vdots & 0 \\
0 & 1 & 0 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Hence, $x_{1}=x_{3}$ and $x_{2}=0$, so the eigenvectors are all scalar multiples of

$$
\mathbf{x}_{1}=\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right]
$$

Therefore

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{2 t}
$$

is a solution of (10.5.11).
We now find a second solution of (10.5.11) in the form

$$
\mathbf{y}_{2}=\mathbf{u} e^{2 t}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] t e^{2 t}
$$

where $\mathbf{u}$ satisfies $(A-2 I) \mathbf{u}=\mathbf{x}_{1}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrrcr}
-1 & 1 & 1 & \vdots & 1 \\
1 & 1 & -1 & \vdots & 0 \\
0 & 2 & 0 & \vdots & 1
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -1 & \vdots & -\frac{1}{2} \\
0 & 1 & 0 & \vdots & \frac{1}{2} \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Letting $u_{3}=0$ yields $u_{1}=-1 / 2$ and $u_{2}=1 / 2$; hence,

$$
\mathbf{u}=\frac{1}{2}\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right]
$$

and

$$
\mathbf{y}_{2}=\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{e^{2 t}}{2}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] t e^{2 t}
$$

is a solution of (10.5.11).
We now find a third solution of (10.5.11) in the form

$$
\mathbf{y}_{3}=\mathbf{v} e^{2 t}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{t e^{2 t}}{2}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] \frac{t^{2} e^{2 t}}{2}
$$

where $\mathbf{v}$ satisfies $(A-2 I) \mathbf{v}=\mathbf{u}$. The augmented matrix of this system is

$$
\left[\begin{array}{rrrrr}
-1 & 1 & 1 & \vdots & -\frac{1}{2} \\
1 & 1 & -1 & \vdots & \frac{1}{2} \\
0 & 2 & 0 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -1 & \vdots & \frac{1}{2} \\
0 & 1 & 0 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Letting $v_{3}=0$ yields $v_{1}=1 / 2$ and $v_{2}=0$; hence,

$$
\mathbf{v}=\frac{1}{2}\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right]
$$

Therefore

$$
\mathbf{y}_{3}=\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right] \frac{e^{2 t}}{2}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{t e^{2 t}}{2}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] \frac{t^{2} e^{2 t}}{2}
$$

is a solution of (10.5.11). Since $\mathbf{y}_{1}, \mathbf{y}_{2}$, and $\mathbf{y}_{3}$ are linearly independent by Theorem 10.5.2, they form a fundamental set of solutions of (10.5.11). Therefore the general solution of (10.5.11) is

$$
\begin{aligned}
\mathbf{y}= & c_{1}\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{2 t}+c_{2}\left(\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{e^{2 t}}{2}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] t e^{2 t}\right) \\
& +c_{3}\left(\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right] \frac{e^{2 t}}{2}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{t e^{2 t}}{2}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] \frac{t^{2} e^{2 t}}{2}\right) .
\end{aligned}
$$

Theorem 10.5.3 Suppose the $n \times n$ matrix $A$ has an eigenvalue $\lambda_{1}$ of multiplicity $\geq 3$ and the associated eigenspace is two-dimensional; that is, all eigenvectors of $A$ associated with $\lambda_{1}$ are linear combinations of two linearly independent eigenvectors $\mathbf{x}_{1}$ and $\mathbf{x}_{2}$. Then there are constants $\alpha$ and $\beta$ (not both zero) such that if

$$
\begin{equation*}
\mathbf{x}_{3}=\alpha \mathbf{x}_{1}+\beta \mathbf{x}_{2} \tag{10.5.12}
\end{equation*}
$$

then there are infinitely many vectors $\mathbf{u}$ such that

$$
\begin{equation*}
\left(A-\lambda_{1} I\right) \mathbf{u}=\mathbf{x}_{3} \tag{10.5.13}
\end{equation*}
$$

If $\mathbf{u}$ satisfies (10.5.13), then

$$
\begin{align*}
& \mathbf{y}_{1}=\mathbf{x}_{1} e^{\lambda_{1} t} \\
& \mathbf{y}_{2}=\mathbf{x}_{2} e^{\lambda_{1} t}, \text { and } \\
& \mathbf{y}_{3}=\mathbf{u} e^{\lambda_{1} t}+\mathbf{x}_{3} t e^{\lambda_{1} t}, \tag{10.5.14}
\end{align*}
$$

are linearly independent solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$.
We omit the proof of this theorem.

Example 10.5.5 Use Theorem 10.5.3 to find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
0 & 0 & 1  \tag{10.5.15}\\
-1 & 1 & 1 \\
-1 & 0 & 2
\end{array}\right] \mathbf{y}
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.5.15) is

$$
\left|\begin{array}{ccc}
-\lambda & 0 & 1 \\
-1 & 1-\lambda & 1 \\
-1 & 0 & 2-\lambda
\end{array}\right|=-(\lambda-1)^{3} .
$$

Hence, $\lambda_{1}=1$ is an eigenvalue of multiplicity 3 . The associated eigenvectors satisfy $(A-I) \mathbf{x}=\mathbf{0}$. The augmented matrix of this system is

$$
\left[\begin{array}{ccccc}
-1 & 0 & 1 & \vdots & 0 \\
-1 & 0 & 1 & \vdots & 0 \\
-1 & 0 & 1 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -1 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right] .
$$

Hence, $x_{1}=x_{3}$ and $x_{2}$ is arbitrary, so the eigenvectors are of the form

$$
\mathbf{x}_{1}=\left[\begin{array}{l}
x_{3} \\
x_{2} \\
x_{3}
\end{array}\right]=x_{3}\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right]+x_{2}\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right] .
$$

Therefore the vectors

$$
\mathbf{x}_{1}=\left[\begin{array}{l}
1  \tag{10.5.16}\\
0 \\
1
\end{array}\right] \quad \text { and } \quad \mathbf{x}_{2}=\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right]
$$

form a basis for the eigenspace, and

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{t} \quad \text { and } \quad \mathbf{y}_{2}=\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right] e^{t}
$$

are linearly independent solutions of (10.5.15).
To find a third linearly independent solution of (10.5.15), we must find constants $\alpha$ and $\beta$ (not both zero) such that the system

$$
\begin{equation*}
(A-I) \mathbf{u}=\alpha \mathbf{x}_{1}+\beta \mathbf{x}_{2} \tag{10.5.17}
\end{equation*}
$$

has a solution $\mathbf{u}$. The augmented matrix of this system is

$$
\left[\begin{array}{ccccc}
-1 & 0 & 1 & \vdots & \alpha \\
-1 & 0 & 1 & \vdots & \beta \\
-1 & 0 & 1 & \vdots & \alpha
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrcr}
1 & 0 & -1 & \vdots & -\alpha  \tag{10.5.18}\\
0 & 0 & 0 & \vdots & \beta-\alpha \\
0 & 0 & 0 & \vdots & 0
\end{array}\right] .
$$

Therefore (10.5.17) has a solution if and only if $\beta=\alpha$, where $\alpha$ is arbitrary. If $\alpha=\beta=1$ then (10.5.12) and (10.5.16) yield

$$
\mathbf{x}_{3}=\mathbf{x}_{1}+\mathbf{x}_{2}=\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right]+\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right]=\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right]
$$

and the augmented matrix (10.5.18) becomes

$$
\left[\begin{array}{rrrlr}
1 & 0 & -1 & \vdots & -1 \\
0 & 0 & 0 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

This implies that $u_{1}=-1+u_{3}$, while $u_{2}$ and $u_{3}$ are arbitrary. Choosing $u_{2}=u_{3}=0$ yields

$$
\mathbf{u}=\left[\begin{array}{r}
-1 \\
0 \\
0
\end{array}\right]
$$

Therefore (10.5.14) implies that

$$
\mathbf{y}_{3}=\mathbf{u} e^{t}+\mathbf{x}_{3} t e^{t}=\left[\begin{array}{r}
-1 \\
0 \\
0
\end{array}\right] e^{t}+\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right] t e^{t}
$$

is a solution of (10.5.15). Since $\mathbf{y}_{1}, \mathbf{y}_{2}$, and $\mathbf{y}_{3}$ are linearly independent by Theorem 10.5.3, they form a fundamental set of solutions for (10.5.15). Therefore the general solution of (10.5.15) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{t}+c_{2}\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right] e^{t}+c_{3}\left(\left[\begin{array}{r}
-1 \\
0 \\
0
\end{array}\right] e^{t}+\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right] t e^{t}\right) .
$$

Geometric Properties of Solutions when $n=2$
We'll now consider the geometric properties of solutions of a $2 \times 2$ constant coefficient system

$$
\left[\begin{array}{l}
y_{1}^{\prime}  \tag{10.5.19}\\
y_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]\left[\begin{array}{l}
y_{1} \\
y_{2}
\end{array}\right]
$$

under the assumptions of this section; that is, when the matrix

$$
A=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]
$$

has a repeated eigenvalue $\lambda_{1}$ and the associated eigenspace is one-dimensional. In this case we know from Theorem 10.5.1 that the general solution of (10.5.19) is

$$
\begin{equation*}
\mathbf{y}=c_{1} \mathbf{x} e^{\lambda_{1} t}+c_{2}\left(\mathbf{u} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t}\right) \tag{10.5.20}
\end{equation*}
$$

where $\mathbf{x}$ is an eigenvector of $A$ and $\mathbf{u}$ is any one of the infinitely many solutions of

$$
\begin{equation*}
\left(A-\lambda_{1} I\right) \mathbf{u}=\mathbf{x} \tag{10.5.21}
\end{equation*}
$$

We assume that $\lambda_{1} \neq 0$.


Figure 10.5.1 Positive and negative half-planes

Let $L$ denote the line through the origin parallel to x . By a half-line of $L$ we mean either of the rays obtained by removing the origin from $L$. Eqn. (10.5.20) is a parametric equation of the half-line of $L$ in the direction of $\mathbf{x}$ if $c_{1}>0$, or of the half-line of $L$ in the direction of $-\mathbf{x}$ if $c_{1}<0$. The origin is the trajectory of the trivial solution $\mathbf{y} \equiv \mathbf{0}$.

Henceforth, we assume that $c_{2} \neq 0$. In this case, the trajectory of (10.5.20) can't intersect $L$, since every point of $L$ is on a trajectory obtained by setting $c_{2}=0$. Therefore the trajectory of (10.5.20) must lie entirely in one of the open half-planes bounded by $L$, but does not contain any point on $L$. Since the initial point $\left(y_{1}(0), y_{2}(0)\right)$ defined by $\mathbf{y}(0)=c_{1} \mathbf{x}_{1}+c_{2} \mathbf{u}$ is on the trajectory, we can determine which half-plane contains the trajectory from the sign of $c_{2}$, as shown in Figure 340. For convenience we'll call the half-plane where $c_{2}>0$ the positive half-plane. Similarly, the-half plane where $c_{2}<0$ is the negative half-plane. You should convince yourself (Exercise 35) that even though there are infinitely many vectors $\mathbf{u}$ that satisfy (10.5.21), they all define the same positive and negative half-planes. In the figures simply regard $\mathbf{u}$ as an arrow pointing to the positive half-plane, since wen't attempted to give $\mathbf{u}$ its proper length or direction in comparison with $\mathbf{x}$. For our purposes here, only the relative orientation of $\mathbf{x}$ and $\mathbf{u}$ is important; that is, whether the positive half-plane is to the right of an observer facing the direction of $\mathbf{x}$ (as in Figures 10.5.2 and 10.5.5), or to the left of the observer (as in Figures 10.5.3 and 10.5.4).

Multiplying (10.5.20) by $e^{-\lambda_{1} t}$ yields

$$
e^{-\lambda_{1} t} \mathbf{y}(t)=c_{1} \mathbf{x}+c_{2} \mathbf{u}+c_{2} t \mathbf{x}
$$

Since the last term on the right is dominant when $|t|$ is large, this provides the following information on the direction of $\mathbf{y}(t)$ :
(a) Along trajectories in the positive half-plane $\left(c_{2}>0\right)$, the direction of $\mathbf{y}(t)$ approaches the direction of $\mathbf{x}$ as $t \rightarrow \infty$ and the direction of $-\mathbf{x}$ as $t \rightarrow-\infty$.
(b) Along trajectories in the negative half-plane $\left(c_{2}<0\right)$, the direction of $\mathbf{y}(t)$ approaches the direction of $-\mathbf{x}$ as $t \rightarrow \infty$ and the direction of $\mathbf{x}$ as $t \rightarrow-\infty$.
Since

$$
\lim _{t \rightarrow \infty}\|\mathbf{y}(t)\|=\infty \quad \text { and } \quad \lim _{t \rightarrow-\infty} \mathbf{y}(t)=\mathbf{0} \quad \text { if } \quad \lambda_{1}>0
$$

or

$$
\lim _{t \rightarrow \infty}\|\mathbf{y}(t)\|=\infty \quad \text { and } \quad \lim _{t \rightarrow \infty} \mathbf{y}(t)=\mathbf{0} \quad \text { if } \quad \lambda_{1}<0,
$$

there are four possible patterns for the trajectories of (10.5.19), depending upon the signs of $c_{2}$ and $\lambda_{1}$. Figures 10.5.2-10.5.5 illustrate these patterns, and reveal the following principle:

If $\lambda_{1}$ and $c_{2}$ have the same sign then the direction of the traectory approaches the direction of -x as $\|\mathbf{y}\| \rightarrow 0$ and the direction of $\mathbf{x}$ as $\|\mathbf{y}\| \rightarrow \infty$. If $\lambda_{1}$ and $c_{2}$ have opposite signs then the direction of the trajectory approaches the direction of $\mathbf{x}$ as $\|\mathbf{y}\| \rightarrow 0$ and the direction of $-\mathbf{x}$ as $\|\mathbf{y}\| \rightarrow \infty$.


Figure 10.5.2 Positive eigenvalue; motion away from the origin


Figure 10.5.4 Negative eigenvalue; motion toward the origin


Figure 10.5.3 Positive eigenvalue; motion away from the origin


Figure 10.5.5 Negative eigenvalue; motion toward the origin

### 10.5 Exercises

In Exercises 1-12 find the general solution.

1. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}3 & 4 \\ -1 & 7\end{array}\right] \mathbf{y}$
2. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}0 & -1 \\ 1 & -2\end{array}\right] \mathbf{y}$
3. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-7 & 4 \\ -1 & -11\end{array}\right] \mathbf{y}$
4. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}3 & 1 \\ -1 & 1\end{array}\right] \mathbf{y}$
5. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}4 & 12 \\ -3 & -8\end{array}\right] \mathbf{y}$
6. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-10 & 9 \\ -4 & 2\end{array}\right] \mathbf{y}$
7. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-13 & 16 \\ -9 & 11\end{array}\right] \mathbf{y}$
8. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}0 & 2 & 1 \\ -4 & 6 & 1 \\ 0 & 4 & 2\end{array}\right] \mathbf{y}$
9. $\mathbf{y}^{\prime}=\frac{1}{3}\left[\begin{array}{rrr}1 & 1 & -3 \\ -4 & -4 & 3 \\ -2 & 1 & 0\end{array}\right] \mathbf{y}$
10. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-1 & 1 & -1 \\ -2 & 0 & 2 \\ -1 & 3 & -1\end{array}\right] \mathbf{y}$
11. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}4 & -2 & -2 \\ -2 & 3 & -1 \\ 2 & -1 & 3\end{array}\right] \mathbf{y}$
12. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}6 & -5 & 3 \\ 2 & -1 & 3 \\ 2 & 1 & 1\end{array}\right] \mathbf{y}$

In Exercises 13-23 solve the initial value problem.
13. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-11 & 8 \\ -2 & -3\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}6 \\ 2\end{array}\right]$
14. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}15 & -9 \\ 16 & -9\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}5 \\ 8\end{array}\right]$
15. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-3 & -4 \\ 1 & -7\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}2 \\ 3\end{array}\right]$
16. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-7 & 24 \\ -6 & 17\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}3 \\ 1\end{array}\right]$
17. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-7 & 3 \\ -3 & -1\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}0 \\ 2\end{array}\right]$
18. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-1 & 1 & 0 \\ 1 & -1 & -2 \\ -1 & -1 & -1\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}6 \\ 5 \\ -7\end{array}\right]$
19. $\mathbf{y}^{\prime}=\left[\begin{array}{lll}-2 & 2 & 1 \\ -2 & 2 & 1 \\ -3 & 3 & 2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}-6 \\ -2 \\ 0\end{array}\right]$
20. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-7 & -4 & 4 \\ -1 & 0 & 1 \\ -9 & -5 & 6\end{array}\right] \mathbf{y}, \quad \mathbf{y}(\mathbf{0})=\left[\begin{array}{r}-6 \\ 9 \\ -1\end{array}\right]$
21. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-1 & -4 & -1 \\ 3 & 6 & 1 \\ -3 & -2 & 3\end{array}\right] \mathbf{y}, \quad \mathbf{y}(\mathbf{0})=\left[\begin{array}{r}-2 \\ 1 \\ 3\end{array}\right]$
22. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}4 & -8 & -4 \\ -3 & -1 & -3 \\ 1 & -1 & 9\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}-4 \\ 1 \\ -3\end{array}\right]$
23. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-5 & -1 & 11 \\ -7 & 1 & 13 \\ -4 & 0 & 8\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}0 \\ 2 \\ 2\end{array}\right]$

The coefficient matrices in Exercises 24-32 have eigenvalues of multiplicity 3. Find the general solution.
24. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}5 & -1 & 1 \\ -1 & 9 & -3 \\ -2 & 2 & 4\end{array}\right] \mathbf{y}$
25. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}1 & 10 & -12 \\ 2 & 2 & 3 \\ 2 & -1 & 6\end{array}\right] \mathbf{y}$
26. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-6 & -4 & -4 \\ 2 & -1 & 1 \\ 2 & 3 & 1\end{array}\right] \mathbf{y}$
27. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}0 & 2 & -2 \\ -1 & 5 & -3 \\ 1 & 1 & 1\end{array}\right] \mathbf{y}$
28. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-2 & -12 & 10 \\ 2 & -24 & 11 \\ 2 & -24 & 8\end{array}\right] \mathbf{y}$
29. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-1 & -12 & 8 \\ 1 & -9 & 4 \\ 1 & -6 & 1\end{array}\right] \mathbf{y}$
30. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-4 & 0 & -1 \\ -1 & -3 & -1 \\ 1 & 0 & -2\end{array}\right] \mathbf{y}$
31. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-3 & -3 & 4 \\ 4 & 5 & -8 \\ 2 & 3 & -5\end{array}\right] \mathbf{y}$
32. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-3 & -1 & 0 \\ 1 & -1 & 0 \\ -1 & -1 & -2\end{array}\right] \mathbf{y}$
33. Under the assumptions of Theorem 10.5.1, suppose $\mathbf{u}$ and $\hat{\mathbf{u}}$ are vectors such that

$$
\left(A-\lambda_{1} I\right) \mathbf{u}=\mathbf{x} \quad \text { and } \quad\left(A-\lambda_{1} I\right) \hat{\mathbf{u}}=\mathbf{x}
$$

and let

$$
\mathbf{y}_{2}=\mathbf{u} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t} \quad \text { and } \quad \hat{\mathbf{y}}_{2}=\hat{\mathbf{u}} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t}
$$

Show that $\mathbf{y}_{2}-\hat{\mathbf{y}}_{2}$ is a scalar multiple of $\mathbf{y}_{1}=\mathbf{x} e^{\lambda_{1} t}$.
34. Under the assumptions of Theorem 10.5.2, let

$$
\begin{aligned}
& \mathbf{y}_{1}=\mathbf{x} e^{\lambda_{1} t}, \\
& \mathbf{y}_{2}=\mathbf{u} e^{\lambda_{1} t}+\mathbf{x} t e^{\lambda_{1} t}, \text { and } \\
& \mathbf{y}_{3}=\mathbf{v} e^{\lambda_{1} t}+\mathbf{u} t e^{\lambda_{1} t}+\mathbf{x} \frac{t^{2} e^{\lambda_{1} t}}{2}
\end{aligned}
$$

Complete the proof of Theorem 10.5 .2 by showing that $\mathbf{y}_{3}$ is a solution of $\mathbf{y}^{\prime}=A \mathbf{y}$ and that $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ is linearly independent.
35. Suppose the matrix

$$
A=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]
$$

has a repeated eigenvalue $\lambda_{1}$ and the associated eigenspace is one-dimensional. Let $\mathbf{x}$ be a $\lambda_{1}$ eigenvector of $A$. Show that if $\left(A-\lambda_{1} I\right) \mathbf{u}_{1}=\mathbf{x}$ and $\left(A-\lambda_{1} I\right) \mathbf{u}_{2}=\mathbf{x}$, then $\mathbf{u}_{2}-\mathbf{u}_{1}$ is parallel to $\mathbf{x}$. Conclude from this that all vectors $\mathbf{u}$ such that $\left(A-\lambda_{1} I\right) \mathbf{u}=\mathbf{x}$ define the same positive and negative half-planes with respect to the line $L$ through the origin parallel to $\mathbf{x}$.

In Exercises 36-45 plot trajectories of the given system.
36. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}-3 & -1 \\ 4 & 1\end{array}\right] \mathbf{y}$
37. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}2 & -1 \\ 1 & 0\end{array}\right] \mathbf{y}$
38. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}-1 & -3 \\ 3 & 5\end{array}\right] \mathrm{y}$
39. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-5 & 3 \\ -3 & 1\end{array}\right] \mathbf{y}$
40. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-2 & -3 \\ 3 & 4\end{array}\right] \mathbf{y}$
41. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-4 & -3 \\ 3 & 2\end{array}\right] \mathbf{y}$
42. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}0 & -1 \\ 1 & -2\end{array}\right] \mathbf{y}$
43. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}0 & 1 \\ -1 & 2\end{array}\right] \mathbf{y}$
44. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-2 & 1 \\ -1 & 0\end{array}\right] \mathbf{y}$
45. $\mathrm{C} / \mathrm{G} \mathbf{y}^{\prime}=\left[\begin{array}{ll}0 & -4 \\ 1 & -4\end{array}\right] \mathbf{y}$

### 10.6 CONSTANT COEFFICIENT HOMOGENEOUS SYSTEMS III

We now consider the system $\mathbf{y}^{\prime}=A \mathbf{y}$, where $A$ has a complex eigenvalue $\lambda=\alpha+i \beta$ with $\beta \neq 0$. We continue to assume that $A$ has real entries, so the characteristic polynomial of $A$ has real coefficients. This implies that $\bar{\lambda}=\alpha-i \beta$ is also an eigenvalue of $A$.

An eigenvector $\mathbf{x}$ of $A$ associated with $\lambda=\alpha+i \beta$ will have complex entries, so we'll write

$$
\mathbf{x}=\mathbf{u}+i \mathbf{v}
$$

where $\mathbf{u}$ and $\mathbf{v}$ have real entries; that is, $\mathbf{u}$ and $\mathbf{v}$ are the real and imaginary parts of $\mathbf{x}$. Since $A \mathbf{x}=\lambda \mathbf{x}$,

$$
\begin{equation*}
A(\mathbf{u}+i \mathbf{v})=(\alpha+i \beta)(\mathbf{u}+i \mathbf{v}) \tag{10.6.1}
\end{equation*}
$$

Taking complex conjugates here and recalling that $A$ has real entries yields

$$
A(\mathbf{u}-i \mathbf{v})=(\alpha-i \beta)(\mathbf{u}-i \mathbf{v})
$$

which shows that $\mathbf{x}=\mathbf{u}-i \mathbf{v}$ is an eigenvector associated with $\bar{\lambda}=\alpha-i \beta$. The complex conjugate eigenvalues $\lambda$ and $\bar{\lambda}$ can be separately associated with linearly independent solutions $\mathbf{y}^{\prime}=A \mathbf{y}$; however, we won't pursue this approach, since solutions obtained in this way turn out to be complex-valued. Instead, we'll obtain solutions of $\mathbf{y}^{\prime}=A y$ in the form

$$
\begin{equation*}
\mathbf{y}=f_{1} \mathbf{u}+f_{2} \mathbf{v} \tag{10.6.2}
\end{equation*}
$$

where $f_{1}$ and $f_{2}$ are real-valued scalar functions. The next theorem shows how to do this.
Theorem 10.6.1 Let $A$ be an $n \times n$ matrix with real entries. Let $\lambda=\alpha+i \beta(\beta \neq 0)$ be a complex eigenvalue of $A$ and let $\mathbf{x}=\mathbf{u}+i \mathbf{v}$ be an associated eigenvector, where $\mathbf{u}$ and $\mathbf{v}$ have real components. Then $\mathbf{u}$ and $\mathbf{v}$ are both nonzero and

$$
\mathbf{y}_{1}=e^{\alpha t}(\mathbf{u} \cos \beta t-\mathbf{v} \sin \beta t) \quad \text { and } \quad \mathbf{y}_{2}=e^{\alpha t}(\mathbf{u} \sin \beta t+\mathbf{v} \cos \beta t),
$$

which are the real and imaginary parts of

$$
\begin{equation*}
e^{\alpha t}(\cos \beta t+i \sin \beta t)(\mathbf{u}+i \mathbf{v}) \tag{10.6.3}
\end{equation*}
$$

are linearly independent solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$.
Proof A function of the form (10.6.2) is a solution of $\mathbf{y}^{\prime}=A \mathbf{y}$ if and only if

$$
\begin{equation*}
f_{1}^{\prime} \mathbf{u}+f_{2}^{\prime} \mathbf{v}=f_{1} A \mathbf{u}+f_{2} A \mathbf{v} \tag{10.6.4}
\end{equation*}
$$

Carrying out the multiplication indicated on the right side of (10.6.1) and collecting the real and imaginary parts of the result yields

$$
A(\mathbf{u}+i \mathbf{v})=(\alpha \mathbf{u}-\beta \mathbf{v})+i(\alpha \mathbf{v}+\beta \mathbf{u})
$$

Equating real and imaginary parts on the two sides of this equation yields

$$
\begin{aligned}
& A \mathbf{u}=\alpha \mathbf{u}-\beta \mathbf{v} \\
& A \mathbf{v}=\alpha \mathbf{v}+\beta \mathbf{u}
\end{aligned}
$$

We leave it to you (Exercise 25) to show from this that $\mathbf{u}$ and $\mathbf{v}$ are both nonzero. Substituting from these equations into (10.6.4) yields

$$
\begin{aligned}
f_{1}^{\prime} \mathbf{u}+f_{2}^{\prime} \mathbf{v} & =f_{1}(\alpha \mathbf{u}-\beta \mathbf{v})+f_{2}(\alpha \mathbf{v}+\beta \mathbf{u}) \\
& =\left(\alpha f_{1}+\beta f_{2}\right) \mathbf{u}+\left(-\beta f_{1}+\alpha f_{2}\right) \mathbf{v}
\end{aligned}
$$

This is true if

$$
\begin{aligned}
& f_{1}^{\prime}=\alpha f_{1}+\beta f_{2} \\
& f_{2}^{\prime}=-\beta f_{1}+\alpha f_{2},
\end{aligned} \quad \text { or, equivalently, } \quad \begin{aligned}
& f_{1}^{\prime}-\alpha f_{1}=\beta f_{2} \\
& f_{2}^{\prime}-\alpha f_{2}=-\beta f_{1}
\end{aligned}
$$

If we let $f_{1}=g_{1} e^{\alpha t}$ and $f_{2}=g_{2} e^{\alpha t}$, where $g_{1}$ and $g_{2}$ are to be determined, then the last two equations become

$$
\begin{aligned}
& g_{1}^{\prime}=\beta g_{2} \\
& g_{2}^{\prime}=-\beta g_{1}
\end{aligned}
$$

which implies that

$$
g_{1}^{\prime \prime}=\beta g_{2}^{\prime}=-\beta^{2} g_{1}
$$

so

$$
g_{1}^{\prime \prime}+\beta^{2} g_{1}=0
$$

The general solution of this equation is

$$
g_{1}=c_{1} \cos \beta t+c_{2} \sin \beta t
$$

Moreover, since $g_{2}=g_{1}^{\prime} / \beta$,

$$
g_{2}=-c_{1} \sin \beta t+c_{2} \cos \beta t
$$

Multiplying $g_{1}$ and $g_{2}$ by $e^{\alpha t}$ shows that

$$
\begin{aligned}
& f_{1}=e^{\alpha t}\left(c_{1} \cos \beta t+c_{2} \sin \beta t\right) \\
& f_{2}=e^{\alpha t}\left(-c_{1} \sin \beta t+c_{2} \cos \beta t\right)
\end{aligned}
$$

Substituting these into (10.6.2) shows that

$$
\begin{align*}
\mathbf{y} & =e^{\alpha t}\left[\left(c_{1} \cos \beta t+c_{2} \sin \beta t\right) \mathbf{u}+\left(-c_{1} \sin \beta t+c_{2} \cos \beta t\right) \mathbf{v}\right]  \tag{10.6.5}\\
& =c_{1} e^{\alpha t}(\mathbf{u} \cos \beta t-\mathbf{v} \sin \beta t)+c_{2} e^{\alpha t}(\mathbf{u} \sin \beta t+\mathbf{v} \cos \beta t)
\end{align*}
$$

is a solution of $\mathbf{y}^{\prime}=A \mathbf{y}$ for any choice of the constants $c_{1}$ and $c_{2}$. In particular, by first taking $c_{1}=1$ and $c_{2}=0$ and then taking $c_{1}=0$ and $c_{2}=1$, we see that $\mathbf{y}_{1}$ and $\mathbf{y}_{2}$ are solutions of $\mathbf{y}^{\prime}=A \mathbf{y}$. We leave it to you to verify that they are, respectively, the real and imaginary parts of (10.6.3) (Exercise 26), and that they are linearly independent (Exercise 27).

Example 10.6.1 Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
4 & -5  \tag{10.6.6}\\
5 & -2
\end{array}\right] \mathbf{y}
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.6.6) is

$$
\left|\begin{array}{cc}
4-\lambda & -5 \\
5 & -2-\lambda
\end{array}\right|=(\lambda-1)^{2}+16
$$

Hence, $\lambda=1+4 i$ is an eigenvalue of $A$. The associated eigenvectors satisfy $(A-(1+4 i) I) \mathbf{x}=\mathbf{0}$. The augmented matrix of this system is

$$
\left[\begin{array}{cccc}
3-4 i & -5 & \vdots & 0 \\
5 & -3-4 i & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{cccc}
1 & -\frac{3+4 i}{5} & \vdots & 0 \\
0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore $x_{1}=(3+4 i) x_{2} / 5$. Taking $x_{2}=5$ yields $x_{1}=3+4 i$, so

$$
\mathbf{x}=\left[\begin{array}{c}
3+4 i \\
5
\end{array}\right]
$$

is an eigenvector. The real and imaginary parts of

$$
e^{t}(\cos 4 t+i \sin 4 t)\left[\begin{array}{c}
3+4 i \\
5
\end{array}\right]
$$

are

$$
\mathbf{y}_{1}=e^{t}\left[\begin{array}{c}
3 \cos 4 t-4 \sin 4 t \\
5 \cos 4 t
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=e^{t}\left[\begin{array}{c}
3 \sin 4 t+4 \cos 4 t \\
5 \sin 4 t
\end{array}\right],
$$

which are linearly independent solutions of (10.6.6). The general solution of (10.6.6) is

$$
\mathbf{y}=c_{1} e^{t}\left[\begin{array}{c}
3 \cos 4 t-4 \sin 4 t \\
5 \cos 4 t
\end{array}\right]+c_{2} e^{t}\left[\begin{array}{c}
3 \sin 4 t+4 \cos 4 t \\
5 \sin 4 t
\end{array}\right] .
$$

Example 10.6.2 Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rr}
-14 & 39  \tag{10.6.7}\\
-6 & 16
\end{array}\right] \mathbf{y} .
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.6.7) is

$$
\left|\begin{array}{cc}
-14-\lambda & 39 \\
-6 & 16-\lambda
\end{array}\right|=(\lambda-1)^{2}+9
$$

Hence, $\lambda=1+3 i$ is an eigenvalue of $A$. The associated eigenvectors satisfy $(A-(1+3 i) I) \mathbf{x}=\mathbf{0}$. The augmented augmented matrix of this system is

$$
\left[\begin{array}{cccc}
-15-3 i & 39 & \vdots & 0 \\
-6 & 15-3 i & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{cccc}
1 & \frac{-5+i}{2} & \vdots & 0 \\
0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore $x_{1}=(5-i) / 2$. Taking $x_{2}=2$ yields $x_{1}=5-i$, so

$$
\mathbf{x}=\left[\begin{array}{c}
5-i \\
2
\end{array}\right]
$$

is an eigenvector. The real and imaginary parts of

$$
e^{t}(\cos 3 t+i \sin 3 t)\left[\begin{array}{c}
5-i \\
2
\end{array}\right]
$$

are

$$
\mathbf{y}_{1}=e^{t}\left[\begin{array}{c}
\sin 3 t+5 \cos 3 t \\
2 \cos 3 t
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=e^{t}\left[\begin{array}{c}
-\cos 3 t+5 \sin 3 t \\
2 \sin 3 t
\end{array}\right]
$$

which are linearly independent solutions of (10.6.7). The general solution of (10.6.7) is

$$
\mathbf{y}=c_{1} e^{t}\left[\begin{array}{c}
\sin 3 t+5 \cos 3 t \\
2 \cos 3 t
\end{array}\right]+c_{2} e^{t}\left[\begin{array}{c}
-\cos 3 t+5 \sin 3 t \\
2 \sin 3 t
\end{array}\right] .
$$

Example 10.6.3 Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
-5 & 5 & 4  \tag{10.6.8}\\
-8 & 7 & 6 \\
1 & 0 & 0
\end{array}\right] \mathbf{y} .
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.6.8) is

$$
\left|\begin{array}{ccc}
-5-\lambda & 5 & 4 \\
-8 & 7-\lambda & 6 \\
1 & 0 & -\lambda
\end{array}\right|=-(\lambda-2)\left(\lambda^{2}+1\right)
$$

Hence, the eigenvalues of $A$ are $\lambda_{1}=2, \lambda_{2}=i$, and $\lambda_{3}=-i$. The augmented matrix of $(A-2 I) \mathbf{x}=\mathbf{0}$ is

$$
\left[\begin{array}{rrrrr}
-7 & 5 & 4 & \vdots & 0 \\
-8 & 5 & 6 & \vdots & 0 \\
1 & 0 & -2 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -2 & \vdots & 0 \\
0 & 1 & -2 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore $x_{1}=x_{2}=2 x_{3}$. Taking $x_{3}=1$ yields

$$
\mathbf{x}_{1}=\left[\begin{array}{l}
2 \\
2 \\
1
\end{array}\right]
$$

so

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
2 \\
2 \\
1
\end{array}\right] e^{2 t}
$$

is a solution of (10.6.8).
The augmented matrix of $(A-i I) \mathbf{x}=\mathbf{0}$ is

$$
\left[\begin{array}{ccccc}
-5-i & 5 & 4 & \vdots & 0 \\
-8 & 7-i & 6 & \vdots & 0 \\
1 & 0 & -i & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{ccccc}
1 & 0 & -i & \vdots & 0 \\
0 & 1 & 1-i & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore $x_{1}=i x_{3}$ and $x_{2}=-(1-i) x_{3}$. Taking $x_{3}=1$ yields the eigenvector

$$
\mathbf{x}_{2}=\left[\begin{array}{c}
i \\
-1+i \\
1
\end{array}\right]
$$

The real and imaginary parts of

$$
(\cos t+i \sin t)\left[\begin{array}{c}
i \\
-1+i \\
1
\end{array}\right]
$$

are

$$
\mathbf{y}_{2}=\left[\begin{array}{c}
-\sin t \\
-\cos t-\sin t \\
\cos t
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{3}=\left[\begin{array}{c}
\cos t \\
\cos t-\sin t \\
\sin t
\end{array}\right],
$$

which are solutions of (10.6.8). Since the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ at $t=0$ is

$$
\left|\begin{array}{rrr}
2 & 0 & 1 \\
2 & -1 & 1 \\
1 & 1 & 0
\end{array}\right|=1
$$

$\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ is a fundamental set of solutions of (10.6.8). The general solution of (10.6.8) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{l}
2 \\
2 \\
1
\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{c}
-\sin t \\
-\cos t-\sin t \\
\cos t
\end{array}\right]+c_{3}\left[\begin{array}{c}
\cos t \\
\cos t-\sin t \\
\sin t
\end{array}\right]
$$

Example 10.6.4 Find the general solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
1 & -1 & -2  \tag{10.6.9}\\
1 & 3 & 2 \\
1 & -1 & 2
\end{array}\right] \mathbf{y}
$$

Solution The characteristic polynomial of the coefficient matrix $A$ in (10.6.9) is

$$
\left|\begin{array}{ccc}
1-\lambda & -1 & -2 \\
1 & 3-\lambda & 2 \\
1 & -1 & 2-\lambda
\end{array}\right|=-(\lambda-2)\left((\lambda-2)^{2}+4\right) .
$$

Hence, the eigenvalues of $A$ are $\lambda_{1}=2, \lambda_{2}=2+2 i$, and $\lambda_{3}=2-2 i$. The augmented matrix of $(A-2 I) \mathbf{x}=\mathbf{0}$ is

$$
\left[\begin{array}{rrrcc}
-1 & -1 & -2 & \vdots & 0 \\
1 & 1 & 2 & \vdots & 0 \\
1 & -1 & 0 & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{ccccc}
1 & 0 & 1 & \vdots & 0 \\
0 & 1 & 1 & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore $x_{1}=x_{2}=-x_{3}$. Taking $x_{3}=1$ yields

$$
\mathbf{x}_{1}=\left[\begin{array}{r}
-1 \\
-1 \\
1
\end{array}\right]
$$

so

$$
\mathbf{y}_{1}=\left[\begin{array}{r}
-1 \\
-1 \\
1
\end{array}\right] e^{2 t}
$$

is a solution of (10.6.9).
The augmented matrix of $(A-(2+2 i) I) \mathbf{x}=\mathbf{0}$ is

$$
\left[\begin{array}{ccrcc}
-1-2 i & -1 & -2 & \vdots & 0 \\
1 & 1-2 i & 2 & \vdots & 0 \\
1 & -1 & -2 i & \vdots & 0
\end{array}\right]
$$

which is row equivalent to

$$
\left[\begin{array}{rrrrr}
1 & 0 & -i & \vdots & 0 \\
0 & 1 & i & \vdots & 0 \\
0 & 0 & 0 & \vdots & 0
\end{array}\right]
$$

Therefore $x_{1}=i x_{3}$ and $x_{2}=-i x_{3}$. Taking $x_{3}=1$ yields the eigenvector

$$
\mathbf{x}_{2}=\left[\begin{array}{r}
i \\
-i \\
1
\end{array}\right]
$$

The real and imaginary parts of

$$
e^{2 t}(\cos 2 t+i \sin 2 t)\left[\begin{array}{r}
i \\
-i \\
1
\end{array}\right]
$$

are

$$
\mathbf{y}_{2}=e^{2 t}\left[\begin{array}{r}
-\sin 2 t \\
\sin 2 t \\
\cos 2 t
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=e^{2 t}\left[\begin{array}{r}
\cos 2 t \\
-\cos 2 t \\
\sin 2 t
\end{array}\right]
$$

which are solutions of (10.6.9). Since the Wronskian of $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ at $t=0$ is

$$
\left|\begin{array}{rrr}
-1 & 0 & 1 \\
-1 & 0 & -1 \\
1 & 1 & 0
\end{array}\right|=-2,
$$

$\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \mathbf{y}_{3}\right\}$ is a fundamental set of solutions of (10.6.9). The general solution of (10.6.9) is

$$
\mathbf{y}=c_{1}\left[\begin{array}{r}
-1 \\
-1 \\
1
\end{array}\right] e^{2 t}+c_{2} e^{2 t}\left[\begin{array}{r}
-\sin 2 t \\
\sin 2 t \\
\cos 2 t
\end{array}\right]+c_{3} e^{2 t}\left[\begin{array}{r}
\cos 2 t \\
-\cos 2 t \\
\sin 2 t
\end{array}\right]
$$

Geometric Properties of Solutions when $n=2$
We'll now consider the geometric properties of solutions of a $2 \times 2$ constant coefficient system

$$
\left[\begin{array}{l}
y_{1}^{\prime}  \tag{10.6.10}\\
y_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]\left[\begin{array}{l}
y_{1} \\
y_{2}
\end{array}\right]
$$

under the assumptions of this section; that is, when the matrix

$$
A=\left[\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right]
$$

has a complex eigenvalue $\lambda=\alpha+i \beta(\beta \neq 0)$ and $\mathbf{x}=\mathbf{u}+i \mathbf{v}$ is an associated eigenvector, where $\mathbf{u}$ and $\mathbf{v}$ have real components. To describe the trajectories accurately it's necessary to introduce a new rectangular coordinate system in the $y_{1}-y_{2}$ plane. This raises a point that hasn't come up before: It is always possible to choose $\mathbf{x}$ so that $(\mathbf{u}, \mathbf{v})=0$. A special effort is required to do this, since not every eigenvector has this property. However, if we know an eigenvector that doesn't, we can multiply it by a suitable complex constant to obtain one that does. To see this, note that if $\mathbf{x}$ is a $\lambda$-eigenvector of $A$ and $k$ is an arbitrary real number, then

$$
\mathbf{x}_{1}=(1+i k) \mathbf{x}=(1+i k)(\mathbf{u}+i \mathbf{v})=(\mathbf{u}-k \mathbf{v})+i(\mathbf{v}+k \mathbf{u})
$$

is also a $\lambda$-eigenvector of $A$, since

$$
A \mathbf{x}_{1}=A((1+i k) \mathbf{x})=(1+i k) A \mathbf{x}=(1+i k) \lambda \mathbf{x}=\lambda((1+i k) \mathbf{x})=\lambda \mathbf{x}_{1} .
$$

The real and imaginary parts of $\mathbf{x}_{1}$ are

$$
\begin{equation*}
\mathbf{u}_{1}=\mathbf{u}-k \mathbf{v} \quad \text { and } \quad \mathbf{v}_{1}=\mathbf{v}+k \mathbf{u} \tag{10.6.11}
\end{equation*}
$$

so

$$
\left(\mathbf{u}_{1}, \mathbf{v}_{1}\right)=(\mathbf{u}-k \mathbf{v}, \mathbf{v}+k \mathbf{u})=-\left[(\mathbf{u}, \mathbf{v}) k^{2}+\left(\|\mathbf{v}\|^{2}-\|\mathbf{u}\|^{2}\right) k-(\mathbf{u}, \mathbf{v})\right] .
$$

Therefore $\left(\mathbf{u}_{1}, \mathbf{v}_{1}\right)=0$ if

$$
\begin{equation*}
(\mathbf{u}, \mathbf{v}) k^{2}+\left(\|\mathbf{v}\|^{2}-\|\mathbf{u}\|^{2}\right) k-(\mathbf{u}, \mathbf{v})=0 \tag{10.6.12}
\end{equation*}
$$

If $(\mathbf{u}, \mathbf{v}) \neq 0$ we can use the quadratic formula to find two real values of $k$ such that $\left(\mathbf{u}_{1}, \mathbf{v}_{1}\right)=0$ (Exercise 28).

Example 10.6.5 In Example 10.6.1 we found the eigenvector

$$
\mathbf{x}=\left[\begin{array}{c}
3+4 i \\
5
\end{array}\right]=\left[\begin{array}{l}
3 \\
5
\end{array}\right]+i\left[\begin{array}{l}
4 \\
0
\end{array}\right]
$$

for the matrix of the system (10.6.6). Here $\mathbf{u}=\left[\begin{array}{l}3 \\ 5\end{array}\right]$ and $\mathbf{v}=\left[\begin{array}{l}4 \\ 0\end{array}\right]$ are not orthogonal, since $(\mathbf{u}, \mathbf{v})=12$. Since $\|\mathbf{v}\|^{2}-\|\mathbf{u}\|^{2}=-18$, (10.6.12) is equivalent to

$$
2 k^{2}-3 k-2=0
$$

The zeros of this equation are $k_{1}=2$ and $k_{2}=-1 / 2$. Letting $k=2$ in (10.6.11) yields

$$
\mathbf{u}_{1}=\mathbf{u}-2 \mathbf{v}=\left[\begin{array}{r}
-5 \\
5
\end{array}\right] \quad \text { and } \quad \mathbf{v}_{1}=\mathbf{v}+2 \mathbf{u}=\left[\begin{array}{l}
10 \\
10
\end{array}\right]
$$

and $\left(\mathbf{u}_{1}, \mathbf{v}_{1}\right)=0$. Letting $k=-1 / 2$ in (10.6.11) yields

$$
\mathbf{u}_{1}=\mathbf{u}+\frac{\mathbf{v}}{2}=\left[\begin{array}{l}
5 \\
5
\end{array}\right] \quad \text { and } \quad \mathbf{v}_{1}=\mathbf{v}-\frac{\mathbf{u}}{2}=\frac{1}{2}\left[\begin{array}{r}
-5 \\
5
\end{array}\right]
$$

and again $\left(\mathbf{u}_{1}, \mathbf{v}_{1}\right)=0$.
(The numbers don't always work out as nicely as in this example. You'll need a calculator or computer to do Exercises 29-40.)

Henceforth, we'll assume that $(\mathbf{u}, \mathbf{v})=0$. Let $\mathbf{U}$ and $\mathbf{V}$ be unit vectors in the directions of $\mathbf{u}$ and $\mathbf{v}$, respectively; that is, $\mathbf{U}=\mathbf{u} /\|\mathbf{u}\|$ and $\mathbf{V}=\mathbf{v} /\|\mathbf{v}\|$. The new rectangular coordinate system will have the same origin as the $y_{1}-y_{2}$ system. The coordinates of a point in this system will be denoted by $\left(z_{1}, z_{2}\right)$, where $z_{1}$ and $z_{2}$ are the displacements in the directions of $\mathbf{U}$ and $\mathbf{V}$, respectively.

From (10.6.5), the solutions of (10.6.10) are given by

$$
\begin{equation*}
\mathbf{y}=e^{\alpha t}\left[\left(c_{1} \cos \beta t+c_{2} \sin \beta t\right) \mathbf{u}+\left(-c_{1} \sin \beta t+c_{2} \cos \beta t\right) \mathbf{v}\right] . \tag{10.6.13}
\end{equation*}
$$

For convenience, let's call the curve traversed by $e^{-\alpha t} \mathbf{y}(t)$ a shadow trajectory of (10.6.10). Multiplying (10.6.13) by $e^{-\alpha t}$ yields

$$
e^{-\alpha t} \mathbf{y}(t)=z_{1}(t) \mathbf{U}+z_{2}(t) \mathbf{V},
$$

where

$$
\begin{aligned}
& z_{1}(t)=\|\mathbf{u}\|\left(c_{1} \cos \beta t+c_{2} \sin \beta t\right) \\
& z_{2}(t)=\|\mathbf{v}\|\left(-c_{1} \sin \beta t+c_{2} \cos \beta t\right) .
\end{aligned}
$$

Therefore

$$
\frac{\left(z_{1}(t)\right)^{2}}{\|\mathbf{u}\|^{2}}+\frac{\left(z_{2}(t)\right)^{2}}{\|\mathbf{v}\|^{2}}=c_{1}^{2}+c_{2}^{2}
$$

(verify!), which means that the shadow trajectories of (10.6.10) are ellipses centered at the origin, with axes of symmetry parallel to $\mathbf{U}$ and $\mathbf{V}$. Since

$$
z_{1}^{\prime}=\frac{\beta\|\mathbf{u}\|}{\|\mathbf{v}\|} z_{2} \quad \text { and } \quad z_{2}^{\prime}=-\frac{\beta\|\mathbf{v}\|}{\|\mathbf{u}\|} z_{1},
$$

the vector from the origin to a point on the shadow ellipse rotates in the same direction that $\mathbf{V}$ would have to be rotated by $\pi / 2$ radians to bring it into coincidence with $\mathbf{U}$ (Figures 10.6 .1 and 10.6.2).

If $\alpha=0$, then any trajectory of (10.6.10) is a shadow trajectory of (10.6.10); therefore, if $\lambda$ is purely imaginary, then the trajectories of (10.6.10) are ellipses traversed periodically as indicated in Figures 10.6.1 and 10.6.2.

If $\alpha>0$, then

$$
\lim _{t \rightarrow \infty}\|\mathbf{y}(t)\|=\infty \quad \text { and } \quad \lim _{t \rightarrow-\infty} \mathbf{y}(t)=0
$$

so the trajectory spirals away from the origin as $t$ varies from $-\infty$ to $\infty$. The direction of the spiral depends upon the relative orientation of $\mathbf{U}$ and $\mathbf{V}$, as shown in Figures 10.6.3 and 10.6.4.


Figure 10.6.1 Shadow trajectories traversed clockwise


Figure 10.6.2 Shadow trajectories traversed counterclockwise

If $\alpha<0$, then

$$
\lim _{t \rightarrow-\infty}\|\mathbf{y}(t)\|=\infty \quad \text { and } \quad \lim _{t \rightarrow \infty} \mathbf{y}(t)=0
$$

so the trajectory spirals toward the origin as $t$ varies from $-\infty$ to $\infty$. Again, the direction of the spiral depends upon the relative orientation of $\mathbf{U}$ and $\mathbf{V}$, as shown in Figures 10.6.5 and 10.6.6.


Figure 10.6.3 $\alpha>0$; shadow trajectory spiraling outward


Figure 10.6.5 $\alpha<0$; shadow trajectory spiraling inward


Figure 10.6.4 $\alpha>0$; shadow trajectory spiraling outward


Figure $10.6 .6 \alpha<0$; shadow trajectory spiraling inward

### 10.6 Exercises

In Exercises 1-16 find the general solution.

1. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}-1 & 2 \\ -5 & 5\end{array}\right] \mathbf{y}$
2. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}-11 & 4 \\ -26 & 9\end{array}\right] \mathbf{y}$
3. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}1 & 2 \\ -4 & 5\end{array}\right] \mathbf{y}$
4. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}5 & -6 \\ 3 & -1\end{array}\right] \mathbf{y}$
5. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & -3 & 1 \\ 0 & 2 & 2 \\ 5 & 1 & 1\end{array}\right] \mathbf{y}$
6. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-3 & 3 & 1 \\ 1 & -5 & -3 \\ -3 & 7 & 3\end{array}\right] \mathbf{y}$
7. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}2 & 1 & -1 \\ 0 & 1 & 1 \\ 1 & 0 & 1\end{array}\right] \mathbf{y}$
8. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-3 & 1 & -3 \\ 4 & -1 & 2 \\ 4 & -2 & 3\end{array}\right] \mathbf{y}$
9. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}5 & -4 \\ 10 & 1\end{array}\right] \mathbf{y}$
10. $\quad \mathbf{y}^{\prime}=\frac{1}{3}\left[\begin{array}{rr}7 & -5 \\ 2 & 5\end{array}\right] \mathbf{y}$
11. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}3 & 2 \\ -5 & 1\end{array}\right] \mathbf{y}$
12. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}34 & 52 \\ -20 & -30\end{array}\right] \mathbf{y}$
13. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}1 & 1 & 2 \\ 1 & 0 & -1 \\ -1 & -2 & -1\end{array}\right] \mathbf{y}$
14. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & -4 & -2 \\ -5 & 7 & -8 \\ -10 & 13 & -8\end{array}\right] \mathbf{y}$
15. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}6 & 0 & -3 \\ -3 & 3 & 3 \\ 1 & -2 & 6\end{array}\right] \mathbf{y}^{\prime}$
16. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}1 & 2 & -2 \\ 0 & 2 & -1 \\ 1 & 0 & 0\end{array}\right] \mathbf{y}^{\prime}$

In Exercises 17-24 solve the initial value problem.
17. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}4 & -6 \\ 3 & -2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}5 \\ 2\end{array}\right]$
18. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}7 & 15 \\ -3 & 1\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}5 \\ 1\end{array}\right]$
19. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}7 & -15 \\ 3 & -5\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}17 \\ 7\end{array}\right]$
20. $\quad \mathbf{y}^{\prime}=\frac{1}{6}\left[\begin{array}{rr}4 & -2 \\ 5 & 2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}1 \\ -1\end{array}\right]$
21. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}5 & 2 & -1 \\ -3 & 2 & 2 \\ 1 & 3 & 2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}4 \\ 0 \\ 6\end{array}\right]$
22. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}4 & 4 & 0 \\ 8 & 10 & -20 \\ 2 & 3 & -2\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}8 \\ 6 \\ 5\end{array}\right]$
23. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}1 & 15 & -15 \\ -6 & 18 & -22 \\ -3 & 11 & -15\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{l}15 \\ 17 \\ 10\end{array}\right]$
24. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}4 & -4 & 4 \\ -10 & 3 & 15 \\ 2 & -3 & 1\end{array}\right] \mathbf{y}, \quad \mathbf{y}(0)=\left[\begin{array}{r}16 \\ 14 \\ 6\end{array}\right]$
25. Suppose an $n \times n$ matrix $A$ with real entries has a complex eigenvalue $\lambda=\alpha+i \beta(\beta \neq 0)$ with associated eigenvector $\mathbf{x}=\mathbf{u}+i \mathbf{v}$, where $\mathbf{u}$ and $\mathbf{v}$ have real components. Show that $\mathbf{u}$ and $\mathbf{v}$ are both nonzero.
26. Verify that

$$
\mathbf{y}_{1}=e^{\alpha t}(\mathbf{u} \cos \beta t-\mathbf{v} \sin \beta t) \quad \text { and } \quad \mathbf{y}_{2}=e^{\alpha t}(\mathbf{u} \sin \beta t+\mathbf{v} \cos \beta t)
$$

are the real and imaginary parts of

$$
e^{\alpha t}(\cos \beta t+i \sin \beta t)(\mathbf{u}+i \mathbf{v})
$$

27. Show that if the vectors $\mathbf{u}$ and $\mathbf{v}$ are not both $\mathbf{0}$ and $\beta \neq 0$ then the vector functions

$$
\mathbf{y}_{1}=e^{\alpha t}(\mathbf{u} \cos \beta t-\mathbf{v} \sin \beta t) \quad \text { and } \quad \mathbf{y}_{2}=e^{\alpha t}(\mathbf{u} \sin \beta t+\mathbf{v} \cos \beta t)
$$

are linearly independent on every interval. Hint: There are two cases to consider: (i) $\{\mathbf{u}, \mathbf{v}\}$ linearly independent, and (ii) $\{\mathbf{u}, \mathbf{v}\}$ linearly dependent. In either case, exploit the the linear independence of $\{\cos \beta t, \sin \beta t\}$ on every interval.
28. Suppose $\mathbf{u}=\left[\begin{array}{l}u_{1} \\ u_{2}\end{array}\right]$ and $\mathbf{v}=\left[\begin{array}{l}v_{1} \\ v_{2}\end{array}\right]$ are not orthogonal; that is, $(\mathbf{u}, \mathbf{v}) \neq 0$.
(a) Show that the quadratic equation

$$
(\mathbf{u}, \mathbf{v}) k^{2}+\left(\|\mathbf{v}\|^{2}-\|\mathbf{u}\|^{2}\right) k-(\mathbf{u}, \mathbf{v})=0
$$

has a positive root $k_{1}$ and a negative root $k_{2}=-1 / k_{1}$.
(b) Let $\mathbf{u}_{1}^{(1)}=\mathbf{u}-k_{1} \mathbf{v}, \mathbf{v}_{1}^{(1)}=\mathbf{v}+k_{1} \mathbf{u}, \mathbf{u}_{1}^{(2)}=\mathbf{u}-k_{2} \mathbf{v}$, and $\mathbf{v}_{1}^{(2)}=\mathbf{v}+k_{2} \mathbf{u}$, so that $\left(\mathbf{u}_{1}^{(1)}, \mathbf{v}_{1}^{(1)}\right)=\left(\mathbf{u}_{1}^{(2)}, \mathbf{v}_{1}^{(2)}\right)=0$, from the discussion given above. Show that

$$
\mathbf{u}_{1}^{(2)}=\frac{\mathbf{v}_{1}^{(1)}}{k_{1}} \quad \text { and } \quad \mathbf{v}_{1}^{(2)}=-\frac{\mathbf{u}_{1}^{(1)}}{k_{1}}
$$

(c) Let $\mathbf{U}_{1}, \mathbf{V}_{1}, \mathbf{U}_{2}$, and $\mathbf{V}_{2}$ be unit vectors in the directions of $\mathbf{u}_{1}^{(1)}, \mathbf{v}_{1}^{(1)}, \mathbf{u}_{1}^{(2)}$, and $\mathbf{v}_{1}^{(2)}$, respectively. Conclude from (a) that $\mathbf{U}_{2}=\mathbf{V}_{1}$ and $\mathbf{V}_{2}=-\mathbf{U}_{1}$, and that therefore the counterclockwise angles from $\mathbf{U}_{1}$ to $\mathbf{V}_{1}$ and from $\mathbf{U}_{2}$ to $\mathbf{V}_{2}$ are both $\pi / 2$ or both $-\pi / 2$.

In Exercises 29-32 find vectors $\mathbf{U}$ and $\mathbf{V}$ parallel to the axes of symmetry of the trajectories, and plot some typical trajectories.
29. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}3 & -5 \\ 5 & -3\end{array}\right] \mathbf{y}$
30. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-15 & 10 \\ -25 & 15\end{array}\right] \mathbf{y}$
31. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-4 & 8 \\ -4 & 4\end{array}\right] \mathbf{y}$
32. $\overline{C / G} \mathbf{y}^{\prime}=\left[\begin{array}{rr}-3 & -15 \\ 3 & 3\end{array}\right] \mathbf{y}$

In Exercises 33-40 find vectors $\mathbf{U}$ and $\mathbf{V}$ parallel to the axes of symmetry of the shadow trajectories, and plot a typical trajectory.
33. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-5 & 6 \\ -12 & 7\end{array}\right] \mathbf{y}$
34. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}5 & -12 \\ 6 & -7\end{array}\right] \mathbf{y}$
35. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}4 & -5 \\ 9 & -2\end{array}\right] \mathbf{y}$
36. $\mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{ll}-4 & 9 \\ -5 & 2\end{array}\right] \mathbf{y}$
37. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-1 & 10 \\ -10 & -1\end{array}\right] \mathbf{y}$
38. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-1 & -5 \\ 20 & -1\end{array}\right] \mathbf{y}$

40. $\quad \mathrm{C} / \mathrm{G} \quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-7 & 6 \\ -12 & 5\end{array}\right] \mathbf{y}$

### 10.7 VARIATION OF PARAMETERS FOR NONHOMOGENEOUS LINEAR SYSTEMS

We now consider the nonhomogeneous linear system

$$
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t)
$$

where $A$ is an $n \times n$ matrix function and $\mathbf{f}$ is an $n$-vector forcing function. Associated with this system is the complementary system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$.

The next theorem is analogous to Theorems 5.3.2 and ??. It shows how to find the general solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t)$ if we know a particular solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t)$ and a fundamental set of solutions of the complementary system. We leave the proof as an exercise (Exercise 21).

Theorem 10.7.1 Suppose the $n \times n$ matrix function $A$ and the $n$-vector function $\mathbf{f}$ are continuous on $(a, b)$. Let $\mathbf{y}_{p}$ be a particular solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t)$ on $(a, b)$, and let $\left\{\mathbf{y}_{1}, \mathbf{y}_{2}, \ldots, \mathbf{y}_{n}\right\}$ be a fundamental set of solutions of the complementary equation $\mathbf{y}^{\prime}=A(t) \mathbf{y}$ on $(a, b)$. Then $\mathbf{y}$ is a solution of $\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t)$ on $(a, b)$ if and only if

$$
\mathbf{y}=\mathbf{y}_{p}+c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}+\cdots+c_{n} \mathbf{y}_{n}
$$

where $c_{1}, c_{2}, \ldots, c_{n}$ are constants.
Finding a Particular Solution of a Nonhomogeneous System
We now discuss an extension of the method of variation of parameters to linear nonhomogeneous systems. This method will produce a particular solution of a nonhomogenous system $\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t)$ provided that we know a fundamental matrix for the complementary system. To derive the method, suppose $Y$ is a fundamental matrix for the complementary system; that is,

$$
Y=\left[\begin{array}{cccc}
y_{11} & y_{12} & \cdots & y_{1 n} \\
y_{21} & y_{22} & \cdots & y_{2 n} \\
\vdots & \vdots & \ddots & \vdots \\
y_{n 1} & y_{n 2} & \cdots & y_{n n}
\end{array}\right]
$$

where

$$
\mathbf{y}_{1}=\left[\begin{array}{c}
y_{11} \\
y_{21} \\
\vdots \\
y_{n 1}
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
y_{12} \\
y_{22} \\
\vdots \\
y_{n 2}
\end{array}\right], \quad \cdots, \quad \mathbf{y}_{n}=\left[\begin{array}{c}
y_{1 n} \\
y_{2 n} \\
\vdots \\
y_{n n}
\end{array}\right]
$$

is a fundamental set of solutions of the complementary system. In Section 10.3 we saw that $Y^{\prime}=A(t) Y$. We seek a particular solution of

$$
\begin{equation*}
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t) \tag{10.7.1}
\end{equation*}
$$

of the form

$$
\begin{equation*}
\mathbf{y}_{p}=Y \mathbf{u} \tag{10.7.2}
\end{equation*}
$$

where $\mathbf{u}$ is to be determined. Differentiating (10.7.2) yields

$$
\begin{aligned}
\mathbf{y}_{p}^{\prime} & =Y^{\prime} \mathbf{u}+Y \mathbf{u}^{\prime} \\
& =A Y \mathbf{u}+Y \mathbf{u}^{\prime}\left(\text { since } Y^{\prime}=A Y\right) \\
& =A \mathbf{y}_{p}+Y \mathbf{u}^{\prime}\left(\text { since } Y \mathbf{u}=\mathbf{y}_{p}\right)
\end{aligned}
$$

Comparing this with (10.7.1) shows that $\mathbf{y}_{p}=Y \mathbf{u}$ is a solution of (10.7.1) if and only if

$$
Y \mathbf{u}^{\prime}=\mathbf{f}
$$

Thus, we can find a particular solution $\mathbf{y}_{p}$ by solving this equation for $\mathbf{u}^{\prime}$, integrating to obtain $\mathbf{u}$, and computing $Y \mathbf{u}$. We can take all constants of integration to be zero, since any particular solution will suffice.

Exercise 22 sketches a proof that this method is analogous to the method of variation of parameters discussed in Sections 5.7 and 9.4 for scalar linear equations.

## Example 10.7.1

(a) Find a particular solution of the system

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
1 & 2  \tag{10.7.3}\\
2 & 1
\end{array}\right] \mathbf{y}+\left[\begin{array}{c}
2 e^{4 t} \\
e^{4 t}
\end{array}\right]
$$

which we considered in Example 10.2.1.
(b) Find the general solution of (10.7.3).
$\underline{\text { SOLUTION(a) The complementary system is }}$

$$
\mathbf{y}^{\prime}=\left[\begin{array}{ll}
1 & 2  \tag{10.7.4}\\
2 & 1
\end{array}\right] \mathbf{y} .
$$

The characteristic polynomial of the coefficient matrix is

$$
\left|\begin{array}{cc}
1-\lambda & 2 \\
2 & 1-\lambda
\end{array}\right|=(\lambda+1)(\lambda-3) .
$$

Using the method of Section 10.4, we find that

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
e^{3 t} \\
e^{3 t}
\end{array}\right] \quad \text { and } \quad \mathbf{y}_{2}=\left[\begin{array}{r}
e^{-t} \\
-e^{-t}
\end{array}\right]
$$

are linearly independent solutions of (10.7.4). Therefore

$$
Y=\left[\begin{array}{rr}
e^{3 t} & e^{-t} \\
e^{3 t} & -e^{-t}
\end{array}\right]
$$

is a fundamental matrix for (10.7.4). We seek a particular solution $\mathbf{y}_{p}=Y \mathbf{u}$ of (10.7.3), where $Y \mathbf{u}^{\prime}=\mathbf{f}$; that is,

$$
\left[\begin{array}{cc}
e^{3 t} & e^{-t} \\
e^{3 t} & -e^{-t}
\end{array}\right]\left[\begin{array}{l}
u_{1}^{\prime} \\
u_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{c}
2 e^{4 t} \\
e^{4 t}
\end{array}\right]
$$

The determinant of $Y$ is the Wronskian

$$
\left|\begin{array}{rr}
e^{3 t} & e^{-t} \\
e^{3 t} & -e^{-t}
\end{array}\right|=-2 e^{2 t}
$$

By Cramer's rule,

$$
\begin{aligned}
& u_{1}^{\prime}=-\frac{1}{2 e^{2 t}}\left|\begin{array}{cc}
2 e^{4 t} & e^{-t} \\
e^{4 t} & -e^{-t}
\end{array}\right|=\frac{3 e^{3 t}}{2 e^{2 t}}=\frac{3}{2} e^{t} \\
& u_{2}^{\prime}=-\frac{1}{2 e^{2 t}}\left|\begin{array}{cc}
e^{3 t} & 2 e^{4 t} \\
e^{3 t} & e^{4 t}
\end{array}\right|=\frac{e^{7 t}}{2 e^{2 t}}=\frac{1}{2} e^{5 t} .
\end{aligned}
$$

Therefore

$$
\mathbf{u}^{\prime}=\frac{1}{2}\left[\begin{array}{l}
3 e^{t} \\
e^{5 t}
\end{array}\right]
$$

Integrating and taking the constants of integration to be zero yields

$$
\mathbf{u}=\frac{1}{10}\left[\begin{array}{c}
15 e^{t} \\
e^{5 t}
\end{array}\right]
$$

so

$$
\mathbf{y}_{p}=Y \mathbf{u}=\frac{1}{10}\left[\begin{array}{rr}
e^{3 t} & e^{-t} \\
e^{3 t} & -e^{-t}
\end{array}\right]\left[\begin{array}{c}
15 e^{t} \\
e^{5 t}
\end{array}\right]=\frac{1}{5}\left[\begin{array}{c}
8 e^{4 t} \\
7 e^{4 t}
\end{array}\right]
$$

is a particular solution of (10.7.3).

SOLUTION(b) From Theorem 10.7.1, the general solution of (10.7.3) is

$$
\mathbf{y}=\mathbf{y}_{p}+c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}=\frac{1}{5}\left[\begin{array}{c}
8 e^{4 t}  \tag{10.7.5}\\
7 e^{4 t}
\end{array}\right]+c_{1}\left[\begin{array}{l}
e^{3 t} \\
e^{3 t}
\end{array}\right]+c_{2}\left[\begin{array}{r}
e^{-t} \\
-e^{-t}
\end{array}\right]
$$

which can also be written as

$$
\mathbf{y}=\frac{1}{5}\left[\begin{array}{l}
8 e^{4 t} \\
7 e^{4 t}
\end{array}\right]+\left[\begin{array}{rr}
e^{3 t} & e^{-t} \\
e^{3 t} & -e^{-t}
\end{array}\right] \mathbf{c},
$$

where $\mathbf{c}$ is an arbitrary constant vector.
Writing (10.7.5) in terms of coordinates yields

$$
\begin{aligned}
& y_{1}=\frac{8}{5} e^{4 t}+c_{1} e^{3 t}+c_{2} e^{-t} \\
& y_{2}=\frac{7}{5} e^{4 t}+c_{1} e^{3 t}-c_{2} e^{-t}
\end{aligned}
$$

so our result is consistent with Example 10.2.1.
If $A$ isn't a constant matrix, it's usually difficult to find a fundamental set of solutions for the system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$. It is beyond the scope of this text to discuss methods for doing this. Therefore, in the following examples and in the exercises involving systems with variable coefficient matrices we'll provide fundamental matrices for the complementary systems without explaining how they were obtained.

Example 10.7.2 Find a particular solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{cc}
2 & 2 e^{-2 t}  \tag{10.7.6}\\
2 e^{2 t} & 4
\end{array}\right] \mathbf{y}+\left[\begin{array}{l}
1 \\
1
\end{array}\right]
$$

given that

$$
Y=\left[\begin{array}{ll}
e^{4 t} & -1 \\
e^{6 t} & e^{2 t}
\end{array}\right]
$$

is a fundamental matrix for the complementary system.

Solution We seek a particular solution $\mathbf{y}_{p}=Y \mathbf{u}$ of (10.7.6) where $Y \mathbf{u}^{\prime}=\mathbf{f}$; that is,

$$
\left[\begin{array}{cc}
e^{4 t} & -1 \\
e^{6 t} & e^{2 t}
\end{array}\right]\left[\begin{array}{c}
u_{1}^{\prime} \\
u_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{l}
1 \\
1
\end{array}\right]
$$

The determinant of $Y$ is the Wronskian

$$
\left|\begin{array}{cc}
e^{4 t} & -1 \\
e^{6 t} & e^{2 t}
\end{array}\right|=2 e^{6 t}
$$

By Cramer's rule,

$$
\begin{aligned}
& u_{1}^{\prime}=\frac{1}{2 e^{6 t}}\left|\begin{array}{cc}
1 & -1 \\
1 & e^{2 t}
\end{array}\right|=\frac{e^{2 t}+1}{2 e^{6 t}}=\frac{e^{-4 t}+e^{-6 t}}{2} \\
& u_{2}^{\prime}=\frac{1}{2 e^{6 t}}\left|\begin{array}{cc}
e^{4 t} & 1 \\
e^{6 t} & 1
\end{array}\right|=\frac{e^{4 t}-e^{6 t}}{2 e^{6 t}}=\frac{e^{-2 t}-1}{2} .
\end{aligned}
$$

Therefore

$$
\mathbf{u}^{\prime}=\frac{1}{2}\left[\begin{array}{c}
e^{-4 t}+e^{-6 t} \\
e^{-2 t}-1
\end{array}\right]
$$

Integrating and taking the constants of integration to be zero yields

$$
\mathbf{u}=-\frac{1}{24}\left[\begin{array}{c}
3 e^{-4 t}+2 e^{-6 t} \\
6 e^{-2 t}+12 t
\end{array}\right]
$$

so

$$
\mathbf{y}_{p}=Y \mathbf{u}=-\frac{1}{24}\left[\begin{array}{ll}
e^{4 t} & -1 \\
e^{6 t} & e^{2 t}
\end{array}\right]\left[\begin{array}{c}
3 e^{-4 t}+2 e^{-6 t} \\
6 e^{-2 t}+12 t
\end{array}\right]=\frac{1}{24}\left[\begin{array}{c}
4 e^{-2 t}+12 t-3 \\
-3 e^{2 t}(4 t+1)-8
\end{array}\right]
$$

is a particular solution of (10.7.6).

Example 10.7.3 Find a particular solution of

$$
\mathbf{y}^{\prime}=-\frac{2}{t^{2}}\left[\begin{array}{cc}
t & -3 t^{2}  \tag{10.7.7}\\
1 & -2 t
\end{array}\right] \mathbf{y}+t^{2}\left[\begin{array}{l}
1 \\
1
\end{array}\right]
$$

given that

$$
Y=\left[\begin{array}{cc}
2 t & 3 t^{2} \\
1 & 2 t
\end{array}\right]
$$

is a fundamental matrix for the complementary system on $(-\infty, 0)$ and $(0, \infty)$.

Solution We seek a particular solution $\mathbf{y}_{p}=Y \mathbf{u}$ of (10.7.7) where $Y \mathbf{u}^{\prime}=\mathbf{f}$; that is,

$$
\left[\begin{array}{cc}
2 t & 3 t^{2} \\
1 & 2 t
\end{array}\right]\left[\begin{array}{l}
u_{1}^{\prime} \\
u_{2}^{\prime}
\end{array}\right]=\left[\begin{array}{c}
t^{2} \\
t^{2}
\end{array}\right] .
$$

The determinant of $Y$ is the Wronskian

$$
\left|\begin{array}{cc}
2 t & 3 t^{2} \\
1 & 2 t
\end{array}\right|=t^{2}
$$

By Cramer's rule,

$$
\begin{aligned}
& u_{1}^{\prime}=\frac{1}{t^{2}}\left|\begin{array}{cc}
t^{2} & 3 t^{2} \\
t^{2} & 2 t
\end{array}\right|=\frac{2 t^{3}-3 t^{4}}{t^{2}}=2 t-3 t^{2} \\
& u_{2}^{\prime}=\frac{1}{t^{2}}\left|\begin{array}{cc}
2 t & t^{2} \\
1 & t^{2}
\end{array}\right|=\frac{2 t^{3}-t^{2}}{t^{2}}=2 t-1
\end{aligned}
$$

Therefore

$$
\mathbf{u}^{\prime}=\left[\begin{array}{c}
2 t-3 t^{2} \\
2 t-1
\end{array}\right]
$$

Integrating and taking the constants of integration to be zero yields

$$
\mathbf{u}=\left[\begin{array}{c}
t^{2}-t^{3} \\
t^{2}-t
\end{array}\right]
$$

so

$$
\mathbf{y}_{p}=Y \mathbf{u}=\left[\begin{array}{cc}
2 t & 3 t^{2} \\
1 & 2 t
\end{array}\right]\left[\begin{array}{c}
t^{2}-t^{3} \\
t^{2}-t
\end{array}\right]=\left[\begin{array}{c}
t^{3}(t-1) \\
t^{2}(t-1)
\end{array}\right]
$$

is a particular solution of (10.7.7).
Example 10.7.4
(a) Find a particular solution of

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
2 & -1 & -1  \tag{10.7.8}\\
1 & 0 & -1 \\
1 & -1 & 0
\end{array}\right] \mathbf{y}+\left[\begin{array}{c}
e^{t} \\
0 \\
e^{-t}
\end{array}\right]
$$

(b) Find the general solution of (10.7.8).
$\underline{\operatorname{SOLUTION}(\mathbf{a})}$ The complementary system for (10.7.8) is

$$
\mathbf{y}^{\prime}=\left[\begin{array}{rrr}
2 & -1 & -1  \tag{10.7.9}\\
1 & 0 & -1 \\
1 & -1 & 0
\end{array}\right] \mathbf{y} .
$$

The characteristic polynomial of the coefficient matrix is

$$
\left|\begin{array}{ccc}
2-\lambda & -1 & -1 \\
1 & -\lambda & -1 \\
1 & -1 & -\lambda
\end{array}\right|=-\lambda(\lambda-1)^{2} .
$$

Using the method of Section 10.4, we find that

$$
\mathbf{y}_{1}=\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right], \quad \mathbf{y}_{2}=\left[\begin{array}{c}
e^{t} \\
e^{t} \\
0
\end{array}\right], \quad \text { and } \quad \mathbf{y}_{3}=\left[\begin{array}{c}
e^{t} \\
0 \\
e^{t}
\end{array}\right]
$$

are linearly independent solutions of (10.7.9). Therefore

$$
Y=\left[\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & e^{t} & 0 \\
1 & 0 & e^{t}
\end{array}\right]
$$

is a fundamental matrix for (10.7.9). We seek a particular solution $\mathbf{y}_{p}=Y \mathbf{u}$ of (10.7.8), where $Y \mathbf{u}^{\prime}=\mathbf{f}$; that is,

$$
\left[\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & e^{t} & 0 \\
1 & 0 & e^{t}
\end{array}\right]\left[\begin{array}{l}
u_{1}^{\prime} \\
u_{2}^{\prime} \\
u_{3}^{\prime}
\end{array}\right]=\left[\begin{array}{c}
e^{t} \\
0 \\
e^{-t}
\end{array}\right]
$$

The determinant of $Y$ is the Wronskian

$$
\left|\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & e^{t} & 0 \\
1 & 0 & e^{t}
\end{array}\right|=-e^{2 t}
$$

Thus, by Cramer's rule,

$$
\begin{aligned}
& u_{1}^{\prime}=-\frac{1}{e^{2 t}}\left|\begin{array}{ccc}
e^{t} & e^{t} & e^{t} \\
0 & e^{t} & 0 \\
e^{-t} & 0 & e^{t}
\end{array}\right|=-\frac{e^{3 t}-e^{t}}{e^{2 t}}=e^{-t}-e^{t} \\
& u_{2}^{\prime}=-\frac{1}{e^{2 t}}\left|\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & 0 & 0 \\
1 & e^{-t} & e^{t}
\end{array}\right|=-\frac{1-e^{2 t}}{e^{2 t}}=1-e^{-2 t} \\
& u_{3}^{\prime}=-\frac{1}{e^{2 t}}\left|\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & e^{t} & 0 \\
1 & 0 & e^{-t}
\end{array}\right|=\frac{e^{2 t}}{e^{2 t}}=1 .
\end{aligned}
$$

Therefore

$$
\mathbf{u}^{\prime}=\left[\begin{array}{c}
e^{-t}-e^{t} \\
1-e^{-2 t} \\
1
\end{array}\right]
$$

Integrating and taking the constants of integration to be zero yields

$$
\mathbf{u}=\left[\begin{array}{c}
-e^{t}-e^{-t} \\
\frac{e^{-2 t}}{2}+t \\
t
\end{array}\right]
$$

so

$$
\mathbf{y}_{p}=Y \mathbf{u}=\left[\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & e^{t} & 0 \\
1 & 0 & e^{t}
\end{array}\right]\left[\begin{array}{c}
-e^{t}-e^{-t} \\
\frac{e^{-2 t}}{2}+t \\
t
\end{array}\right]=\left[\begin{array}{c}
e^{t}(2 t-1)-\frac{e^{-t}}{2} \\
e^{t}(t-1)-\frac{e^{-t}}{2} \\
e^{t}(t-1)-e^{-t}
\end{array}\right]
$$

is a particular solution of (10.7.8).
$\underline{\text { SOLUTION(a) From Theorem 10.7.1 the general solution of (10.7.8) is }}$

$$
\mathbf{y}=\mathbf{y}_{p}+c_{1} \mathbf{y}_{1}+c_{2} \mathbf{y}_{2}+c_{3} \mathbf{y}_{3}=\left[\begin{array}{c}
e^{t}(2 t-1)-\frac{e^{-t}}{2} \\
e^{t}(t-1)-\frac{e^{-t}}{2} \\
e^{t}(t-1)-e^{-t}
\end{array}\right]+c_{1}\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right]+c_{2}\left[\begin{array}{c}
e^{t} \\
e^{t} \\
0
\end{array}\right]+c_{3}\left[\begin{array}{c}
e^{t} \\
0 \\
e^{t}
\end{array}\right],
$$

which can be written as

$$
\mathbf{y}=\mathbf{y}_{p}+Y \mathbf{c}=\left[\begin{array}{c}
e^{t}(2 t-1)-\frac{e^{-t}}{2} \\
e^{t}(t-1)-\frac{e^{-t}}{2} \\
e^{t}(t-1)-e^{-t}
\end{array}\right]+\left[\begin{array}{ccc}
1 & e^{t} & e^{t} \\
1 & e^{t} & 0 \\
1 & 0 & e^{t}
\end{array}\right] \mathbf{c}
$$

where $\mathbf{c}$ is an arbitrary constant vector.
Example 10.7.5 Find a particular solution of

$$
\mathbf{y}^{\prime}=\frac{1}{2}\left[\begin{array}{ccc}
3 & e^{-t} & -e^{2 t}  \tag{10.7.10}\\
0 & 6 & 0 \\
-e^{-2 t} & e^{-3 t} & -1
\end{array}\right] \mathbf{y}+\left[\begin{array}{c}
1 \\
e^{t} \\
e^{-t}
\end{array}\right]
$$

given that

$$
Y=\left[\begin{array}{ccc}
e^{t} & 0 & e^{2 t} \\
0 & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right]
$$

is a fundamental matrix for the complementary system.

Solution We seek a particular solution of (10.7.10) in the form $\mathbf{y}_{p}=Y \mathbf{u}$, where $Y \mathbf{u}^{\prime}=\mathbf{f}$; that is,

$$
\left[\begin{array}{ccc}
e^{t} & 0 & e^{2 t} \\
0 & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right]\left[\begin{array}{c}
u_{1}^{\prime} \\
u_{2}^{\prime} \\
u_{3}^{\prime}
\end{array}\right]=\left[\begin{array}{c}
1 \\
e^{t} \\
e^{-t}
\end{array}\right] .
$$

The determinant of $Y$ is the Wronskian

$$
\left|\begin{array}{ccc}
e^{t} & 0 & e^{2 t} \\
0 & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right|=-2 e^{4 t}
$$

By Cramer's rule,

$$
\begin{aligned}
u_{1}^{\prime}=-\frac{1}{2 e^{4 t}}\left|\begin{array}{ccc}
1 & 0 & e^{2 t} \\
e^{t} & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right|=\frac{e^{4 t}}{2 e^{4 t}}=\frac{1}{2} \\
u_{2}^{\prime}=-\frac{1}{2 e^{4 t}}\left|\begin{array}{ccc}
e^{t} & 1 & e^{2 t} \\
0 & e^{t} & e^{3 t} \\
e^{-t} & e^{-t} & 0
\end{array}\right|=\frac{e^{3 t}}{2 e^{4 t}}=\frac{1}{2} e^{-t} \\
u_{3}^{\prime}=-\frac{1}{2 e^{4 t}}\left|\begin{array}{ccc}
e^{t} & 0 & 1 \\
0 & e^{3 t} & e^{t} \\
e^{-t} & 1 & e^{-t}
\end{array}\right|=-\frac{e^{3 t}-2 e^{2 t}}{2 e^{4 t}}=\frac{2 e^{-2 t}-e^{-t}}{2} .
\end{aligned}
$$

Therefore

$$
\mathbf{u}^{\prime}=\frac{1}{2}\left[\begin{array}{c}
1 \\
e^{-t} \\
2 e^{-2 t}-e^{-t}
\end{array}\right]
$$

Integrating and taking the constants of integration to be zero yields

$$
\mathbf{u}=\frac{1}{2}\left[\begin{array}{c}
t \\
-e^{-t} \\
e^{-t}-e^{-2 t}
\end{array}\right]
$$

so

$$
\mathbf{y}_{p}=Y \mathbf{u}=\frac{1}{2}\left[\begin{array}{ccc}
e^{t} & 0 & e^{2 t} \\
0 & e^{3 t} & e^{3 t} \\
e^{-t} & 1 & 0
\end{array}\right]\left[\begin{array}{c}
t \\
-e^{-t} \\
e^{-t}-e^{-2 t}
\end{array}\right]=\frac{1}{2}\left[\begin{array}{c}
e^{t}(t+1)-1 \\
-e^{t} \\
e^{-t}(t-1)
\end{array}\right]
$$

is a particular solution of (10.7.10).

### 10.7 Exercises

In Exercises 1-10 find a particular solution.

1. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}-1 & -4 \\ -1 & -1\end{array}\right] \mathbf{y}+\left[\begin{array}{l}21 e^{4 t} \\ 8 e^{-3 t}\end{array}\right]$
2. $\mathbf{y}^{\prime}=\frac{1}{5}\left[\begin{array}{rr}-4 & 3 \\ -2 & -11\end{array}\right] \mathbf{y}+\left[\begin{array}{c}50 e^{3 t} \\ 10 e^{-3 t}\end{array}\right]$
3. $\mathbf{y}^{\prime}=\left[\begin{array}{ll}1 & 2 \\ 2 & 1\end{array}\right] \mathbf{y}+\left[\begin{array}{l}1 \\ t\end{array}\right]$
4. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-4 & -3 \\ 6 & 5\end{array}\right] \mathbf{y}+\left[\begin{array}{c}2 \\ -2 e^{t}\end{array}\right]$
5. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rr}-6 & -3 \\ 1 & -2\end{array}\right] \mathbf{y}+\left[\begin{array}{l}4 e^{-3 t} \\ 4 e^{-5 t}\end{array}\right]$
6. $\mathbf{y}^{\prime}=\left[\begin{array}{rr}0 & 1 \\ -1 & 0\end{array}\right] \mathbf{y}+\left[\begin{array}{l}1 \\ t\end{array}\right]$
7. $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & 1 & -1 \\ 3 & 5 & 1 \\ -6 & 2 & 4\end{array}\right] \mathbf{y}+\left[\begin{array}{l}3 \\ 6 \\ 3\end{array}\right]$
8. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & -1 & -1 \\ -2 & 3 & 2 \\ 4 & -1 & -2\end{array}\right] \mathbf{y}+\left[\begin{array}{c}1 \\ e^{t} \\ e^{t}\end{array}\right]$
9. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{rrr}-3 & 2 & 2 \\ 2 & -3 & 2 \\ 2 & 2 & -3\end{array}\right] \mathbf{y}+\left[\begin{array}{c}e^{t} \\ e^{-5 t} \\ e^{t}\end{array}\right]$
10. $\mathbf{y}^{\prime}=\frac{1}{3}\left[\begin{array}{rrr}1 & 1 & -3 \\ -4 & -4 & 3 \\ -2 & 1 & 0\end{array}\right] \mathbf{y}+\left[\begin{array}{c}e^{t} \\ e^{t} \\ e^{t}\end{array}\right]$

In Exercises 11-20 find a particular solution, given that $Y$ is a fundamental matrix for the complementary system.
11. $\quad \mathbf{y}^{\prime}=\frac{1}{t}\left[\begin{array}{rr}1 & t \\ -t & 1\end{array}\right] \mathbf{y}+t\left[\begin{array}{c}\cos t \\ \sin t\end{array}\right] ; \quad Y=t\left[\begin{array}{rr}\cos t & \sin t \\ -\sin t & \cos t\end{array}\right]$
12. $\mathbf{y}^{\prime}=\frac{1}{t}\left[\begin{array}{ll}1 & t \\ t & 1\end{array}\right] \mathbf{y}+\left[\begin{array}{c}t \\ t^{2}\end{array}\right] ; \quad Y=t\left[\begin{array}{cc}e^{t} & e^{-t} \\ e^{t} & -e^{-t}\end{array}\right]$
13. $\quad \mathbf{y}^{\prime}=\frac{1}{t^{2}-1}\left[\begin{array}{rr}t & -1 \\ -1 & t\end{array}\right] \mathbf{y}+t\left[\begin{array}{r}1 \\ -1\end{array}\right] ; \quad Y=\left[\begin{array}{ll}t & 1 \\ 1 & t\end{array}\right]$
14. $\quad \mathbf{y}^{\prime}=\frac{1}{3}\left[\begin{array}{cc}1 & -2 e^{-t} \\ 2 e^{t} & -1\end{array}\right] \mathbf{y}+\left[\begin{array}{c}e^{2 t} \\ e^{-2 t}\end{array}\right] ; \quad Y=\left[\begin{array}{cc}2 & e^{-t} \\ e^{t} & 2\end{array}\right]$
15. $\quad \mathbf{y}^{\prime}=\frac{1}{2 t^{4}}\left[\begin{array}{cc}3 t^{3} & t^{6} \\ 1 & -3 t^{3}\end{array}\right] \mathbf{y}+\frac{1}{t}\left[\begin{array}{c}t^{2} \\ 1\end{array}\right] ; \quad Y=\frac{1}{t^{2}}\left[\begin{array}{rr}t^{3} & t^{4} \\ -1 & t\end{array}\right]$
16. $\mathbf{y}^{\prime}=\left[\begin{array}{cc}\frac{1}{t-1} & -\frac{e^{-t}}{t-1} \\ \frac{e^{t}}{t+1} & \frac{1}{t+1}\end{array}\right] \mathbf{y}+\left[\begin{array}{c}t^{2}-1 \\ t^{2}-1\end{array}\right] ; \quad Y=\left[\begin{array}{cc}t & e^{-t} \\ e^{t} & t\end{array}\right]$
17. $\mathbf{y}^{\prime}=\frac{1}{t}\left[\begin{array}{rrr}1 & 1 & 0 \\ 0 & 2 & 1 \\ -2 & 2 & 2\end{array}\right] \mathbf{y}+\left[\begin{array}{l}1 \\ 2 \\ 1\end{array}\right] \quad Y=\left[\begin{array}{rrr}t^{2} & t^{3} & 1 \\ t^{2} & 2 t^{3} & -1 \\ 0 & 2 t^{3} & 2\end{array}\right]$
18. $\quad \mathbf{y}^{\prime}=\left[\begin{array}{ccc}3 & e^{t} & e^{2 t} \\ e^{-t} & 2 & e^{t} \\ e^{-2 t} & e^{-t} & 1\end{array}\right] \mathbf{y}+\left[\begin{array}{c}e^{3 t} \\ 0 \\ 0\end{array}\right] ; \quad Y=\left[\begin{array}{rrr}e^{5 t} & e^{2 t} & 0 \\ e^{4 t} & 0 & e^{t} \\ e^{3 t} & -1 & -1\end{array}\right]$
19. $\mathbf{y}^{\prime}=\frac{1}{t}\left[\begin{array}{rrr}1 & t & 0 \\ 0 & 1 & t \\ 0 & -t & 1\end{array}\right] \mathbf{y}+\left[\begin{array}{l}t \\ t \\ t\end{array}\right] ; \quad Y=t\left[\begin{array}{rrr}1 & \cos t & \sin t \\ 0 & -\sin t & \cos t \\ 0 & -\cos t & -\sin t\end{array}\right]$
20. $\mathbf{y}^{\prime}=-\frac{1}{t}\left[\begin{array}{rrr}e^{-t} & -t & 1-e^{-t} \\ e^{-t} & 1 & -t-e^{-t} \\ e^{-t} & -t & 1-e^{-t}\end{array}\right] \mathbf{y}+\frac{1}{t}\left[\begin{array}{c}e^{t} \\ 0 \\ e^{t}\end{array}\right] ; \quad Y=\frac{1}{t}\left[\begin{array}{ccc}e^{t} & e^{-t} & t \\ e^{t} & -e^{-t} & e^{-t} \\ e^{t} & e^{-t} & 0\end{array}\right]$
21. Prove Theorem 10.7.1.
22. (a) Convert the scalar equation

$$
\begin{equation*}
P_{0}(t) y^{(n)}+P_{1}(t) y^{(n-1)}+\cdots+P_{n}(t) y=F(t) \tag{A}
\end{equation*}
$$

into an equivalent $n \times n$ system

$$
\begin{equation*}
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t) \tag{B}
\end{equation*}
$$

(b) Suppose (A) is normal on an interval $(a, b)$ and $\left\{y_{1}, y_{2}, \ldots, y_{n}\right\}$ is a fundamental set of solutions of

$$
\begin{equation*}
P_{0}(t) y^{(n)}+P_{1}(t) y^{(n-1)}+\cdots+P_{n}(t) y=0 \tag{C}
\end{equation*}
$$

on $(a, b)$. Find a corresponding fundamental matrix $Y$ for

$$
\begin{equation*}
\mathbf{y}^{\prime}=A(t) \mathbf{y} \tag{D}
\end{equation*}
$$

on $(a, b)$ such that

$$
y=c_{1} y_{1}+c_{2} y_{2}+\cdots+c_{n} y_{n}
$$

is a solution of (C) if and only if $\mathbf{y}=Y \mathbf{c}$ with

$$
\mathbf{c}=\left[\begin{array}{c}
c_{1} \\
c_{2} \\
\vdots \\
c_{n}
\end{array}\right]
$$

is a solution of (D).
(c) Let $y_{p}=u_{1} y_{1}+u_{1} y_{2}+\cdots+u_{n} y_{n}$ be a particular solution of (A), obtained by the method of variation of parameters for scalar equations as given in Section 9.4, and define

$$
\mathbf{u}=\left[\begin{array}{c}
u_{1} \\
u_{2} \\
\vdots \\
u_{n}
\end{array}\right]
$$

Show that $\mathbf{y}_{p}=Y \mathbf{u}$ is a solution of ( B ).
(d) Let $\mathbf{y}_{p}=Y \mathbf{u}$ be a particular solution of (B), obtained by the method of variation of parameters for systems as given in this section. Show that $y_{p}=u_{1} y_{1}+u_{1} y_{2}+\cdots+u_{n} y_{n}$ is a solution of (A).
23. Suppose the $n \times n$ matrix function $A$ and the $n$-vector function $\mathbf{f}$ are continuous on $(a, b)$. Let $t_{0}$ be in $(a, b)$, let $\mathbf{k}$ be an arbitrary constant vector, and let $Y$ be a fundamental matrix for the homogeneous system $\mathbf{y}^{\prime}=A(t) \mathbf{y}$. Use variation of parameters to show that the solution of the initial value problem

$$
\mathbf{y}^{\prime}=A(t) \mathbf{y}+\mathbf{f}(t), \quad \mathbf{y}\left(t_{0}\right)=\mathbf{k}
$$

is

$$
\mathbf{y}(t)=Y(t)\left(Y^{-1}\left(t_{0}\right) \mathbf{k}+\int_{t_{0}}^{t} Y^{-1}(s) \mathbf{f}(s) d s\right)
$$

## A BRIEF TABLE OF INTEGRALS

$$
\begin{aligned}
& \int u^{\alpha} d u=\frac{u^{\alpha+1}}{\alpha+1}+c, \quad \alpha \neq-1 \\
& \int \frac{d u}{u}=\ln |u|+c \\
& \int \cos u d u=\sin u+c \\
& \int \sin u d u=-\cos u+c \\
& \int \tan u d u=-\ln |\cos u|+c \\
& \int \cot u d u=\ln |\sin u|+c \\
& \int \sec ^{2} u d u=\tan u+c \\
& \int \csc ^{2} u d u=-\cot u+c \\
& \int \sec u d u=\ln |\sec u+\tan u|+c \\
& \int \cos ^{2} u d u=\frac{u}{2}+\frac{1}{4} \sin 2 u+c \\
& \int \sin ^{2} u d u=\frac{u}{2}-\frac{1}{4} \sin 2 u+c \\
& \int \frac{d u}{1+u^{2}} d u=\tan ^{-1} u+c \\
& \int \frac{d u}{\sqrt{1-u^{2}}} d u=\sin ^{-1} u+c \\
& \int \frac{1}{u^{2}-1} d u=\frac{1}{2} \ln \left|\frac{u-1}{u+1}\right|+c \\
& \int \cosh u d u=\sinh u+c \\
& \int \sinh u d u=\cosh u+c \\
& \int u d v=u v-\int v d u \\
& \int u \cos u d u=u \sin u+\cos u+c \\
& \int u \sin u d u=-u \cos u+\sin u+c \\
& \int u e^{u} d u=u e^{u}-e^{u}+c \\
& \int e^{\lambda u} \cos \omega u d u=\frac{e^{\lambda u}(\lambda \cos \omega u+\omega \sin \omega u)}{\lambda^{2}+\omega^{2}}+c \\
& \int e^{\lambda u} \sin \omega u d u=\frac{e^{\lambda u}(\lambda \sin \omega u-\omega \cos \omega u)}{\lambda^{2}+\omega^{2}}+c
\end{aligned}
$$

$$
\begin{aligned}
& \int \ln |u| d u=u \ln |u|-u+c \\
& \int u \ln |u| d u=\frac{u^{2} \ln |u|}{2}-\frac{u^{2}}{4}+c \\
& \int \cos \omega_{1} u \cos \omega_{2} u d u=\frac{\sin \left(\omega_{1}+\omega_{2}\right) u}{2\left(\omega_{1}+\omega_{2}\right)}+\frac{\sin \left(\omega_{1}-\omega_{2}\right) u}{2\left(\omega_{1}-\omega_{2}\right)}+c \quad\left(\omega_{1} \neq \pm \omega_{2}\right) \\
& \int \sin \omega_{1} u \sin \omega_{2} u d u=-\frac{\sin \left(\omega_{1}+\omega_{2}\right) u}{2\left(\omega_{1}+\omega_{2}\right)}+\frac{\sin \left(\omega_{1}-\omega_{2}\right) u}{2\left(\omega_{1}-\omega_{2}\right)}+c \quad\left(\omega_{1} \neq \pm \omega_{2}\right) \\
& \int \sin \omega_{1} u \cos \omega_{2} u d u=-\frac{\cos \left(\omega_{1}+\omega_{2}\right) u}{2\left(\omega_{1}+\omega_{2}\right)}-\frac{\cos \left(\omega_{1}-\omega_{2}\right) u}{2\left(\omega_{1}-\omega_{2}\right)}+c \quad\left(\omega_{1} \neq \pm \omega_{2}\right)
\end{aligned}
$$

# Answers to Selected 

## Exercises

## Section 1.2 Answers, pp. 12-13

1.2.1 (p. 12) (a) 3 (b) 2 (c) 1 (d) 2
1.2.3 (p. 12) (a) $y=-\frac{x^{2}}{2}+c$ (b) $y=x \cos x-\sin x+c$
(c) $y=\frac{x^{2}}{2} \ln x-\frac{x^{2}}{4}+c$ (d) $y=-x \cos x+2 \sin x+c_{1}+c_{2} x$
(e) $y=(2 x-4) e^{x}+c_{1}+c_{2} x \quad$ (f) $y=\frac{x^{3}}{3}-\sin x+e^{x}+c_{1}+c_{2} x$
(g) $y=\sin x+c_{1}+c_{2} x+c_{3} x^{2}$ (h) $y=-\frac{x^{5}}{60}+e^{x}+c_{1}+c_{2} x+c_{3} x^{2}$
(i) $y=\frac{7}{64} e^{4 x}+c_{1}+c_{2} x+c_{3} x^{2}$
1.2.4 (p. 12) (a) $y=-(x-1) e^{x}$ (b) $y=1-\frac{1}{2} \cos x^{2}$ (c) $y=3-\ln (\sqrt{2} \cos x)$
(d) $y=-\frac{47}{15}-\frac{37}{5}(x-2)+\frac{x^{5}}{30}$ (e) $y=\frac{1}{4} x e^{2 x}-\frac{1}{4} e^{2 x}+\frac{29}{4}$ (f) $y=x \sin x+2 \cos x-3 x-1$ (g) $y=\left(x^{2}-6 x+12\right) e^{x}+\frac{x^{2}}{2}-8 x-11$
(h) $y=\frac{x^{3}}{3}+\frac{\cos 2 x}{6}+\frac{7}{4} x^{2}-6 x+\frac{7}{8}$ (i) $y=\frac{x^{4}}{12}+\frac{x^{3}}{6}+\frac{1}{2}(x-2)^{2}-\frac{26}{3}(x-2)-\frac{5}{3}$
1.2 .7 (p.13) (a) 576 ft (b) 10 s 1.2 .8 (p.13) (b) $y=01.2 .10$ (p.13) (a) $(-2 c-2, \infty)(-\infty, \infty)$

## Section 2.1 Answers, pp. 36-39

2.1.1 (p. 36) $y=e^{-a x} 2.1 .2$ (p. 36) $y=c e^{-x^{3}} \quad 2.1 .3$ (p. 36) $y=c e^{-(\ln x)^{2} / 2}$
2.1.4 (p. 36) $y=\frac{c}{x^{3}} 2.1 .5$ (p. 36) $y=c e^{1 / x} 2.1 .6$ (p.36) $y=\frac{e^{-(x-1)}}{x} \quad 2.1 .7$ (p. 36) $y=\frac{e}{x \ln x}$
2.1.8 (p. 37) $y=\frac{\pi}{x \sin x} 2.1 .9$ (p. 37) $y=2\left(1+x^{2}\right)$ 2.1.10 (p. 37) $y=3 x^{-k}$
2.1.11 (p. 37) $y=c(\cos k x)^{1 / k} 2.1 .12$ (p. 37) $y=\frac{1}{3}+c e^{-3 x}$ 2.1.13 (p. 37) $y=\frac{2}{x}+\frac{c}{x} e^{x}$
2.1.14(p.37) $y=e^{-x^{2}}\left(\frac{x^{2}}{2}+c\right)$ 2.1.15 (p.37) $y=-\frac{e^{-x}+c}{1+x^{2}} 2.1 .16$ (p.37) $y=\frac{7 \ln |x|}{x}+\frac{3}{2} x+\frac{c}{x}$
2.1.17 (p. 37) $y=(x-1)^{-4}(\ln |x-1|-\cos x+c)$ 2.1.18 (p.37) $y=e^{-x^{2}}\left(\frac{x^{3}}{4}+\frac{c}{x}\right)$
2.1.19 (p.37) $y=\frac{2 \ln |x|}{x^{2}}+\frac{1}{2}+\frac{c}{x^{2}}$ 2.1.20 (p.37) $y=(x+c) \cos x$ 2.1.21 (p.37) $y=\frac{c-\cos x}{(1+x)^{2}}$
2.1.22 (p. 37) $y=-\frac{1}{2} \frac{(x-2)^{3}}{(x-1)}+c \frac{(x-2)^{5}}{(x-1)} \mathbf{2 . 1 . 2 3}(\mathbf{p} .37) y=(x+c) e^{-\sin ^{2} x}$
2.1.24 (p. 37) $y=\frac{e^{x}}{x^{2}}-\frac{e^{x}}{x^{3}}+\frac{c}{x^{2}} . y=\frac{e^{3 x}-e^{-7 x}}{10}$ 2.1.26 (p. 37) $\frac{2 x+1}{\left(1+x^{2}\right)^{2}}$
2.1.27(p.37) $y=\frac{1}{x^{2}} \ln \left(\frac{1+x^{2}}{2}\right) \mathbf{2 . 1 . 2 9}(\mathbf{p . 3 7}) y=\frac{2 \ln |x|}{x}+\frac{x}{2}-\frac{1}{2 x} \mathbf{2 . 1 . 2 8 ( \mathbf { p . 3 7 } )} y=\frac{1}{2}(\sin x+\csc x)$
2.1.29 (p. 37) $y=\frac{2 \ln |x|}{x}+\frac{x}{2}-\frac{1}{2 x}$ 2.1.30 (p.37) $y=(x-1)^{-3}[\ln (1-x)-\cos x]$
2.1.31 (p.37) $y=2 x^{2}+\frac{1}{x^{2}} \quad(0, \infty)$ 2.1.32 (p.37) $y=x^{2}(1-\ln x) 2.1 .33\left(\right.$ p.37) $y=\frac{1}{2}+\frac{5}{2} e^{-x^{2}}$
2.1.34 (p. 37) $y=\frac{\ln |x-1|+\tan x+1}{(x-1)^{3}} \quad 2.1 .35$ (p. 37) $y=\frac{\ln |x|+x^{2}+1}{(x+2)^{4}}$
2.1.36 (p. 37) $y=\left(x^{2}-1\right)\left(\frac{1}{2} \ln \left|x^{2}-1\right|-4\right)$
2.1.37 (p. 37) $y=-\left(x^{2}-5\right)\left(7+\ln \left|x^{2}-5\right|\right) 2.1 .38$ (p. 38) $y=e^{-x^{2}}\left(3+\int_{0}^{x} t^{2} e^{t^{2}} d t\right)$
2.1.39 (p. 38) $y=\frac{1}{x}\left(2+\int_{1}^{x} \frac{\sin t}{t} d t\right)$ 2.1.40 (p.38) $y=e^{-x} \int_{1}^{x} \frac{\tan t}{t} d t$
2.1.41(p.38) $y=\frac{1}{1+x^{2}}\left(1+\int_{0}^{x} \frac{e^{t}}{1+t^{2}} d t\right) \mathbf{2 . 1 . 4 2 ( p . 3 8 )} y=\frac{1}{x}\left(2 e^{-(x-1)}+e^{-x} \int_{1}^{x} e^{t} e^{t^{2}} d t\right)$
2.1.43(p.38) $G=\frac{r}{\lambda}+\left(G_{0}-\frac{r}{\lambda}\right) e^{-\lambda t} \lim _{t \rightarrow \infty} G(t)=\frac{r}{\lambda}$ 2.1.45 (p.38) (a) $y=y_{0} e^{-a\left(x-x_{0}\right)}+e^{-a x} \int_{x_{0}}^{x} e^{a t} f(t) d t$
2.1.48 (p. 39) (a) $y=\tan ^{-1}\left(\frac{1}{3}+c e^{3 x}\right) \quad$ (b) $y= \pm\left[\ln \left(\frac{1}{x}+\frac{c}{x^{2}}\right)\right]^{1 / 2}$
$\begin{array}{ll}\text { (c) } y=\exp \left(x^{2}+\frac{c}{x^{2}}\right) & \text { (d) } y=-1+\frac{x}{c+3 \ln |x|}\end{array}$
Section 2.2 Answers, pp. 46-48
2.2.1 (p. 46) $y=2 \pm \sqrt{2\left(x^{3}+x^{2}+x+c\right)}$
2.2.2 (p. 46) $\ln (|\sin y|)=\cos x+c ; y \equiv k \pi, k=$ integer
$2.2 .3\left(\mathbf{p . 4 6 )} y=\frac{c}{x-c} \quad y \equiv-1 \quad \mathbf{2 . 2 . 4}\left(\mathbf{p . 4 6 )} \quad \frac{(\ln y)^{2}}{2}=-\frac{x^{3}}{3}+c\right.\right.$
2.2.5 (p. 46) $y^{3}+3 \sin y+\ln |y|+\ln \left(1+x^{2}\right)+\tan ^{-1} x=c ; y \equiv 0$
2.2.6 (p. 46) $y= \pm\left(1+\left(\frac{x}{1+c x}\right)^{2}\right)^{1 / 2} ; y \equiv \pm 1$
2.2.7(p.46) $y=\tan \left(\frac{x^{3}}{3}+c\right) \mathbf{2 . 2 . 8}\left(\mathbf{p . 4 6 )} y=\frac{c}{\sqrt{1+x^{2}}} \quad 2.2 .9\left(\mathbf{p . 4 6 )} y=\frac{2-c e^{(x-1)^{2} / 2}}{1-c e^{(x-1)^{2} / 2}} ; \quad y \equiv 1\right.\right.$
2.2.10 (p. 46) $y=1+\left(3 x^{2}+9 x+c\right)^{1 / 3}$
2.2.11 (p. 46) $y=2+\sqrt{\frac{2}{3} x^{3}+3 x^{2}+4 x-\frac{11}{3}} 2.2 .12$ (p. 46) $y=\frac{e^{-\left(x^{2}-4\right) / 2}}{2-e^{-\left(x^{2}-4\right) / 2}}$
2.2.13(p.46) $y^{3}+2 y^{2}+x^{2}+\sin x=3 \mathbf{2 . 2 . 1 4}(\mathbf{p . 4 6})(y+1)(y-1)^{-3}(y-2)^{2}=-256(x+1)^{-6}$
2.2.15(p.46) $y=-1+3 e^{-x^{2}} 2.2 .16\left(\mathbf{p . 4 6 )} y=\frac{1}{\sqrt{2 e^{-2 x^{2}}-1}} \mathbf{2 . 2 . 1 7}(\mathbf{p . 4 6}) y \equiv-1 ; \quad(-\infty, \infty)\right.$
$2.2 .18\left(\mathbf{p . 4 6 )} y=\frac{4-e^{-x^{2}}}{2-e^{-x^{2}}} ; \quad(-\infty, \infty) \mathbf{2 . 2 . 1 9}\left(\mathbf{p . 4 6 )} y=\frac{-1+\sqrt{4 x^{2}-15}}{2} ; \quad\left(\frac{\sqrt{15}}{2}, \infty\right)\right.\right.$
$2.2 .20\left(\right.$ p. 46) $y=\frac{2}{1+e^{-2 x}} \quad(-\infty, \infty) 2.2 .21\left(\right.$ p. 46) $y=-\sqrt{25-x^{2}} ; \quad(-5,5)$
$\mathbf{2 . 2 . 2 2}(\mathbf{p . 4 6}) y \equiv 2, \quad(-\infty, \infty) \mathbf{2 . 2 . 2 3}\left(\mathbf{p . 4 6 )} y=3\left(\frac{x+1}{2 x-4}\right)^{1 / 3} ; \quad(-\infty, 2)\right.$
2.2.24(p. 46) $y=\frac{x+c}{1-c x} 2.2 .25\left(\mathbf{p . 4 6 )} y=-x \cos c+\sqrt{1-x^{2}} \sin c ; \quad y \equiv 1 ; y \equiv-1\right.$
2.2 .26 (p. 47) $y=-x+3 \pi / 22.2 .28$ (p. 47) $P=\frac{P_{0}}{\alpha P_{0}+\left(1-\alpha P_{0}\right) e^{-a t}} ; \lim _{t \rightarrow \infty} P(t)=1 / \alpha$
2.2.29(p. 47) $I=\frac{S I_{0}}{I_{0}+\left(S-I_{0}\right) e^{-r S t}}$
2.2.30 (p. 47) If $q=r S$ then $I=\frac{I_{0}}{1+r I_{0} t}$ and $\lim _{t \rightarrow \infty} I(t)=0$. If $q \neq R s$, then $\quad I=$
$\frac{\alpha I_{0}}{I_{0}+\left(\alpha-I_{0}\right) e^{-r \alpha t}}$. If $q<r s$, then $\lim _{t \rightarrow \infty} I(t)=\alpha=S-\frac{q}{r}$
if $q>r S$, then $\lim _{t \rightarrow \infty} I(t)=0 \quad 2.2 .34$ (p. 48) $f=a p, \quad$ where $a=$ constant
2.2 .35 (p. 48) $y=e^{-x}\left(-1 \pm \sqrt{2 x^{2}+c}\right) \quad 2.2 .36\left(\mathbf{p . 4 8 )} y=x^{2}\left(-1+\sqrt{x^{2}+c}\right)\right.$
$\mathbf{2 . 2 . 3 7}\left(\mathbf{p . 4 8 )} y=e^{x}\left(-1+\left(3 x e^{x}+c\right)^{1 / 3}\right)\right.$
2.2.38(p.48) $y=e^{2 x}\left(1 \pm \sqrt{c-x^{2}}\right) \mathbf{2 . 2 . 3 9}\left(\mathbf{p . 4 8 )}\right.$ (a) $y_{1}=1 / x ; \quad g(x)=h(x)$
(b) $y_{1}=x ; \quad g(x)=h(x) / x^{2} \quad$ (c) $y_{1}=e^{-x} ; g(x)=e^{x} h(x)$
(d) $y_{1}=x^{-r} ; g(x)=x^{r-1} h(x)$ (e) $y_{1}=1 / v(x) ; g(x)=v(x) h(x)$

## Section 2.3 Answers, pp. 53-54

2.3.1 (p. 53) (a), (b) $x_{0} \neq k \pi\left(k=\right.$ integer) 2.3 .2 (p. 53) (a), (b) $\left(x_{0}, y_{0}\right) \neq(0,0)$
2.3.3 (p. 53) (a), (b) $x_{0} y_{0} \neq(2 k+1) \frac{\pi}{2}\left(k=\right.$ integer) 2.3.4 (p. 53) (a), (b) $x_{0} y_{0}>0$ and $x_{0} y_{0} \neq 1$
2.3.5 (p. 53) (a) all $\left(x_{0}, y_{0}\right) \quad$ (b) $\left(x_{0}, y_{0}\right)$ with $y_{0} \neq 0 \quad 2.3 .6$ (p. 53) (a), (b) all $\left(x_{0}, y_{0}\right)$
2.3.7 (p. 53) (a), (b) all $\left(x_{0}, y_{0}\right) 2.3 .8$ (p.53) (a), (b) $\left(x_{0}, y_{0}\right)$ such that $x_{0} \neq 4 y_{0}$
2.3 .9 (p. 53) (a) all $\left(x_{0}, y_{0}\right)$ (b) all $\left(x_{0}, y_{0}\right) \neq(0,0) 2.3 .10$ (p. 53) (a) all $\left(x_{0}, y_{0}\right)$
(b) all $\left(x_{0}, y_{0}\right)$ with $y_{0} \neq \pm 12.3 .11$ (p. 53) (a), (b) all $\left(x_{0}, y_{0}\right)$
2.3.12 (p. 53) (a), (b) all $\left(x_{0}, y_{0}\right)$ such that $x_{0}+y_{0}>0$
2.3.13 (p. 53) (a), (b) all $\left(x_{0}, y_{0}\right)$ with $x_{0} \neq 1, \quad y_{0} \neq(2 k+1) \frac{\pi}{2}(k=$ integer $)$
2.3.16 (p. 54) $y=\left(\frac{3}{5} x+1\right)^{5 / 3},-\infty<x<\infty$, is a solution.

Also,

$$
y=\left\{\begin{array}{cl}
0, & -\infty<x \leq-\frac{5}{3} \\
\left(\frac{3}{5} x+1\right)^{5 / 3}, & -\frac{5}{3}<x<\infty
\end{array}\right.
$$

is a solution, For every $a \geq \frac{5}{3}$, the following function is also a solution:

$$
y=\left\{\begin{array}{cl}
\left(\frac{3}{5}(x+a)\right)^{5 / 3}, & -\infty<x<-a \\
0, & -a \leq x \leq-\frac{5}{3} \\
\left(\frac{3}{5} x+1\right)^{5 / 3}, & -\frac{5}{3}<x<\infty
\end{array}\right.
$$

2.3.17 (p. 54) (a) all $\left(x_{0}, y_{0}\right)$ (b) all $\left(x_{0}, y_{0}\right)$ with $y_{0} \neq 1$
2.3.18 (p. 54) $y_{1} \equiv 1 ; y_{2}=1+|x|^{3} ; y_{3}=1-|x|^{3} ; y_{4}=1+x^{3} ; y_{5}=1-x^{3}$

$$
\begin{aligned}
& y_{6}=\left\{\begin{array}{cc}
1+x^{3}, & x \geq 0, \\
1, & x<0
\end{array} ; \quad y_{7}=\left\{\begin{array}{cc}
1-x^{3}, & x \geq 0, \\
1, & x<0
\end{array},\right.\right. \\
& y_{8}=\left\{\begin{array}{cc}
1, & x \geq 0, \\
1+x^{3}, & x<0
\end{array} ; \quad y_{9}=\left\{\begin{array}{cc}
1, & x \geq 0 \\
1-x^{3}, & x<0
\end{array}\right.\right.
\end{aligned}
$$

2.3.19 (p. 54) $y=1+\left(x^{2}+4\right)^{3 / 2}, \quad-\infty<x<\infty$
2.3.20 (p. 54) (a) The solution is unique on $(0, \infty)$. It is given by

$$
y=\left\{\begin{array}{cc}
1, & 0<x \leq \sqrt{5} \\
1-\left(x^{2}-5\right)^{3 / 2}, & \sqrt{5}<x<\infty
\end{array}\right.
$$

(b)

$$
y=\left\{\begin{array}{cl}
1, & -\infty<x \leq \sqrt{5} \\
1-\left(x^{2}-5\right)^{3 / 2}, & \sqrt{5}<x<\infty
\end{array}\right.
$$

is a solution of $(\mathrm{A})$ on $(-\infty, \infty)$. If $\alpha \geq 0$, then

$$
y=\left\{\begin{array}{cl}
1+\left(x^{2}-\alpha^{2}\right)^{3 / 2}, & -\infty<x<-\alpha \\
1, & -\alpha \leq x \leq \sqrt{5} \\
1-\left(x^{2}-5\right)^{3 / 2}, & \sqrt{5}<x<\infty
\end{array}\right.
$$

and

$$
y=\left\{\begin{array}{cl}
1-\left(x^{2}-\alpha^{2}\right)^{3 / 2}, & -\infty<x<-\alpha \\
1, & -\alpha \leq x \leq \sqrt{5} \\
1-\left(x^{2}-5\right)^{3 / 2}, & \sqrt{5}<x<\infty
\end{array}\right.
$$

are also solutions of (A) on $(-\infty, \infty)$.

## Section 2.5 Answers, pp. 60-63

2.5.1 (p. 60) $2 x^{3} y^{2}=c \mathbf{2 . 5 . 2}$ (p. 60) $3 y \sin x+2 x^{2} e^{x}+3 y=c \mathbf{2 . 5 . 3}$ (p. 60) Not exact
2.5.4(p.60) $x^{2}-2 x y^{2}+4 y^{3}=c \mathbf{2 . 5 . 5}(\mathbf{p} .60) x+y=c$ 2.5.6(p. 60) Not exact
$\mathbf{2 . 5 . 7}(\mathbf{p} .60) 2 y^{2} \cos x+3 x y^{3}-x^{2}=c \mathbf{2 . 5 . 8}(\mathbf{p} .60)$ Not exact
2.5.9 (p. 60) $x^{3}+x^{2} y+4 x y^{2}+9 y^{2}=c$ 2.5.10 (p.60) Not exact 2.5 .11 (p. 60) $\ln |x y|+x^{2}+y^{2}=c$
2.5.12 (p. 60) Not exact 2.5.13 (p. 60) $x^{2}+y^{2}=c \mathbf{2 . 5 . 1 4}$ (p. 60) $x^{2} y^{2} e^{x}+2 y+3 x^{2}=c$
2.5.15 (p. 60) $x^{3} e^{x^{2}+y}-4 y^{3}+2 x^{2}=c \mathbf{2 . 5 . 1 6}$ (p. 60) $x^{4} e^{x y}+3 x y=c$
$\mathbf{2 . 5 . 1 7}$ (p. 60) $x^{3} \cos x y+4 y^{2}+2 x^{2}=c \mathbf{2 . 5 . 1 8}$ (p. 60) $y=\frac{x+\sqrt{2 x^{2}+3 x-1}}{x^{2}}$
2.5.19 (p. 60) $y=\sin x-\sqrt{1-\frac{\tan x}{2}} 2.5 .20$ (p.60) $y=\left(\frac{e^{x}-1}{e^{x}+1}\right)^{1 / 3}$
2.5.21 (p. 60) $y=1+2 \tan x$ 2.5.22 (p. 60) $y=\frac{x^{2}-x+6}{(x+2)(x-3)}$
$\mathbf{2 . 5 . 2 3}\left(\mathbf{p . 6 0 )} \frac{7 x^{2}}{2}+4 x y+\frac{3 y^{2}}{2}=c \mathbf{2 . 5 . 2 4}\left(\mathbf{p . 6 1 )}\left(x^{4} y^{2}+1\right) e^{x}+y^{2}=c\right.\right.$
2.5.29 (p. 61) (a) $M(x, y)=2 x y+f(x)$ (b) $M(x, y)=2(\sin x+x \cos x)(y \sin y+\cos y)+f(x)$
(c) $M(x, y)=y e^{x}-e^{y} \cos x+f(x)$
2.5.30 (p. 61) (a) $N(x, y)=\frac{x^{4} y}{2}+x^{2}+6 x y+g(y)$ (b) $N(x, y)=\frac{x}{y}+2 y \sin x+g(y)$
(c) $N(x, y)=x(\sin y+y \cos y)+g(y)$
2.5.33(p. 62) $B=C$ 2.5.34 (p. 62) $B=2 D, \quad E=2 C$
2.5.37 (p. 62) (a) $2 x^{2}+x^{4} y^{4}+y^{2}=c$ (b) $x^{3}+3 x y^{2}=c$ (c) $x^{3}+y^{2}+2 x y=c$
2.5.38 (p. 62) $y=-1-\frac{1}{x^{2}} \mathbf{2 . 5 . 3 9}$ (p. 62) $y=x^{3}\left(\frac{-3\left(x^{2}+1\right)+\sqrt{9 x^{4}+34 x^{2}+21}}{2}\right)$
2.5.40 (p. 62) $y=-e^{-x^{2}}\left(\frac{2 x+\sqrt{9-5 x^{2}}}{3}\right)$.
2.5.44 (p. 63) (a) $G(x, y)=2 x y+c$ (b) $G(x, y)=e^{x} \sin y+c$
(c) $G(x, y)=3 x^{2} y-y^{3}+c(\mathbf{d}) G(x, y)=-\sin x \sinh y+c$
(e) $G(x, y)=\cos x \sinh y+c$

## Section 2.6 Answers, pp. 70-71

2.6.3 (p. 70) $\mu(x)=1 / x^{2} ; y=c x$ and $\mu(y)=1 / y^{2} ; x=c y$
$2.6 .4\left(\mathbf{p . 7 0 )} \mu(x)=x^{-3 / 2} ; x^{3 / 2} y=c 2.6 .5\right.$ (p. 70) $\mu(y)=1 / y^{3} ; y^{3} e^{2 x}=c$
$2.6 .6\left(\mathbf{p . 7 0 )} \mu(x)=e^{5 x / 2} ; e^{5 x / 2}(x y+1)=c 2.6 .7\left(\mathbf{p . 7 1 )} \mu(x)=e^{x} ; e^{x}(x y+y+x)=c\right.\right.$
2.6.8(p.71) $\mu(x)=x ; x^{2} y^{2}(9 x+4 y)=c 2.6 .9\left(\mathbf{p . 7 1 )} \mu(y)=y^{2} ; y^{3}\left(3 x^{2} y+2 x+1\right)=c \mathbf{2 . 6 . 1 0}\right.$
(p. 71) $\mu(y)=y e^{y} ; e^{y}\left(x y^{3}+1\right)=c \mathbf{2 . 6 . 1 1}\left(\mathbf{p . 7 1 )} \mu(y)=y^{2} ; y^{3}\left(3 x^{4}+8 x^{3} y+y\right)=c\right.$
2.6.12 (p. 71) $\mu(x)=x e^{x} ; x^{2} y(x+1) e^{x}=c$
2.6.13(p. 71) $\mu(x)=\left(x^{3}-1\right)^{-4 / 3} ; x y\left(x^{3}-1\right)^{-1 / 3}=c$ and $x \equiv 1$
2.6.14(p.71) $\mu(y)=e^{y} ; e^{y}(\sin x \cos y+y-1)=c$ 2.6.15 (p.71) $\mu(y)=e^{-y^{2}} ; x y e^{-y^{2}}(x+y)=c$
2.6.16 (p. 71) $\frac{x y}{\sin y}=c$ and $y=k \pi(k=$ integer $) 2.6 .17$ (p. 71) $\mu(x, y)=x^{4} y^{3} ; x^{5} y^{4} \ln x=c$
2.6.18 (p. 71) $\mu(x, y)=1 / x y ;|x|^{\alpha}|y|^{\beta} e^{\gamma x} e^{\delta y}=c$ and $x \equiv 0, y \equiv 0$
2.6.19 (p. 71) $\mu(x, y)=x^{-2} y^{-3} ; 3 x^{2} y^{2}+y=1+c x y^{2}$ and $x \equiv 0, y \equiv 0$
2.6.20 (p. 71) $\mu(x, y)=x^{-2} y^{-1} ;-\frac{2}{x}+y^{3}+3 \ln |y|=c$ and $x \equiv 0, y \equiv 0$
2.6.21 (p. 71) $\mu(x, y)=e^{a x} e^{b y} ; e^{a x} e^{b y} \cos x y=c$
2.6.22(p.71) $\mu(x, y)=x^{-4} y^{-3}$ (and others) $x y=c 2.6 .23\left(\mathbf{p . 7 1 )} \mu(x, y)=x e^{y} ; x^{2} y e^{y} \sin x=c\right.$
2.6.24(p.71) $\mu(x)=1 / x^{2} ; \frac{x^{3} y^{3}}{3}-\frac{y}{x}=c 2.6 .25\left(\right.$ p. 71) $\mu(x)=x+1 ; y(x+1)^{2}(x+y)=c$
2.6.26(p. 71) $\mu(x, y)=x^{2} y^{2} ; x^{3} y^{3}\left(3 x+2 y^{2}\right)=c$
2.6.27 (p. 71) $\mu(x, y)=x^{-2} y^{-2} ; 3 x^{2} y=c x y+2$ and $x \equiv 0, y \equiv 0$

## Section 3.1 Answers, pp. 82-84

3.1.1 (p. 82) $y_{1}=1.450000000, y_{2}=2.085625000, y_{3}=3.079099746$
3.1.2 (p. 82) $y_{1}=1.200000000, y_{2}=1.440415946, y_{3}=1.729880994$
3.1.3 (p. 82) $y_{1}=1.900000000, y_{2}=1.781375000, y_{3}=1.646612970$
3.1.4 (p. 82) $y_{1}=2.962500000, y_{2}=2.922635828, y_{3}=2.880205639$
3.1.5 (p. 83) $y_{1}=2.513274123, y_{2}=1.814517822, y_{3}=1.216364496$
3.1.6 (p. 83)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 1.0 | 48.298147362 | 51.492825643 | 53.076673685 | 54.647937102 |

3.1.7 (p. 83)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | 1.390242009 | 1.370996758 | 1.361921132 | 1.353193719 |

3.1.8 (p. 83)

| $x$ | $h=0.05$ | $h=0.025$ | $h=0.0125$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 1.50 | 7.886170437 | 8.852463793 | 9.548039907 | 10.500000000 |

3.1.9 (p. 83)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | $h=0.1$ | $h=0.05$ | $h=0.025$ |  |  |  |  |  |
| :---: | :---: | :---: | :---: | ---: | ---: | ---: | :---: | :---: | :---: | :---: | :---: |
| 3.0 | 1.469458241 | 1.462514486 | 1.459217010 | 0.3210 | 0.1537 | 0.0753 |  |  |  |  |  |
| Approximate Solutions |  |  |  |  |  |  |  | Residuals |  |  |  |

3.1.10 (p. 83)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | $h=0.1$ | $h=0.05$ | $h=0.025$ |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | ---: | :---: | :---: | :---: | :---: | :---: |
| 2.0 | 0.473456737 | 0.483227470 | 0.487986391 | -0.3129 | -0.1563 | -0.0781 |  |  |  |  |  |
| Approximate Solutions |  |  |  |  |  |  |  | Residuals |  |  |  |

3.1.11 (p. 84)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 1.0 | 0.691066797 | 0.676269516 | 0.668327471 | 0.659957689 |

3.1.12 (p. 84)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | -0.772381768 | -0.761510960 | -0.756179726 | -0.750912371 |

3.1.13 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| 1.0 | 0.538871178 | 0.593002325 | 0.620131525 | 0.647231889 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| 1.0 | 0.647231889 | 0.647231889 | 0.647231889 | 0.647231889 |

Applying variation of parameters to the given initial value problem yields
$y=u e^{-3 x}$, where (A) $u^{\prime}=7, \quad u(0)=6$. Since $u^{\prime \prime}=0$, Euler's method yields the exact
solution of (A). Therefore the Euler semilinear method produces the exact solution of the
given problem
3.1.14 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 12.804226135 | 13.912944662 | 14.559623055 | 15.282004826 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 15.354122287 | 15.317257705 | 15.299429421 | 15.282004826 |

3.1.15 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.867565004 | 0.885719263 | 0.895024772 | 0.904276722 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.569670789 | 0.720861858 | 0.808438261 | 0.904276722 |

3.1.16 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 3.0 | 0.922094379 | 0.945604800 | 0.956752868 | 0.967523153 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 3.0 | 0.993954754 | 0.980751307 | 0.974140320 | 0.967523153 |

3.1.17 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 0.319892131 | 0.330797109 | 0.337020123 | 0.343780513 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 0.305596953 | 0.323340268 | 0.333204519 | 0.343780513 |

3.1.18 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.754572560 | 0.743869878 | 0.738303914 | 0.732638628 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.722610454 | 0.727742966 | 0.730220211 | 0.732638628 |

3.1.19 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 2.175959970 | 2.210259554 | 2.227207500 | 2.244023982 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 2.117953342 | 2.179844585 | 2.211647904 | 2.244023982 |

3.1.20 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 0.032105117 | 0.043997045 | 0.050159310 | 0.056415515 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 0.056020154 | 0.056243980 | 0.056336491 | 0.056415515 |

3.1.21 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 28.987816656 | 38.426957516 | 45.367269688 | 54.729594761 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 54.709134946 | 54.724150485 | 54.728228015 | 54.729594761 |

3.1.22 (p. 84)

| Euler's method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 1.361427907 | 1.361320824 | 1.361332589 | 1.361383810 |


| Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 1.291345518 | 1.326535737 | 1.344004102 | 1.361383810 |

Section 3.2 Answers, pp. 92-84
3.2.1 (p. 92) $y_{1}=1.542812500, y_{2}=2.421622101, y_{3}=4.208020541$
3.2.2 (p. 92) $y_{1}=1.220207973, y_{2}=1.489578775 y_{3}=1.819337186$
3.2.3 (p. 92) $y_{1}=1.890687500, y_{2}=1.763784003, y_{3}=1.622698378$
3.2.4 (p.92) $y_{1}=2.961317914 y_{2}=2.920132727 y_{3}=2.876213748$.
3.2.5 (p. 92) $y_{1}=2.478055238, y_{2}=1.844042564, y_{3}=1.313882333$
3.2 .6 (p. 92)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 1.0 | 56.134480009 | 55.003390448 | 54.734674836 | 54.647937102 |

3.2.7 (p. 92)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | 1.353501839 | 1.353288493 | 1.353219485 | 1.353193719 |

3.2 .8 (p. 92)

| $x$ | $h=0.05$ | $h=0.025$ | $h=0.0125$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 1.50 | 10.141969585 | 10.396770409 | 10.472502111 | 10.500000000 |

3.2 .9 (p. 93)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | $h=0.1$ | $h=0.05$ | $h=0.025$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 3.0 | 1.455674816 | 1.455935127 | 1.456001289 | -0.00818 | -0.00207 | -0.000518 |  |  |
| Approximate Solutions |  |  |  |  |  | Residuals |  |  |

3.2 .10 (p. 93)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | $h=0.1$ | $h=0.05$ | $h=0.025$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 2.0 | 0.492862999 | 0.492709931 | 0.492674855 | 0.00335 | 0.000777 | 0.000187 |  |  |  |  |  |  |  |
| Approximate Solutions |  |  |  |  |  |  |  |  |  | Residuals |  |  |  |

3.2.11 (p. 93)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 1.0 | 0.660268159 | 0.660028505 | 0.659974464 | 0.659957689 |

3.2.12 (p. 93)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | -0.749751364 | -0.750637632 | -0.750845571 | -0.750912371 |

3.2.13 (p. 93) Applying variation of parameters to the given initial value problem
$y=u e^{-3 x}$, where (A) $u^{\prime}=1-2 x, \quad u(0)=2$. Since $u^{\prime \prime \prime}=0$, the improved Euler method yields the exact solution of (A). Therefore the improved Euler semilinear method produces the exact solution of the given problem.

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| 1.0 | 0.105660401 | 0.100924399 | 0.099893685 | 0.099574137 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| 1.0 | 0.099574137 | 0.099574137 | 0.099574137 | 0.099574137 |

3.2.14 (p. 93)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 15.107600968 | 15.234856000 | 15.269755072 | 15.282004826 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 15.285231726 | 15.282812424 | 15.282206780 | 15.282004826 |

3.2.15 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.924335375 | 0.907866081 | 0.905058201 | 0.904276722 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.969670789 | 0.920861858 | 0.908438261 | 0.904276722 |

3.2.16 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 3.0 | 0.967473721 | 0.967510790 | 0.967520062 | 0.967523153 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 3.0 | 0.967473721 | 0.967510790 | 0.967520062 | 0.967523153 |

3.2.17 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 0.349176060 | 0.345171664 | 0.344131282 | 0.343780513 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 0.349350206 | 0.345216894 | 0.344142832 | 0.343780513 |

3.2 .18 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.732679223 | 0.732721613 | 0.732667905 | 0.732638628 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.2$ | $h=0.1$ | $h=0.05$ | "Exact" |
| 2.0 | 0.732166678 | 0.732521078 | 0.732609267 | 0.732638628 |

3.2.19 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 2.247880315 | 2.244975181 | 2.244260143 | 2.244023982 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.0500$ | $h=0.0250$ | $h=0.0125$ | "Exact" |
| 1.50 | 2.248603585 | 2.245169707 | 2.244310465 | 2.244023982 |

3.2.20 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 0.059071894 | 0.056999028 | 0.056553023 | 0.056415515 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 0.056295914 | 0.056385765 | 0.056408124 | 0.056415515 |

3.2.21 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 50.534556346 | 53.483947013 | 54.391544440 | 54.729594761 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 1.0 | 54.709041434 | 54.724083572 | 54.728191366 | 54.729594761 |

3.2 .22 (p. 94)

| Improved Euler method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 1.361395309 | 1.361379259 | 1.361382239 | 1.361383810 |


| Improved Euler semilinear method |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: |
| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| 3.0 | 1.375699933 | 1.364730937 | 1.362193997 | 1.361383810 |

3.2.23 (p. 94)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | 1.349489056 | 1.352345900 | 1.352990822 | 1.353193719 |

3.2 .24 (p. 94)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | 1.350890736 | 1.352667599 | 1.353067951 | 1.353193719 |

3.2.25 (p. 94)

| $x$ | $h=0.05$ | $h=0.025$ | $h=0.0125$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 1.50 | 10.133021311 | 10.391655098 | 10.470731411 | 10.500000000 |

3.2 .26 (p. 94)

| $x$ | $h=0.05$ | $h=0.025$ | $h=0.0125$ | Exact |
| :---: | :---: | :---: | :---: | :---: |
| 1.50 | 10.136329642 | 10.393419681 | 10.470731411 | 10.500000000 |

3.2.27 (p. 94)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 1.0 | 0.660846835 | 0.660189749 | 0.660016904 | 0.659957689 |

3.2.28 (p. 94)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 1.0 | 0.660658411 | 0.660136630 | 0.660002840 | 0.659957689 |

3.2.29 (p. 94)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | -0.750626284 | -0.750844513 | -0.750895864 | -0.751331499 |

3.2 .30 (p. 94)

| $x$ | $h=0.1$ | $h=0.05$ | $h=0.025$ | "Exact" |
| :---: | :---: | :---: | :---: | :---: |
| 2.0 | -0.750335016 | -0.750775571 | -0.750879100 | -0.751331499 |

Section 4.2 Answers, pp. 103-105
$4.2 .1(\mathbf{p} .103) \approx 15.15^{\circ} \mathrm{F} 4.2 .2(\mathbf{p} .103) T=-10+110 e^{-t \ln \frac{11}{9}} 4.2 .3(\mathbf{p} .103) \approx 24.33^{\circ} \mathrm{F}$
4.2.4 (p. 103) (a) $91.30^{\circ} \mathrm{F}$ (b) 8.99 minutes after being placed outside (c) never
4.2.5 (p.103) (a) $12: 11: 32$ (b) $12: 47: 334.2 .6$ (p.103) $(85 / 3)^{\circ} C 4.2 .7$ (p.103) $32^{\circ} \mathrm{F} 4.2 .8(\mathbf{p} .103) Q(t)=40\left(1-e^{-3 t / 40}\right)$
4.2 .9 (p. 103) $Q(t)=30-20 e^{-t / 10} 4.2 .10(\mathbf{p} .103) K(t)=.3-.2 e^{-t / 20} 4.2 .11(\mathbf{p} .103) Q(50)=47.5$
(pounds)
4.2 .12 (p.103) 50 gallons 4.2 .13 (p.103) $\min q_{2}=q_{1} / \bar{c} 4.2 .14$ (p.103) $Q=t+300-\frac{234 \times 10^{5}}{(t+300)^{2}}, \quad 0 \leq t \leq 300$
4.2.15 (p. 104) (a) $Q^{\prime}+\frac{2}{25} Q=6-2 e^{-t / 25}$ (b) $Q=75-50 e^{-t / 25}-25 e^{-2 t / 25}$ (c) 75
4.2.16 (p. 104) (a) $T=T_{m}+\left(T_{0}-T_{m}\right) e^{-k t}+\frac{k\left(S_{0}-T_{m}\right)}{\left(k-k_{m}\right)}\left(e^{-k_{m} t}-e^{-k t}\right)$
(b) $T=T_{m}+k\left(S_{0}-T_{m}\right) t e^{-k t}+\left(T_{0}-T_{m}\right) e^{-k t} \quad$ (c) $\lim _{t \rightarrow \infty} T(t)=\lim _{t \rightarrow \infty} S(t)=T_{m}$
4.2 .17 (p. 104) (a) $T^{\prime}=-k\left(1+\frac{a}{a_{m}}\right) T+k\left(T_{m 0}+\frac{a}{a_{m}} T_{0}\right) \quad$ (b) $T=\frac{a T_{0}+a_{m} T_{m 0}}{a+a_{m}}+\frac{a_{m}\left(T_{0}-T_{m 0}\right)}{a+a_{m}} e^{-k\left(1+a / a_{m}\right) t}$, $T_{m}=\frac{a T_{0}+a_{m} T_{m 0}}{a+a_{m}}+\frac{a\left(T_{m 0}-T_{0}\right)}{a+a_{m}} e^{-k\left(1+a / a_{m}\right) t} ;(\mathbf{c}) \lim _{t \rightarrow \infty} T(t)=\lim _{t \rightarrow \infty} T_{m}(t)=\frac{a T_{0}+a_{m} T_{m 0}}{a+a_{m}}$
4.2 .18 (p. 104) $V=\frac{a}{b} \frac{V_{0}}{V_{0}-\left(V_{0}-a / b\right) e^{-a t}}, \quad \lim _{t \rightarrow \infty} V(t)=a / b$
4.2.19 (p. 104) $c_{1}=c\left(1-e^{-r t / W}\right), c_{2}=c\left(1-e^{-r t / W}-\frac{r}{W} t e^{-r t / W}\right)$.
4.2 .20 (p. 104) (a) $c_{n}=c\left(1-e^{-r t / W} \sum_{j=0}^{n-1} \frac{1}{j!}\left(\frac{r t}{W}\right)^{j}\right)$ (b) $c$ (c) 0
4.2.21 (p. 105) Let $c_{\infty}=\frac{c_{1} W_{1}+c_{2} W_{2}}{W_{1}+W_{2}}, \alpha=\frac{c_{2} W_{2}^{2}-c_{1} W_{1}^{2}}{W_{1}+W_{2}}$, and $\beta=\frac{W_{1}+W_{2}}{W_{1} W_{2}}$. Then:
(a) $c_{1}(t)=c_{\infty}+\frac{\alpha}{W_{1}} e^{-r \beta t}, c_{2}(t)=c_{\infty}-\frac{\alpha}{W_{2}} e^{-r \beta t}$
(b) $\lim _{t \rightarrow \infty} c_{1}(t)=\lim _{t \rightarrow \infty} c_{2}(t)=c_{\infty}$

## Section 4.3 Answers, pp. 113-115

4.3.1 (p. 113) $v=-\frac{384}{5}\left(1-e^{-5 t / 12}\right) ;-\frac{384}{5} \mathrm{ft} / \mathrm{s} 4.3 .2$ (p. 114) $k=12 ; \quad v=-16\left(1-e^{-2 t}\right)$
4.3.3 (p. 114) $v=25\left(1-e^{-t}\right) ; 25 \mathrm{ft} / \mathrm{s} 4.3 .4\left(\mathbf{p}\right.$. 114) $v=20-27 e^{-t / 40} 4.3 .5(\mathbf{p} .114) \approx 17.10 \mathrm{ft}$
4.3 .6 (p. 114) $v=-\frac{40\left(13+3 e^{-4 t / 5}\right)}{13-3 e^{-4 t / 5}} ;-40 \mathrm{ft} / \mathrm{s} 4.3 .7$ (p. 114) $v=-128\left(1-e^{-t / 4}\right)$
4.3.9 (p. 114) $T=\frac{m}{k} \ln \left(1+\frac{v_{0} k}{m g}\right) ; \quad y_{m}=y_{0}+\frac{m}{k}\left[v_{0}-\frac{m g}{k} \ln \left(1+\frac{v_{0} k}{m g}\right)\right]$
4.3.10 (p. 114) $v=-\frac{64\left(1-e^{-t}\right)}{1+e^{-t}} ;-64 \mathrm{ft} / \mathrm{s}$
4.3.11 (p. 114) $v=\alpha \frac{v_{0}\left(1+e^{-\beta t}\right)-\alpha\left(1-e^{-\beta t}\right)}{\alpha\left(1+e^{-\beta t}\right)-v_{0}\left(1-e^{-\beta t}\right)} ; \quad-\alpha$, where $\alpha=\sqrt{\frac{m g}{k}}$ and $\beta=2 \sqrt{\frac{k g}{m}}$.
4.3.12 (p. 114) $T=\sqrt{\frac{m}{k g}} \tan ^{-1}\left(v_{0} \sqrt{\frac{k}{m g}}\right) v=-\sqrt{\frac{m g}{k}} ; \frac{1-e^{-2 \sqrt{\frac{g k}{m}}(t-T)}}{1+e^{-2 \sqrt{\frac{g k}{m}}(t-T)}}$
4.3 .13 (p. 114) $s^{\prime}=m g-\frac{a s}{s+1} ; a_{0}=m g .4 .3 .14$ (p.114) (a) $m s^{\prime}=m g-f(s)$
4.3 .15 (p. 115) (a) $v^{\prime}=-9.8+v^{4} / 81$ (b) $v_{T} \approx-5.308 \mathrm{~m} / \mathrm{s}$
4.3.16 (p. 115) (a) $v^{\prime}=-32+8 \sqrt{|v|} ; v_{T}=-16 \mathrm{ft} / \mathrm{s}$ (b) From Exercise 4.3.14(c), $v_{T}$ is the negative number such that $-32+8 \sqrt{\left|v_{T}\right|}=0$; thus, $v_{T}=-16 \mathrm{ft} / \mathrm{s}$.
$4.3 .17(\mathbf{p} .115) \approx 6.76 \mathrm{miles} / \mathrm{s} 4.3 .18(\mathbf{p} .115) \approx 1.47 \mathrm{miles} / \mathrm{s} 4.3 .20(\mathbf{p} .115) \alpha=\frac{g R^{2}}{\left(y_{m}+R\right)^{2}}$
Section 4.4 Answers, pp. 128-129
4.4.1 (p. 128) $\bar{y}=0$ is a stable equilibrium; trajectories are $v^{2}+\frac{y^{4}}{4}=c$
4.4.2 (p. 128) $\bar{y}=0$ is an unstable equilibrium; trajectories are $v^{2}+\frac{2 y^{3}}{3}=c$
4.4.3 (p. 128) $\bar{y}=0$ is a stable equilibrium; trajectories are $v^{2}+\frac{2|y|^{3}}{3}=c$
4.4.4 (p. 128) $\bar{y}=0$ is a stable equilibrium; trajectories are $v^{2}-e^{-y}(y+1)=c$
4.4.5 (p. 128) equilibria: 0 (stable) and $-2,2$ (unstable); trajectories: $2 v^{2}-y^{4}+8 y^{2}=c$; separatrix: $2 v^{2}-y^{4}+8 y^{2}=16$
4.4.6 (p. 128) equilibria: 0 (unstable) and $-2,2$ (stable); trajectories: $2 v^{2}+y^{4}-8 y^{2}=c$; separatrix: $2 v^{2}+y^{4}-8 y^{2}=0$
4.4.7 (p. 128) equilibria: $0,-2,2$ (stable), $-1,1$ (unstable); trajectories:

$$
6 v^{2}+y^{2}\left(2 y^{4}-15 y^{2}+24\right)=c \text {; separatrix: } 6 v^{2}+y^{2}\left(2 y^{4}-15 y^{2}+24\right)=11
$$

4.4.8 (p. 128) equilibria: 0,2 (stable) and $-2,1$ (unstable);
trajectories: $30 v^{2}+y^{2}\left(12 y^{3}-15 y^{2}-80 y+120\right)=c$;
separatrices: $30 v^{2}+y^{2}\left(12 y^{3}-15 y^{2}-80 y+120\right)=496$ and
$30 v^{2}+y^{2}\left(12 y^{3}-15 y^{2}-80 y+120\right)=37$
4.4.9 (p. 128) No equilibria if $a<0 ; 0$ is unstable if $a=0 ; \sqrt{a}$ is stable and $-\sqrt{a}$ is unstable if $a>0$.

* 4.4 .10 (p. 128) 0 is a stable equilibrium if $a \leq 0 ;-\sqrt{a}$ and $\sqrt{a}$ are stable and 0 is unstable if $a>0$.
4.4.11 (p. 128) 0 is unstable if $a \leq 0 ;-\sqrt{a}$ and $\sqrt{a}$ are unstable and 0 is stable if $a>0$.
4.4.12 (p. 128) 0 is stable if $a \leq 0 ; 0$ is stable and $-\sqrt{a}$ and $\sqrt{a}$ are unstable if $a \leq 0$.
4.4.22 (p. 129) An equilibrium solution $\bar{y}$ of $y^{\prime \prime}+p(y)=0$ is unstable if there's an $\epsilon>0$ such that, for every $\delta>0$, there's a solution of (A) with $\sqrt{(y(0)-\bar{y})^{2}+v^{2}(0)}<\delta$, but $\sqrt{(y(t)-\bar{y})^{2}+v^{2}(t)} \geq$ $\epsilon$ for some $t>0$.


## Section 5.1 Answers, pp. 140-145

5.1.1 (p. 140) (c) $y=-2 e^{2 x}+e^{5 x}$ (d) $y=\left(5 k_{0}-k_{1}\right) \frac{e^{2 x}}{3}+\left(k_{1}-2 k_{0}\right) \frac{e^{5 x}}{3}$.
5.1.2 (p. 140) (c) $y=e^{x}(3 \cos x-5 \sin x)$ (d) $y=e^{x}\left(k_{0} \cos x+\left(k_{1}-k_{0}\right) \sin x\right)$
5.1.3 (p. 140) (c) $y=e^{x}(7-3 x)$ (d) $y=e^{x}\left(k_{0}+\left(k_{1}-k_{0}\right) x\right)$
5.1.4 (p. 140) (a) $y=\frac{c_{1}}{x-1}+\frac{c_{2}}{x+1}$ (b) $y=\frac{2}{x-1}-\frac{3}{x+1} ;(-1,1)$
5.1.5 (p. 141) (a) $e^{x}$ (b) $e^{2 x} \cos x$ (c) $x^{2}+2 x-2$ (d) $-\frac{5}{6} x^{-5 / 6}$ (e) $-\frac{1}{x^{2}}$ (f) $(x \ln |x|)^{2}(\mathbf{g}) \frac{e^{2 x}}{2 \sqrt{x}}$
5.1.6 (p.141) 0 5.1.7 (p. 141) $W(x)=\left(1-x^{2}\right)^{-1} 5.1 .8$ (p.141) $W(x)=\frac{1}{x} 5.1 .10$ (p. 141) $y_{2}=e^{-x}$
5.1.11 (p. 141) $y_{2}=x e^{3 x} 5.1 .12$ (p. 141) $y_{2}=x e^{a x} 5.1 .13$ (p. 141) $y_{2}=\frac{1}{x}$ 5.1.14 (p. 141) $y_{2}=x \ln x$
5.1.15 (p. 141) $y_{2}=x^{a} \ln x \quad 5.1 .16$ (p. 141) $y_{2}=x^{1 / 2} e^{-2 x}$ 5.1.17 (p. 141) $y_{2}=x$ 5.1.18 (p. 141) $y_{2}=x \sin x$ 5.1.19 (p. 141) $y_{2}=x^{1 / 2} \cos x$ 5.1.20 (p. 141) $y_{2}=x e^{-x} \quad 5.1 .21$ (p. 141) $y_{2}=\frac{1}{x^{2}-4}$
5.1.22 (p. 142) $y_{2}=e^{2 x}$
5.1.23 (p. 142) $y_{2}=x^{2} 5.1 .35$ (p. 143) (a) $y^{\prime \prime}-2 y^{\prime}+5 y=0$ (b) $(2 x-1) y^{\prime \prime}-4 x y^{\prime}+4 y=0$ (c) $x^{2} y^{\prime \prime}-x y^{\prime}+y=0$
(d) $x^{2} y^{\prime \prime}+x y^{\prime}+y=0$ (e) $y^{\prime \prime}-y=0$ (f) $x y^{\prime \prime}-y^{\prime}=0$
5.1.37 (p. 143) (c) $y=k_{0} y_{1}+k_{1} y_{2}$ 5.1.38 (p. 144) $y_{1}=1, y_{2}=x-x_{0} ; y=k_{0}+k_{1}\left(x-x_{0}\right)$
5.1.39 (p. 144) $y_{1}=\cosh \left(x-x_{0}\right), y_{2}=\sinh \left(x-x_{0}\right) ; y=k_{0} \cosh \left(x-x_{0}\right)+k_{1} \sinh \left(x-x_{0}\right)$
5.1.40 (p. 144) $y_{1}=\cos \omega\left(x-x_{0}\right), y_{2}=\frac{1}{\omega} \sin \omega\left(x-x_{0}\right) y=k_{0} \cos \omega\left(x-x_{0}\right)+\frac{k_{1}}{\omega} \sin \omega\left(x-x_{0}\right)$
5.1.41 (p. 144) $y_{1}=\frac{1}{1-x^{2}}, y_{2}=\frac{x}{1-x^{2}} y=\frac{k_{0}+k_{1} x}{1-x^{2}}$
5.1.42 (p. 144) (c) $k_{0}=k_{1}=0 ; y=\left\{\begin{array}{c}c_{1} x^{2}+c_{2} x^{3}, x \geq 0, \\ c_{1} x^{2}+c_{3} x^{3}, x<0\end{array}\right.$
(d) $(0, \infty)$ if $x_{0}>0,(-\infty, 0)$ if $x_{0}<0$
5.1.43 (p. 145) (c) $k_{0}=0, k_{1}$ arbitrary $y=k_{1} x+c_{2} x^{2}$
5.1.44 (p. 145) (c) $k_{0}=k_{1}=0 y=\left\{\begin{array}{c}a_{1} x^{3}+a_{2} x^{4}, x \geq 0, \\ b_{1} x^{3}+b_{2} x^{4}, x<0\end{array}\right.$
(d) $(0, \infty)$ if $x_{0}>0,(-\infty, 0)$ if $x_{0}<0$

Section 5.2 Answers, pp. 152-154
5.2.1(p. 152) $y=c_{1} e^{-6 x}+c_{2} e^{x} 5.2 .2$ (p. 152) $y=e^{2 x}\left(c_{1} \cos x+c_{2} \sin x\right)$ 5.2.3(p.152) $y=c_{1} e^{-7 x}+c_{2} e^{-x}$
5.2.4(p. 152) $y=e^{2 x}\left(c_{1}+c_{2} x\right)$ 5.2.5 (p. 152) $y=e^{-x}\left(c_{1} \cos 3 x+c_{2} \sin 3 x\right)$
5.2.6 (p. 152) $y=e^{-3 x}\left(c_{1} \cos x+c_{2} \sin x\right) 5.2 .7$ (p. 152) $y=e^{4 x}\left(c_{1}+c_{2} x\right) \quad 5.2 .8$ (p. 152) $y=c_{1}+c_{2} e^{-x}$
5.2.9 (p. 152) $y=e^{x}\left(c_{1} \cos \sqrt{2} x+c_{2} \sin \sqrt{2} x\right)$ 5.2.10 (p. 152) $y=e^{-3 x}\left(c_{1} \cos 2 x+c_{2} \sin 2 x\right)$
5.2.11 (p. 152) $y=e^{-x / 2}\left(c_{1} \cos \frac{3 x}{2}+c_{2} \sin \frac{3 x}{2}\right) \mathbf{5 . 2 . 1 2}$ (p. 152) $y=c_{1} e^{-x / 5}+c_{2} e^{x / 2}$
5.2.13 (p. 152) $y=e^{-7 x}(2 \cos x-3 \sin x) \mathbf{5 . 2 . 1 4}$ (p. 152) $y=4 e^{x / 2}+6 e^{-x / 3}$ 5.2.15 (p. 152) $y=$ $3 e^{x / 3}-4 e^{-x / 2}$
$\mathbf{5 . 2 . 1 6}$ (p. 152) $y=\frac{e^{-x / 2}}{3}+\frac{3 e^{3 x / 2}}{4} \mathbf{5 . 2 . 1 7}$ (p. 152) $y=e^{3 x / 2}(3-2 x) \mathbf{5 . 2 . 1 8}$ (p.152) $y=3 e^{-4 x}-4 e^{-3 x}$
5.2.19 (p. 152) $y=2 x e^{3 x} 5.2 .20$ (p.152) $y=e^{x / 6}(3+2 x) 5.2 .21$ (p.152) $y=e^{-2 x}\left(3 \cos \sqrt{6} x+\frac{2 \sqrt{6}}{3} \sin \sqrt{6} x\right)$
5.2 .23 (p. 153) $y=2 e^{-(x-1)}-3 e^{-2(x-1)} \mathbf{5 . 2 . 2 4}$ (p. 153) $y=\frac{1}{3} e^{-(x-2)}-\frac{2}{3} e^{7(x-2)}$
5.2.25 (p. 153) $y=e^{7(x-1)}(2-3(x-1)) \mathbf{5 . 2 . 2 6}\left(\mathbf{p . 1 5 3 )} y=e^{-(x-2) / 3}(2-4(x-2))\right.$
$\mathbf{5 . 2 . 2 7}$ (p. 153) $y=2 \cos \frac{2}{3}\left(x-\frac{\pi}{4}\right)-3 \sin \frac{2}{3}\left(x-\frac{\pi}{4}\right) \mathbf{5 . 2 . 2 8}$ (p.153) $y=2 \cos \sqrt{3}\left(x-\frac{\pi}{3}\right)-\frac{1}{\sqrt{3}} \sin \sqrt{3}\left(x-\frac{\pi}{3}\right)$
$\mathbf{5 . 2 . 3 0}$ (p. 153) $y=\frac{k_{0}}{r_{2}-r_{1}}\left(r_{2} e^{r_{1}\left(x-x_{0}\right)}-r_{1} e^{r_{2}\left(x-x_{0}\right)}\right)+\frac{k_{1}}{r_{2}-r_{1}}\left(e^{r_{2}\left(x-x_{0}\right)}-e^{r_{1}\left(x-x_{0}\right)}\right)$
5.2.31 (p. 153) $y=e^{r_{1}\left(x-x_{0}\right)}\left[k_{0}+\left(k_{1}-r_{1} k_{0}\right)\left(x-x_{0}\right)\right]$
5.2.32 (p. 153) $y=e^{\lambda\left(x-x_{0}\right)}\left[k_{0} \cos \omega\left(x-x_{0}\right)+\left(\frac{k_{1}-\lambda k_{0}}{\omega}\right) \sin \omega\left(x-x_{0}\right)\right]$

## Section 5.3 Answers, pp. 160-162

5.3.1 (p. 160) $y_{p}=-1+2 x+3 x^{2} ; y=-1+2 x+3 x^{2}+c_{1} e^{-6 x}+c_{2} e^{x}$
5.3.2 (p. 160) $y_{p}=1+x ; y=1+x+e^{2 x}\left(c_{1} \cos x+c_{2} \sin x\right)$
5.3.3(p. 160) $y_{p}=-x+x^{3} ; y=-x+x^{3}+c_{1} e^{-7 x}+c_{2} e^{-x}$
5.3.4 (p. 160) $y_{p}=1-x^{2} ; y=1-x^{2}+e^{2 x}\left(c_{1}+c_{2} x\right)$
5.3.5 (p. 160) $y_{p}=2 x+x^{3} ; y=2 x+x^{3}+e^{-x}\left(c_{1} \cos 3 x+c_{2} \sin 3 x\right)$;
$y=2 x+x^{3}+e^{-x}(2 \cos 3 x+3 \sin 3 x)$
5.3.6 (p. 160) $y_{p}=1+2 x ; y=1+2 x+e^{-3 x}\left(c_{1} \cos x+c_{2} \sin x\right) ; y=1+2 x+e^{-3 x}(\cos x-\sin x)$
5.3.8 (p.160) $y_{p}=\frac{2}{x} 5.3 .9\left(\mathbf{p . 1 6 0 )} y_{p}=4 x^{1 / 2} 5.3 .10\left(\mathbf{p . 1 6 0 )} y_{p}=\frac{x^{3}}{2} \quad 5.3 .11\right.\right.$ (p. 160) $y_{p}=\frac{1}{x^{3}}$
5.3.12 (p. 160) $y_{p}=9 x^{1 / 3} 5.3 .13$ (p.160) $y_{p}=\frac{2 x^{4}}{13} 5.3 .16$ (p. 160) $y_{p}=\frac{e^{3 x}}{3} ; y=\frac{e^{3 x}}{3}+c_{1} e^{-6 x}+c_{2} e^{x}$
5.3.17 (p. 160) $y_{p}=e^{2 x} ; y=e^{2 x}\left(1+c_{1} \cos x+c_{2} \sin x\right)$
5.3.18 (p. 161) $y=-2 e^{-2 x} ; y=-2 e^{-2 x}+c_{1} e^{-7 x}+c_{2} e^{-x} ; y=-2 e^{-2 x}-e^{-7 x}+e^{-x}$
5.3.19 (p. 161) $y_{p}=e^{x} ; y=e^{x}+e^{2 x}\left(c_{1}+c_{2} x\right) ; y=e^{x}+e^{2 x}(1-3 x)$
5.3.20 (p. 161) $y_{p}=\frac{4}{45} e^{x / 2} ; y=\frac{4}{45} e^{x / 2}+e^{-x}\left(c_{1} \cos 3 x+c_{2} \sin 3 x\right)$
5.3.21 (p. 161) $y_{p}=e^{-3 x} ; y=e^{-3 x}\left(1+c_{1} \cos x+c_{2} \sin x\right)$
5.3.24 (p. 161) $y_{p}=\cos x-\sin x ; y=\cos x-\sin x+e^{4 x}\left(c_{1}+c_{2} x\right)$
5.3.25 (p. 161) $y_{p}=\cos 2 x-2 \sin 2 x ; y=\cos 2 x-2 \sin 2 x+c_{1}+c_{2} e^{-x}$
5.3.26 (p. 161) $y_{p}=\cos 3 x ; y=\cos 3 x+e^{x}\left(c_{1} \cos \sqrt{2} x+c_{2} \sin \sqrt{2} x\right)$
5.3.27 (p. 161) $y_{p}=\cos x+\sin x ; y=\cos x+\sin x+e^{-3 x}\left(c_{1} \cos 2 x+c_{2} \sin 2 x\right)$
5.3.28 (p. 161) $y_{p}=-2 \cos 2 x+\sin 2 x ; y=-2 \cos 2 x+\sin 2 x+c_{1} e^{-4 x}+c_{2} e^{-3 x}$
$y=-2 \cos 2 x+\sin 2 x+2 e^{-4 x}-3 e^{-3 x}$
5.3.29 (p. 161) $y_{p}=\cos 3 x-\sin 3 x ; y=\cos 3 x-\sin 3 x+e^{3 x}\left(c_{1}+c_{2} x\right)$ $y=\cos 3 x-\sin 3 x+e^{3 x}(1+2 x)$
5.3.30 (p. 161) $y=\frac{1}{\omega_{0}^{2}-\omega^{2}}(M \cos \omega x+N \sin \omega x)+c_{1} \cos \omega_{0} x+c_{2} \sin \omega_{0} x$
5.3.33 (p. 161) $y_{p}=-1+2 x+3 x^{2}+\frac{e^{3 x}}{3} ; y=-1+2 x+3 x^{2}+\frac{e^{3 x}}{3}+c_{1} e^{-6 x}+c_{2} e^{x}$
5.3.34 (p. 161) $y_{p}=1+x+e^{2 x} ; y=1+x+e^{2 x}\left(1+c_{1} \cos x+c_{2} \sin x\right)$
5.3.35(p. 161) $y_{p}=-x+x^{3}-2 e^{-2 x} ; y=-x+x^{3}-2 e^{-2 x}+c_{1} e^{-7 x}+c_{2} e^{-x}$
5.3.36 (p. 161) $y_{p}=1-x^{2}+e^{x} ; y=1-x^{2}+e^{x}+e^{2 x}\left(c_{1}+c_{2} x\right)$
5.3.37 (p. 161) $y_{p}=2 x+x^{3}+\frac{4}{45} e^{x / 2} ; y=2 x+x^{3}+\frac{4}{45} e^{x / 2}+e^{-x}\left(c_{1} \cos 3 x+c_{2} \sin 3 x\right)$
5.3.38 (p. 161) $y_{p}=1+2 x+e^{-3 x} ; y=1+2 x+e^{-3 x}\left(1+c_{1} \cos x+c_{2} \sin x\right)$

Section 5.4 Answers, pp. 167-169
5.4.1 (p. 167) $y_{p}=e^{3 x}\left(-\frac{1}{4}+\frac{x}{2}\right) 5.4 .2$ (p. 167) $y_{p}=e^{-3 x}\left(1-\frac{x}{4}\right) 5.4 .3$ (p. 167) $y_{p}=e^{x}\left(2-\frac{3 x}{4}\right)$
5.4.4 (p. 167) $y_{p}=e^{2 x}\left(1-3 x+x^{2}\right) 5.4 .5(\mathbf{p} .167) y_{p}=e^{-x}\left(1+x^{2}\right) \mathbf{5 . 4 . 6}$ (p.167) $y_{p}=e^{x}\left(-2+x+2 x^{2}\right)$
5.4.7 (p. 167) $y_{p}=x e^{-x}\left(\frac{1}{6}+\frac{x}{2}\right) 5.4 .8\left(\mathbf{p . 1 6 7 )} y_{p}=x e^{x}(1+2 x) 5.4 .9\left(\mathbf{p . 1 6 7 )} y_{p}=x e^{3 x}\left(-1+\frac{x}{2}\right)\right.\right.$
5.4.10(p.167) $y_{p}=x e^{2 x}(-2+x) 5.4 .11$ (p.167) $y_{p}=x^{2} e^{-x}\left(1+\frac{x}{2}\right)$ 5.4.12 (p.167) $y_{p}=x^{2} e^{x}\left(\frac{1}{2}-x\right)$
5.4.13(p. 167) $y_{p}=\frac{x^{2} e^{2 x}}{2}\left(1-x+x^{2}\right) \mathbf{5 . 4 . 1 4}$ (p. 167) $y_{p}=\frac{x^{2} e^{-x / 3}}{27}\left(3-2 x+x^{2}\right)$
5.4.15 (p. 167) $y=\frac{e^{3 x}}{4}(-1+2 x)+c_{1} e^{x}+c_{2} e^{2 x} \quad 5.4 .16$ (p. 167) $y=e^{x}(1-2 x)+c_{1} e^{2 x}+c_{2} e^{4 x}$
5.4.17 (p. 167) $y=\frac{e^{2 x}}{5}(1-x)+e^{-3 x}\left(c_{1}+c_{2} x\right) 5.4 .18$ (p. 167) $y=x e^{x}(1-2 x)+c_{1} e^{x}+c_{2} e^{-3 x}$
5.4.19 (p. 167) $y=e^{x}\left[x^{2}(1-2 x)+c_{1}+c_{2} x\right] \quad 5.4 .20\left(\right.$ p. 167) $y=-e^{2 x}(1+x)+2 e^{-x}-e^{5 x}$
5.4.21(p.167) $y=x e^{2 x}+3 e^{x}-e^{-4 x}$ 5.4.22 (p. 167) $y=e^{-x}\left(2+x-2 x^{2}\right)-e^{-3 x}$
5.4.23 (p. 167) $y=e^{-2 x}(3-x)-2 e^{5 x} \mathbf{5 . 4 . 2 4}$ (p. 167) $y_{p}=-\frac{e^{x}}{3}(1-x)+e^{-x}(3+2 x)$
5.4 .25 (p. 167) $y_{p}=e^{x}(3+7 x)+x e^{3 x} 5.4 .26$ (p. 167) $y_{p}=x^{3} e^{4 x}+1+2 x+x^{2}$
5.4 .27 (p. 167) $y_{p}=x e^{2 x}(1-2 x)+x e^{x} \quad 5.4 .28$ (p. 167) $y_{p}=e^{x}(1+x)+x^{2} e^{-x}$
5.4 .29 (p. 167) $y_{p}=x^{2} e^{-x}+e^{3 x}\left(1-x^{2}\right) \mathbf{5 . 4 . 3 1}$ (p. 168) $y_{p}=2 e^{2 x} \quad 5.4 .32$ (p. 168) $y_{p}=5 x e^{4 x}$
5.4 .33 (p. 168) $y_{p}=x^{2} e^{4 x} 5.4 .34$ (p. 168) $y_{p}=-\frac{e^{3 x}}{4}\left(1+2 x-2 x^{2}\right) 5.4 .35$ (p.168) $y_{p}=x e^{3 x}\left(4-x+2 x^{2}\right)$
5.4.36 (p. 168) $y_{p}=x^{2} e^{-x / 2}\left(-1+2 x+3 x^{2}\right)$
5.4.37 (p. 168) (a) $y=e^{-x}\left(\frac{4}{3} x^{3 / 2}+c_{1} x+c_{2}\right)$ (b) $y=e^{-3 x}\left[\frac{x^{2}}{4}(2 \ln x-3)+c_{1} x+c_{2}\right]$
(c) $y=e^{2 x}\left[(x+1) \ln |x+1|+c_{1} x+c_{2}\right]$ (d) $y=e^{-x / 2}\left(x \ln |x|+\frac{x^{3}}{6}+c_{1} x+c_{2}\right)$
5.4.39 (p. 169) (a) $e^{x}(3+x)+c$ (b) $-e^{-x}(1+x)^{2}+c(\mathbf{c})-\frac{e^{-2 x}}{8}\left(3+6 x+6 x^{2}+4 x^{3}\right)+c$

$$
\text { (d) } e^{x}\left(1+x^{2}\right)+c \text { (e) } e^{3 x}\left(-6+4 x+9 x^{2}\right)+c(\mathbf{f})-e^{-x}\left(1-2 x^{3}+3 x^{4}\right)+c
$$

5.4.40 (p. 169) $\frac{(-1)^{k} k!e^{\alpha x}}{\alpha^{k+1}} \sum_{r=0}^{k} \frac{(-\alpha x)^{r}}{r!}+c$

## Section 5.5 Answers, pp. 174-178

5.5.1 (p. 174) $y_{p}=\cos x+2 \sin x$ 5.5.2 (p. 175) $y_{p}=\cos x+(2-2 x) \sin x$
5.5.3 (p. 175) $y_{p}=e^{x}(-2 \cos x+3 \sin x)$
5.5.4 (p. 175) $y_{p}=\frac{e^{2 x}}{2}(\cos 2 x-\sin 2 x)$ 5.5.5 (p. 175) $y_{p}=-e^{x}(x \cos x-\sin x)$
5.5.6 (p. 175) $y_{p}=e^{-2 x}(1-2 x)(\cos 3 x-\sin 3 x)$ 5.5.7 (p. 175) $y_{p}=x(\cos 2 x-3 \sin 2 x)$
5.5.8 (p. 175) $y_{p}=-x[(2-x) \cos x+(3-2 x) \sin x] 5.5 .9$ (p. 175) $y_{p}=x\left[x \cos \left(\frac{x}{2}\right)-3 \sin \left(\frac{x}{2}\right)\right]$
5.5.10 (p. 175) $y_{p}=x e^{-x}(3 \cos x+4 \sin x)$ 5.5.11 (p. 175) $y_{p}=x e^{x}[(-1+x) \cos 2 x+(1+x) \sin 2 x]$
5.5.12 (p. 175) $y_{p}=-(14-10 x) \cos x-\left(2+8 x-4 x^{2}\right) \sin x$.
5.5.13 (p.175) $y_{p}=\left(1+2 x+x^{2}\right) \cos x+\left(1+3 x^{2}\right) \sin x$ 5.5.14 (p. 175) $y_{p}=\frac{x^{2}}{2}(\cos 2 x-\sin 2 x)$
5.5.15 (p. 175) $y_{p}=e^{x}\left(x^{2} \cos x+2 \sin x\right)$ 5.5.16 (p. 175) $y_{p}=e^{x}\left(1-x^{2}\right)(\cos x+\sin x)$
5.5.17 (p. 175) $y_{p}=e^{x}\left(x^{2}-x^{3}\right)(\cos x+\sin x) \mathbf{5 . 5 . 1 8}$ (p. 175) $y_{p}=e^{-x}[(1+2 x) \cos x-(1-3 x) \sin x]$
5.5.19 (p. 175) $y_{p}=x(2 \cos 3 x-\sin 3 x)$ 5.5.20 (p. 175) $y_{p}=-x^{3} \cos x+\left(x+2 x^{2}\right) \sin x$
5.5.21 (p. 175) $y_{p}=-e^{-x}\left[\left(x+x^{2}\right) \cos x-(1+2 x) \sin x\right]$
5.5.22 (p. 175) $y=e^{x}(2 \cos x+3 \sin x)+3 e^{x}-e^{6 x} 5.5 .23$ (p. 175) $y=e^{x}[(1+2 x) \cos x+(1-3 x) \sin x]$
5.5.24(p.175) $y=e^{x}(\cos x-2 \sin x)+e^{-3 x}(\cos x+\sin x) 5.5 .25$ (p.175) $y=e^{3 x}[(2+2 x) \cos x-(1+3 x) \sin x]$
5.5.26 (p.175) $y=e^{3 x}[(2+3 x) \cos x+(4-x) \sin x]+3 e^{x}-5 e^{2 x} 5.5 .27$ (p. 175) $y_{p}=x e^{3 x}-\frac{e^{x}}{5}(\cos x-2 \sin x)$
5.5.28 (p. 175) $y_{p}=x(\cos x+2 \sin x)-\frac{e^{x}}{2}(1-x)+\frac{e^{-x}}{2}$
5.5.29 (p. 175) $y_{p}=-\frac{x e^{x}}{2}(2+x)+2 x e^{2 x}+\frac{1}{10}(3 \cos x+\sin x)$
5.5.30 (p. 175) $y_{p}=x e^{x}(\cos x+x \sin x)+\frac{e^{-x}}{25}(4+5 x)+1+x+\frac{x^{2}}{2}$
5.5.31 (p. 175) $y_{p}=\frac{x^{2} e^{2 x}}{6}(3+x)-e^{2 x}(\cos x-\sin x)+3 e^{3 x}+\frac{1}{4}(2+x)$
5.5.32 (p. 175) $y=\left(1-2 x+3 x^{2}\right) e^{2 x}+4 \cos x+3 \sin x$ 5.5.33 (p. 175) $y=x e^{-2 x} \cos x+3 \cos 2 x$
5.5.34 (p. 175) $y=-\frac{3}{8} \cos 2 x+\frac{1}{4} \sin 2 x+e^{-x}-\frac{13}{8} e^{-2 x}-\frac{3}{4} x e^{-2 x}$
5.5.40 (p. 178) (a) $2 x \cos x-\left(2-x^{2}\right) \sin x+c$ (b) $-\frac{e^{x}}{2}\left[\left(1-x^{2}\right) \cos x-(1-x)^{2} \sin x\right]+c$
(c) $-\frac{e^{-x}}{25}[(4+10 x) \cos 2 x-(3-5 x) \sin 2 x]+c$
(d) $-\frac{e^{-x}}{2}\left[(1+x)^{2} \cos x-\left(1-x^{2}\right) \sin x\right]+c$
(e) $-\frac{e^{x}}{2}\left[x\left(3-3 x+x^{2}\right) \cos x-\left(3-3 x+x^{3}\right) \sin x\right]+c$
(f) $-e^{x}[(1-2 x) \cos x+(1+x) \sin x]+c(\mathbf{g}) e^{-x}[x \cos x+x(1+x) \sin x]+c$

## Section 5.7 Answers, pp. 184-186

5.7.1 (p. 184) $y_{p}=\frac{-\cos 3 x \ln |\sec 3 x+\tan 3 x|}{9} \mathbf{5 . 7 . 2}$ (p. 184) $y_{p}=-\frac{\sin 2 x \ln |\cos 2 x|}{4}+\frac{x \cos 2 x}{2}$
5.7.3 (p. 184) $y_{p}=4 e^{x}\left(1+e^{x}\right) \ln \left(1+e^{-x}\right) 5.7 .4$ (p. 184) $y_{p}=3 e^{x}(\cos x \ln |\cos x|+x \sin x)$
5.7.5 (p. 184) $y_{p}=\frac{8}{5} x^{7 / 2} e^{x} 5.7 .6$ (p. 184) $y_{p}=e^{x} \ln \left(1-e^{-2 x}\right)-e^{-x} \ln \left(e^{2 x}-1\right) 5.7 .7$ (p. 184) $y_{p}=$ $\frac{2\left(x^{2}-3\right)}{3}$
5.7 .8 (p. 184) $y_{p}=\frac{e^{2 x}}{x} 5.7 .9$ (p. 184) $y_{p}=x^{1 / 2} e^{x} \ln x 5.7 .10$ (p. 184) $y_{p}=e^{-x(x+2)}$
5.7.11 (p. 184) $y_{p}=-4 x^{5 / 2} 5.7 .12$ (p. 184) $y_{p}=-2 x^{2} \sin x-2 x \cos x 5.7 .13$ (p.184) $y_{p}=-\frac{x e^{-x}(x+1)}{2}$
5.7 .14 (p. 184) $y_{p}=-\frac{\sqrt{x} \cos \sqrt{x}}{2} 5.7 .15$ (p. 184) $y_{p}=\frac{3 x^{4} e^{x}}{2} 5.7 .16$ (p. 184) $y_{p}=x^{a+1}$
5.7.17 (p. 184) $y_{p}=\frac{x^{2} \sin x}{2} 5.7 .18$ (p. 184) $y_{p}=-2 x^{2} 5.7 .19$ (p. 185) $y_{p}=-e^{-x} \sin x$
5.7 .20 (p. 185) $y_{p}=-\frac{\sqrt{x}}{2} 5.7 .21$ (p. 185) $y_{p}=\frac{x^{3 / 2}}{4} 5.7 .22$ (p. 185) $y_{p}=-3 x^{2}$
5.7.23 (p. 185) $y_{p}=\frac{x^{3} e^{x}}{2} 5.7 .24$ (p. 185) $y_{p}=-\frac{4 x^{3 / 2}}{15} \quad 5.7 .25$ (p. 185) $y_{p}=x^{3} e^{x} \quad 5.7 .26$ (p. 185)
$y_{p}=x e^{x}$
5.7.27(p. 185) $y_{p}=x^{2} 5.7 .28\left(\mathbf{p . 1 8 5 )} y_{p}=x e^{x}(x-2) 5.7 .29(p .185) y_{p}=\sqrt{x} e^{x}(x-1) / 4\right.$
5.7.30 (p. 185) $y=\frac{e^{2 x}\left(3 x^{2}-2 x+6\right)}{6}+\frac{x e^{-x}}{3} \mathbf{5 . 7 . 3 1}$ (p. 185) $y=(x-1)^{2} \ln (1-x)+2 x^{2}-5 x+3$
5.7.32(p.185) $y=\left(x^{2}-1\right) e^{x}-5(x-1) \mathbf{5 . 7 . 3 3}(\mathbf{p . 1 8 5}) y=\frac{x\left(x^{2}+6\right)}{3\left(x^{2}-1\right)} \mathbf{5 . 7 . 3 4}$ (p.185) $y=-\frac{x^{2}}{2}+x+\frac{1}{2 x^{2}}$
5.7.35 (p. 185) $y=\frac{x^{2}(4 x+9)}{6(x+1)}$
5.7.38 (p. 186) (a) $y=k_{0} \cosh x+k_{1} \sinh x+\int_{0}^{x} \sinh (x-t) f(t) d t$
(b) $y^{\prime}=k_{0} \sinh x+k_{1} \cosh x+\int_{0}^{x} \cosh (x-t) f(t) d t$
5.7.39 (p. 186) (a) $y(x)=k_{0} \cos x+k_{1} \sin x+\int_{0}^{x} \sin (x-t) f(t) d t$ (b) $y^{\prime}(x)=-k_{0} \sin x+k_{1} \cos x+\int_{0}^{x} \cos (x-t) f(t) d t$

Section 6.1 Answers, pp. 196-197
6.1.1 (p. 196) $y=3 \cos 4 \sqrt{6} t-\frac{1}{2 \sqrt{6}} \sin 4 \sqrt{6} t \mathrm{ft} 6.1 .2$ (p. 196) $y=-\frac{1}{4} \cos 8 \sqrt{5} t-\frac{1}{4 \sqrt{5}} \sin 8 \sqrt{5} t \mathrm{ft}$
6.1.3 (p. 196) $y=1.5 \cos 14 \sqrt{10} t \mathrm{~cm}$
6.1.4 (p. 196) $y=\frac{1}{4} \cos 8 t-\frac{1}{16} \sin 8 t \mathrm{ft} ; R=\frac{\sqrt{17}}{16} \mathrm{ft} ; \omega_{0}=8 \mathrm{rad} / \mathrm{s} ; T=\pi / 4 \mathrm{~s}$; $\phi \approx-.245 \mathrm{rad} \approx-14.04^{\circ} ;$
6.1.5 (p. 196) $y=10 \cos 14 t+\frac{25}{14} \sin 14 t \mathrm{~cm} ; R=\frac{5}{14} \sqrt{809} \mathrm{~cm} ; \omega_{0}=14 \mathrm{rad} / \mathrm{s} ; T=\pi / 7 \mathrm{~s}$; $\phi \approx .177 \mathrm{rad} \approx 10.12^{\circ}$
6.1.6 (p. 196) $y=-\frac{1}{4} \cos \sqrt{70} t+\frac{2}{\sqrt{70}} \sin \sqrt{70} t \mathrm{~m} ; R=\frac{1}{4} \sqrt{\frac{67}{35}} \mathrm{~m} \omega_{0}=\sqrt{70} \mathrm{rad} / \mathrm{s}$;

$$
T=2 \pi / \sqrt{70} \mathrm{~s} ; \phi \approx 2.38 \mathrm{rad} \approx 136.28^{\circ}
$$

6.1.7 (p. 196) $y=\frac{2}{3} \cos 16 t-\frac{1}{4} \sin 16 t \mathrm{ft} \mathbf{6 . 1 . 8}$ (p. 196) $y=\frac{1}{2} \cos 8 t-\frac{3}{8} \sin 8 t \mathrm{ft} \mathbf{6 . 1 . 9}$ (p. 196) .72 m
6.1.10 (p. 196) $y=\frac{1}{3} \sin t+\frac{1}{2} \cos 2 t+\frac{5}{6} \sin 2 t \mathrm{ft} 6.1 .11$ (p. 197) $y=\frac{16}{5}\left(4 \sin \frac{t}{4}-\sin t\right)$
6.1.12 (p. 197) $y=-\frac{1}{16} \sin 8 t+\frac{1}{3} \cos 4 \sqrt{2} t-\frac{1}{8 \sqrt{2}} \sin 4 \sqrt{2} t$
6.1.13 (p. 197) $y=-t \cos 8 t-\frac{1}{6} \cos 8 t+\frac{1}{8} \sin 8 t \mathrm{ft} \mathbf{6 . 1} 14$ (p. 197) $T=4 \sqrt{2} \mathrm{~s}$
6.1.15 (p. 197) $\omega=8 \mathrm{rad} / \mathrm{s} y=-\frac{t}{16}(-\cos 8 t+2 \sin 8 t)+\frac{1}{128} \sin 8 t \mathrm{ft}$
6.1.16 (p. 197) $\omega=4 \sqrt{6} \mathrm{rad} / \mathrm{s} ; \quad y=-\frac{t}{\sqrt{6}}\left[\frac{8}{3} \cos 4 \sqrt{6} t+4 \sin 4 \sqrt{6} t\right]+\frac{1}{9} \sin 4 \sqrt{6} t \mathrm{ft}$
6.1.17 (p. 197) $y=\frac{t}{2} \cos 2 t-\frac{t}{4} \sin 2 t+3 \cos 2 t+2 \sin 2 t \mathrm{~m}$
6.1.18 (p. 197) $y=y_{0} \cos \omega_{0} t+\frac{v_{0}}{\omega_{0}} \sin \omega_{0} t ; R=\frac{1}{\omega_{0}} \sqrt{\left(\omega_{0} y_{0}\right)^{2}+\left(v_{0}\right)^{2}}$;

$$
\cos \phi=\frac{y_{0} \omega_{0}}{\sqrt{\left(\omega_{0} y_{0}\right)^{2}+\left(v_{0}\right)^{2}}} ; \sin \phi=\frac{\omega_{0}}{\sqrt{\left(\omega_{0} y_{0}\right)^{2}+\left(v_{0}\right)^{2}}}
$$

6.1.19 (p. 197) The object with the longer period weighs four times as much as the other.
6.1.20 (p. 197) $T_{2}=\sqrt{2} T_{1}$, where $T_{1}$ is the period of the smaller object.
6.1.21 (p. 197) $k_{1}=9 k_{2}$, where $k_{1}$ is the spring constant of the system with the shorter period.

Section 6.2 Answers, pp. 204-205
6.2.1 (p. 204) $y=\frac{e^{-2 t}}{2}(3 \cos 2 t-\sin 2 t) \mathrm{ft} ; \sqrt{\frac{5}{2}} e^{-2 t} \mathrm{ft}$
6.2.2 (p. 204) $y=-e^{-t}\left(3 \cos 3 t+\frac{1}{3} \sin 3 t\right) \mathrm{ft} \frac{\sqrt{82}}{3} e^{-t} \mathrm{ft}$
6.2.3 (p. 204) $y=e^{-16 t}\left(\frac{1}{4}+10 t\right) \mathrm{ft} \mathbf{6 . 2 . 4}\left(\mathbf{p}\right.$. 204) $y=-\frac{e^{-3 t}}{4}(5 \cos t+63 \sin t) \mathrm{ft}$
6.2.5 (p. 204) $0 \leq c<8 \mathrm{lb}-\mathrm{sec} / \mathrm{ft} 6.2 .6$ (p. 204) $y=\frac{1}{2} e^{-3 t}\left(\cos \sqrt{91} t+\frac{11}{\sqrt{91}} \sin \sqrt{91} t\right) \mathrm{ft}$
6.2.7 (p. 204) $y=-\frac{e^{-4 t}}{3}(2+8 t)$ ft 6.2 .8 (p. 204) $y=e^{-10 t}\left(9 \cos 4 \sqrt{6} t+\frac{45}{2 \sqrt{6}} \sin 4 \sqrt{6} t\right) \mathrm{cm}$
6.2.9 (p. 204) $y=e^{-3 t / 2}\left(\frac{3}{2} \cos \frac{\sqrt{41}}{2} t+\frac{9}{2 \sqrt{41}} \sin \frac{\sqrt{41}}{2} t\right) \mathrm{ft}$
6.2.10 (p. 204) $y=e^{-\frac{3}{2} t}\left(\frac{1}{2} \cos \frac{\sqrt{119}}{2} t-\frac{9}{2 \sqrt{119}} \sin \frac{\sqrt{119}}{2} t\right) \mathrm{ft}$
6.2.11 (p. 204) $y=e^{-8 t}\left(\frac{1}{4} \cos 8 \sqrt{2} t-\frac{1}{4 \sqrt{2}} \sin 8 \sqrt{2} t\right) \mathrm{ft}$
6.2.12 (p. 204) $y=e^{-t}\left(-\frac{1}{3} \cos 3 \sqrt{11} t+\frac{14}{9 \sqrt{11}} \sin 3 \sqrt{11} t\right) \mathrm{ft}$
6.2.13 (p. 204) $y_{p}=\frac{22}{61} \cos 2 t+\frac{2}{61} \sin 2 t \mathrm{ft} 6.2 .14$ (p. 204) $y=-\frac{2}{3}\left(e^{-8 t}-2 e^{-4 t}\right)$
6.2.15 (p. 204) $y=e^{-2 t}\left(\frac{1}{10} \cos 4 t-\frac{1}{5} \sin 4 t\right) \mathrm{m} 6.2 .16$ (p. 204) $y=e^{-3 t}(10 \cos t-70 \sin t) \mathrm{cm}$
6.2.17 (p. 205) $y_{p}=-\frac{2}{15} \cos 3 t+\frac{1}{15} \sin 3 t \mathrm{ft}$
6.2.18 (p. 205) $y_{p}=\frac{11}{100} \cos 4 t+\frac{27}{100} \sin 4 t \mathrm{~cm} 6.2 .19$ (p. 205) $y_{p}=\frac{42}{73} \cos t+\frac{39}{73} \sin t \mathrm{ft}$
6.2.20 (p. 205) $y=-\frac{1}{2} \cos 2 t+\frac{1}{4} \sin 2 t \mathrm{~m} 6.2 .21$ (p.205) $y_{p}=\frac{1}{c \omega_{0}}\left(-\beta \cos \omega_{0} t+\alpha \sin \omega_{0} t\right)$
6.2.24 (p. 205) $y=e^{-c t / 2 m}\left(y_{0} \cos \omega_{1} t+\frac{1}{\omega_{1}}\left(v_{0}+\frac{c y_{0}}{2 m}\right) \sin \omega_{1} t\right)$
6.2.25 (p. 205) $y=\frac{r_{2} y_{0}-v_{0}}{r_{2}-r_{1}} e^{r_{1} t}+\frac{v_{0}-r_{1} y_{0}}{r_{2}-r_{1}} e^{r_{2} t} 6.2 .26$ (p. 205) $y=e^{r_{1} t}\left(y_{0}+\left(v_{0}-r_{1} y_{0}\right) t\right)$

## Section 7.1 Answers, pp. 216-218

7.1.1 (p. 216) (a) $R=2 ; I=(-1,3) ;$ (b) $R=1 / 2 ; I=(3 / 2,5 / 2)$ (c) $R=0$; (d) $R=16$; $I=(-14,18)($ e) $R=\infty ; I=(-\infty, \infty)(\mathbf{f}) R=4 / 3 ; I=(-25 / 3,-17 / 3)$
7.1.3 (p. 216) (a) $R=1 ; I=(0,2)$ (b) $R=\sqrt{2} ; I=(-2-\sqrt{2},-2+\sqrt{2})$; (c) $R=\infty$; $I=(-\infty, \infty)(\mathbf{d}) R=0$ (e) $R=\sqrt{3} ; I=(-\sqrt{3}, \sqrt{3})(\mathbf{f}) R=1 I=(0,2)$
7.1.5 (p. 217) (a) $R=3 ; I=(0,6)$ (b) $R=1 ; I=(-1,1)$ (c) $R=1 / \sqrt{3}$ $I=(3-1 / \sqrt{3}, 3+1 / \sqrt{3})(\mathbf{d}) R=\infty ; I=(-\infty, \infty)($ e) $R=0$ (f) $R=2$; $I=(-1,3)$
7.1.11 (p. 217) $b_{n}=2(n+2)(n+1) a_{n+2}+(n+1) n a_{n+1}+(n+3) a_{n}$
7.1.12 (p. 217) $b_{0}=2 a_{2}-2 a_{0} b_{n}=(n+2)(n+1) a_{n+2}+[3 n(n-1)-2] a_{n}+3(n-1) a_{n-1}, n \geq 1$
7.1.13 (p. 217) $b_{n}=(n+2)(n+1) a_{n+2}+2(n+1) a_{n+1}+\left(2 n^{2}-5 n+4\right) a_{n}$
7.1.14 (p. 217) $b_{n}=(n+2)(n+1) a_{n+2}+2(n+1) a_{n+1}+\left(n^{2}-2 n+3\right) a_{n}$
7.1.15 (p. 217) $b_{n}=(n+2)(n+1) a_{n+2}+\left(3 n^{2}-5 n+4\right) a_{n}$
7.1.16 (p. 217) $b_{0}=-2 a_{2}+2 a_{1}+a_{0}$, $b_{n}=-(n+2)(n+1) a_{n+2}+(n+1)(n+2) a_{n+1}+(2 n+1) a_{n}+a_{n-1}, n \geq 2$
7.1.17 (p. 218) $b_{0}=8 a_{2}+4 a_{1}-6 a_{0}$, $b_{n}=4(n+2)(n+1) a_{n+2}+4(n+1)^{2} a_{n+1}+\left(n^{2}+n-6\right) a_{n}-3 a_{n-1}, n \geq 1$
7.1.21 (p. 218) $b_{0}=(r+1)(r+2) a_{0}$, $b_{n}=(n+r+1)(n+r+2) a_{n}-(n+r-2)^{2} a_{n-1}, n \geq 1$.
7.1.22 (p. 218) $b_{0}=(r-2)(r+2) a_{0}$, $b_{n}=(n+r-2)(n+r+2) a_{n}+(n+r+2)(n+r-3) a_{n-1}, n \geq 14$
7.1.23 (p. 218) $b_{0}=(r-1)^{2} a_{0}, b_{1}=r^{2} a_{1}+(r+2)(r+3) a_{0}$, $b_{n}=(n+r-1)^{2} a_{n}+(n+r+1)(n+r+2) a_{n-1}+(n+r-1) a_{n-2}, n \geq 2$
7.1.24 (p. 218) $b_{0}=r(r+1) a_{0}, b_{1}=(r+1)(r+2) a_{1}+3(r+1)(r+2) a_{0}$, $b_{n}=(n+r)(n+r+1) a_{n}+3(n+r)(n+r+1) a_{n-1}+(n+r) a_{n-2}, n \geq 2$
7.1.25 (p. 218) $b_{0}=(r+2)(r+1) a_{0} \quad b_{1}=(r+3)(r+2) a_{1}$, $b_{n}=(n+r+2)(n+r+1) a_{n}+2(n+r-1)(n+r-3) a_{n-2}, n \geq 2$
7.1.26 (p. 218) $b_{0}=2(r+1)(r+3) a_{0}, b_{1}=2(r+2)(r+4) a_{1}$, $b_{n}=2(n+r+1)(n+r+3) a_{n}+(n+r-3)(n+r) a_{n-2}, n \geq 2$
Section 7.2 Answers, pp. 227-231
7.2.1 (p. 227) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m}(2 m+1) x^{2 m}+a_{1} \sum_{m=0}^{\infty}(-1)^{m}(m+1) x^{2 m+1}$
7.2.2 (p. 227) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m+1} \frac{x^{2 m}}{2 m-1}+a_{1} x$
7.2.3 (p. 227) $y=a_{0}\left(1-10 x^{2}+5 x^{4}\right)+a_{1}\left(x-2 x^{3}+\frac{1}{5} x^{5}\right)$
7.2.4 (p. 227) $y=a_{0} \sum_{m=0}^{\infty}(m+1)(2 m+1) x^{2 m}+\frac{a_{1}}{3} \sum_{m=0}^{\infty}(m+1)(2 m+3) x^{2 m+1}$
7.2.5 (p. 227) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{4 j+1}{2 j+1}\right] x^{2 m}+a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1}(4 j+3)\right] \frac{x^{2 m+1}}{2^{m} m!}$
7.2.6 (p. 227) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{(4 j+1)^{2}}{2 j+1}\right] \frac{x^{2 m}}{8^{m} m!}+a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{(4 j+3)^{2}}{2 j+3}\right] \frac{x^{2 m+1}}{8^{m} m!}$
7.2.7 (p. 227) $y=a_{0} \sum_{m=0}^{\infty} \frac{2^{m} m!}{\prod_{j=0}^{m-1}(2 j+1)} x^{2 m}+a_{1} \sum_{m=0}^{\infty} \frac{\prod_{j=0}^{m-1}(2 j+3)}{2^{m} m!} x^{2 m+1}$
7.2.8 (p. 227) $y=a_{0}\left(1-14 x^{2}+\frac{35}{3} x^{4}\right)+a_{1}\left(x-3 x^{3}+\frac{3}{5} x^{5}+\frac{1}{35} x^{7}\right)$
7.2 .9 (p. 227) (a) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{2 m}}{\prod_{j=0}^{m-1}(2 j+1)}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{2 m+1}}{2^{m} m!}$
7.2 .10 (p. 227) (a) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{4 j+3}{2 j+1}\right] \frac{x^{2 m}}{2^{m} m!}+a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1} \frac{4 j+5}{2 j+3}\right] \frac{x^{2 m+1}}{2^{m} m!}$
7.2.11 (p. 227) $y=2-x-x^{2}+\frac{1}{3} x^{3}+\frac{5}{12} x^{4}-\frac{1}{6} x^{5}-\frac{17}{72} x^{6}+\frac{13}{126} x^{7}+\cdots$
7.2.12 (p. 227) $y=1-x+3 x^{2}-\frac{5}{2} x^{3}+5 x^{4}-\frac{21}{8} x^{5}+3 x^{6}-\frac{11}{16} x^{7}+\cdots$
7.2 .13 (p. 227) $y=2-x-2 x^{2}+\frac{1}{3} x^{3}+3 x^{4}-\frac{5}{6} x^{5}-\frac{49}{5} x^{6}+\frac{45}{14} x^{7}+\cdots$
7.2.16 (p. 228) $y=a_{0} \sum_{m=0}^{\infty} \frac{(x-3)^{2 m}}{(2 m)!}+a_{1} \sum_{m=0}^{\infty} \frac{(x-3)^{2 m+1}}{(2 m+1)!}$
7.2 .17 (p. 228) $y=a_{0} \sum_{m=0}^{\infty} \frac{(x-3)^{2 m}}{2^{m} m!}+a_{1} \sum_{m=0}^{\infty} \frac{(x-3)^{2 m+1}}{\prod_{j=0}^{m-1}(2 j+3)}$
7.2.18 (p. 228) $y=a_{0} \sum_{m=0}^{\infty}\left[\prod_{j=0}^{m-1}(2 j+3)\right] \frac{(x-1)^{2 m}}{m!}+a_{1} \sum_{m=0}^{\infty} \frac{4^{m}(m+1)!}{\prod_{j=0}^{m-1}(2 j+3)}(x-1)^{2 m+1}$
7.2 .19 (p. 228) $y=a_{0}\left(1-6(x-2)^{2}+\frac{4}{3}(x-2)^{4}+\frac{8}{135}(x-2)^{6}\right)+a_{1}\left((x-2)-\frac{10}{9}(x-2)^{3}\right)$
7.2 .20 (p.228) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1}(2 j+1)\right] \frac{3^{m}}{4^{m} m!}(x+1)^{2 m}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{3^{m} m!}{\prod_{j=0}^{m-1}(2 j+3)}(x+1)^{2 m+1}$
7.2.21 (p. 228) $y=-1+2 x+\frac{3}{8} x^{2}-\frac{1}{3} x^{3}-\frac{3}{128} x^{4}-\frac{1}{1024} x^{6}+\cdots$
7.2 .22 (p. 228) $y=-2+3(x-3)+3(x-3)^{2}-2(x-3)^{3}-\frac{5}{4}(x-3)^{4}+\frac{3}{5}(x-3)^{5}+\frac{7}{24}(x-3)^{6}-\frac{4}{35}(x-3)^{7}+\cdots$
$7.2 .23(\mathbf{p . 2 2 8}) y=-1+(x-1)+3(x-1)^{2}-\frac{5}{2}(x-1)^{3}-\frac{27}{4}(x-1)^{4}+\frac{21}{4}(x-1)^{5}+\frac{27}{2}(x-1)^{6}-\frac{81}{8}(x-1)^{7}+\cdots$
$7.2 .24(\mathbf{p . 2 2 8}) y=4-6(x-3)-2(x-3)^{2}+(x-3)^{3}+\frac{3}{2}(x-3)^{4}-\frac{5}{4}(x-3)^{5}-\frac{49}{20}(x-3)^{6}+\frac{135}{56}(x-3)^{7}+\cdots$
$7.2 .25\left(\mathbf{p . 2 2 8 )} y=3-4(x-4)+15(x-4)^{2}-4(x-4)^{3}+\frac{15}{4}(x-4)^{4}-\frac{1}{5}(x-4)^{5}\right.$
7.2 .26 (p. 228) $y=3-3(x+1)-30(x+1)^{2}+\frac{20}{3}(x+1)^{3}+20(x+1)^{4}-\frac{4}{3}(x+1)^{5}-\frac{8}{9}(x+1)^{6}$
7.2 .27 (p. 228) (a) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} x^{2 m}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} x^{2 m+1}$ (b) $y=\frac{a_{0}+a_{1} x}{1+x^{2}}$
7.2.33 (p. 230) $y=a_{0} \sum_{m=0}^{\infty} \frac{x^{3 m}}{3^{m} m!\prod_{j=0}^{m-1}(3 j+2)}+a_{1} \sum_{m=0}^{\infty} \frac{x^{3 m+1}}{3^{m} m!\prod_{j=0}^{m-1}(3 j+4)}$
7.2.34 (p. 230) $y=a_{0} \sum_{m=0}^{\infty}\left(\frac{2}{3}\right)^{m}\left[\prod_{j=0}^{m-1}(3 j+2)\right] \frac{x^{3 m}}{m!}+a_{1} \sum_{m=0}^{\infty} \frac{6^{m} m!}{\prod_{j=0}^{m-1}(3 j+4)} x^{3 m+1}$
7.2 .35 (p. 230) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{3^{m} m!}{\prod_{j=0}^{m-1}(3 j+2)} x^{3 m}+a_{1} \sum_{m=0}^{\infty}(-1)^{m}\left[\prod_{j=0}^{m-1}(3 j+4)\right] \frac{x^{3 m+1}}{3^{m} m!}$
7.2.36 (p. 230) $y=a_{0}\left(1-4 x^{3}+4 x^{6}\right)+a_{1} \sum_{m=0}^{\infty} 2^{m}\left[\prod_{j=0}^{m-1} \frac{3 j-5}{3 j+4}\right] x^{3 m+1}$
7.2 .37 (p. 231) $y=a_{0}\left(1+\frac{21}{2} x^{3}+\frac{42}{5} x^{6}+\frac{7}{20} x^{9}\right)+a_{1}\left(x+4 x^{4}+\frac{10}{7} x^{7}\right)$
7.2 .39 (p. 231) $y=a_{0} \sum_{m=0}^{\infty}(-2)^{m}\left[\prod_{j=0}^{m-1} \frac{5 j+1}{5 j+4}\right] x^{5 m}+a_{1} \sum_{m=0}^{\infty}\left(-\frac{2}{5}\right)^{m}\left[\prod_{j=0}^{m-1}(5 j+2)\right] \frac{x^{5 m+1}}{m!}$
7.2 .40 (p. 231) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{4 m}}{4^{m} m!\prod_{j=0}^{m-1}(4 j+3)}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{4 m+1}}{4^{m} m!\prod_{j=0}^{m-1}(4 j+5)}$
7.2 .41 (p. 231) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{7 m}}{\prod_{j=0}^{m-1}(7 j+6)}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{7 m+1}}{7^{m} m!}$
7.2.42 (p. 231) $y=a_{0}\left(1-\frac{9}{7} x^{8}\right)+a_{1}\left(x-\frac{7}{9} x^{9}\right)$
7.2.43 (p. 231) $y=a_{0} \sum_{m=0}^{\infty} x^{6 m}+a_{1} \sum_{m=0}^{\infty} x^{6 m+1}$
7.2.44 (p. 231) $y=a_{0} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{6 m}}{\prod_{j=0}^{m-1}(6 j+5)}+a_{1} \sum_{m=0}^{\infty}(-1)^{m} \frac{x^{6 m+1}}{6^{m} m!}$

## Section 7.3 Answers, pp. 234-238

7.3.1 (p. 234) $y=2-3 x-2 x^{2}+\frac{7}{2} x^{3}-\frac{55}{12} x^{4}+\frac{59}{8} x^{5}-\frac{83}{6} x^{6}+\frac{9547}{336} x^{7}+\cdots$
7.3.2 (p. 234) $y=-1+2 x-4 x^{3}+4 x^{4}+4 x^{5}-12 x^{6}+4 x^{7}+\cdots$
7.3.3 (p. 234) $y=1+x^{2}-\frac{2}{3} x^{3}+\frac{11}{6} x^{4}-\frac{9}{5} x^{5}+\frac{329}{90} x^{6}-\frac{1301}{315} x^{7}+\cdots$
7.3.4 (p. 234) $y=x-x^{2}-\frac{7}{2} x^{3}+\frac{15}{2} x^{4}+\frac{45}{8} x^{5}-\frac{261}{8} x^{6}+\frac{207}{16} x^{7}+\cdots$
7.3.5 (p. 234) $y=4+3 x-\frac{15}{4} x^{2}+\frac{1}{4} x^{3}+\frac{11}{16} x^{4}-\frac{5}{16} x^{5}+\frac{1}{20} x^{6}+\frac{1}{120} x^{7}+\cdots$
7.3.6 (p.234) $y=7+3 x-\frac{16}{3} x^{2}+\frac{13}{3} x^{3}-\frac{23}{9} x^{4}+\frac{10}{9} x^{5}-\frac{7}{27} x^{6}-\frac{1}{9} x^{7}+\cdots$
7.3.7 (p. 234) $y=2+5 x-\frac{7}{4} x^{2}-\frac{3}{16} x^{3}+\frac{37}{192} x^{4}-\frac{7}{192} x^{5}-\frac{1}{1920} x^{6}+\frac{19}{11520} x^{7}+\cdots$
$7.3 .8\left(\right.$ p. 234) $y=1-(x-1)+\frac{4}{3}(x-1)^{3}-\frac{4}{3}(x-1)^{4}-\frac{4}{5}(x-1)^{5}+\frac{136}{45}(x-1)^{6}-\frac{104}{63}(x-1)^{7}+\cdots$
7.3.9 (p. 235) $y=1-(x+1)+4(x+1)^{2}-\frac{13}{3}(x+1)^{3}+\frac{77}{6}(x+1)^{4}-\frac{278}{15}(x+1)^{5}+\frac{1942}{45}(x+1)^{6}-\frac{23332}{315}(x+1)^{7}+\cdots$
7.3.10 (p. 235) $y=2-(x-1)-\frac{1}{2}(x-1)^{2}+\frac{5}{3}(x-1)^{3}-\frac{19}{12}(x-1)^{4}+\frac{7}{30}(x-1)^{5}+\frac{59}{45}(x-1)^{6}-\frac{1091}{630}(x-1)^{7}+\cdots$
7.3.11 (p. 235) $y=-2+3(x+1)-\frac{1}{2}(x+1)^{2}-\frac{2}{3}(x+1)^{3}+\frac{5}{8}(x+1)^{4}-\frac{11}{30}(x+1)^{5}+\frac{29}{144}(x+1)^{6}-\frac{101}{840}(x+1)^{7}+\cdots$
7.3.12 (p. 235) $y=1-2(x-1)-3(x-1)^{2}+8(x-1)^{3}-4(x-1)^{4}-\frac{42}{5}(x-1)^{5}+19(x-1)^{6}-\frac{604}{35}(x-1)^{7}+\cdots$
7.3.19 (p. 236) $y=2-7 x-4 x^{2}-\frac{17}{6} x^{3}-\frac{3}{4} x^{4}-\frac{9}{40} x^{5}+\cdots$
7.3.20 (p. 236) $y=1-2(x-1)+\frac{1}{2}(x-1)^{2}-\frac{1}{6}(x-1)^{3}+\frac{5}{36}(x-1)^{4}-\frac{73}{1080}(x-1)^{5}+\cdots$
7.3.21 (p. 236) $y=2-(x+2)-\frac{7}{2}(x+2)^{2}+\frac{4}{3}(x+2)^{3}-\frac{1}{24}(x+2)^{4}+\frac{1}{60}(x+2)^{5}+\cdots$
7.3.22 (p. 236) $y=2-2(x+3)-(x+3)^{2}+(x+3)^{3}-\frac{11}{12}(x+3)^{4}+\frac{67}{60}(x+3)^{5}+\cdots$
7.3.23 (p. 236) $y=-1+2 x+\frac{1}{3} x^{3}-\frac{5}{12} x^{4}+\frac{2}{5} x^{5}+\cdots$
$7.3 .24\left(\right.$ p. 236) $y=2-3(x+1)+\frac{7}{2}(x+1)^{2}-5(x+1)^{3}+\frac{197}{24}(x+1)^{4}-\frac{287}{20}(x+1)^{5}+\cdots$
7.3 .25 (p. 236) $y=-2+3(x+2)-\frac{9}{2}(x+2)^{2}+\frac{11}{6}(x+2)^{3}+\frac{5}{24}(x+2)^{4}+\frac{7}{20}(x+2)^{5}+\cdots$
7.3 .26 (p. 236) $y=2-4(x-2)-\frac{1}{2}(x-2)^{2}+\frac{2}{9}(x-2)^{3}+\frac{49}{432}(x-2)^{4}+\frac{23}{1080}(x-2)^{5}+\cdots$
7.3 .27 (p. 236) $y=1+2(x+4)-\frac{1}{6}(x+4)^{2}-\frac{10}{27}(x+4)^{3}+\frac{19}{648}(x+4)^{4}+\frac{13}{324}(x+4)^{5}+\cdots$
$7.3 .28\left(\right.$ p. 236) $y=-1+2(x+1)-\frac{1}{4}(x+1)^{2}+\frac{1}{2}(x+1)^{3}-\frac{65}{96}(x+1)^{4}+\frac{67}{80}(x+1)^{5}+\cdots$
7.3.31 (p. 237) (a) $y=\frac{c_{1}}{1+x}+\frac{c_{2}}{1+2 x}$ (b) $y=\frac{c_{1}}{1-2 x}+\frac{c_{2}}{1-3 x}$ (c) $y=\frac{c_{1}}{1-2 x}+\frac{c_{2} x}{(1-2 x)^{2}}$
(d) $y=\frac{c_{1}}{2+x}+\frac{c_{2} x}{(2+x)^{2}}$ (e) $y=\frac{c_{1}}{2+x}+\frac{c_{2}}{2+3 x}$
7.3.32 (p. 238) $y=1-2 x-\frac{3}{2} x^{2}+\frac{5}{3} x^{3}+\frac{17}{24} x^{4}-\frac{11}{20} x^{5}+\cdots$
7.3.33 (p. 238) $y=1-2 x-\frac{5}{2} x^{2}+\frac{2}{3} x^{3}-\frac{3}{8} x^{4}+\frac{1}{3} x^{5}+\cdots$
7.3.34 (p. 238) $y=6-2 x+9 x^{2}+\frac{2}{3} x^{3}-\frac{23}{4} x^{4}-\frac{3}{10} x^{5}+\cdots$
7.3.35 (p. 238) $y=2-5 x+2 x^{2}-\frac{10}{3} x^{3}+\frac{3}{2} x^{4}-\frac{25}{12} x^{5}+\cdots$
7.3.36 (p. 238) $y=3+6 x-3 x^{2}+x^{3}-2 x^{4}-\frac{17}{20} x^{5}+\cdots$
7.3.37 (p. 238) $y=3-2 x-3 x^{2}+\frac{3}{2} x^{3}+\frac{3}{2} x^{4}-\frac{49}{80} x^{5}+\cdots$
7.3.38 (p. 238) $y=-2+3 x+\frac{4}{3} x^{2}-x^{3}-\frac{19}{54} x^{4}+\frac{13}{60} x^{5}+\cdots$
7.3.39 (p. 238) $y_{1}=\sum_{m=0}^{\infty} \frac{(-1)^{m} x^{2 m}}{m!}=e^{-x^{2}}, \quad y_{2}=\sum_{m=0}^{\infty} \frac{(-1)^{m} x^{2 m+1}}{m!}=x e^{-x^{2}}$
7.3.40 (p. 238) $y=-2+3 x+x^{2}-\frac{1}{6} x^{3}-\frac{3}{4} x^{4}+\frac{31}{120} x^{5}+\cdots$
7.3.41 (p. 238) $y=2+3 x-\frac{7}{2} x^{2}-\frac{5}{6} x^{3}+\frac{41}{24} x^{4}+\frac{41}{120} x^{5}+\cdots$
7.3.42 (p. 238) $y=-3+5 x-5 x^{2}+\frac{23}{6} x^{3}-\frac{23}{12} x^{4}+\frac{11}{30} x^{5}+\cdots$
7.3.43 (p. 238) $y=-2+3(x-1)+\frac{3}{2}(x-1)^{2}-\frac{17}{12}(x-1)^{3}-\frac{1}{12}(x-1)^{4}+\frac{1}{8}(x-1)^{5}+\cdots$
7.3 .44 (p. 238) $y=2-3(x+2)+\frac{1}{2}(x+2)^{2}-\frac{1}{3}(x+2)^{3}+\frac{31}{24}(x+2)^{4}-\frac{53}{120}(x+2)^{5}+\cdots$
7.3.45 (p. 238) $y=1-2 x+\frac{3}{2} x^{2}-\frac{11}{6} x^{3}+\frac{15}{8} x^{4}-\frac{71}{60} x^{5}+\cdots$
7.3.46 (p. 238) $y=2-(x+2)-\frac{7}{2}(x+2)^{2}-\frac{43}{6}(x+2)^{3}-\frac{203}{24}(x+2)^{4}-\frac{167}{30}(x+2)^{5}+\cdots$
7.3 .47 (p. 238) $y=2-x-x^{2}+\frac{7}{6} x^{3}-x^{4}+\frac{89}{120} x^{5}+\cdots$
7.3.48 (p. 238) $y=1+\frac{3}{2}(x-1)^{2}+\frac{1}{6}(x-1)^{3}-\frac{1}{8}(x-1)^{5}+\cdots$
7.3.49 (p. 238) $y=1-2(x-3)+\frac{1}{2}(x-3)^{2}-\frac{1}{6}(x-3)^{3}+\frac{1}{4}(x-3)^{4}-\frac{1}{6}(x-3)^{5}+\cdots$

Section 8.1 Answers, pp. 247-250
8.1.1 (p.247) (a) $\frac{1}{s^{2}}$ (b) $\frac{1}{(s+1)^{2}}$ (c) $\frac{b}{s^{2}-b^{2}}$ (d) $\frac{-2 s+5}{(s-1)(s-2)}$ (e) $\frac{2}{s^{3}}$
8.1.2 (p. 248) (a) $\frac{s^{2}+2}{\left[(s-1)^{2}+1\right]\left[(s+1)^{2}+1\right]}$ (b) $\frac{2}{s\left(s^{2}+4\right)}$ (c) $\frac{s^{2}+8}{s\left(s^{2}+16\right)}$ (d) $\frac{s^{2}-2}{s\left(s^{2}-4\right)}$
(e) $\frac{4 s}{\left(s^{2}-4\right)^{2}}$ (f) $\frac{1}{s^{2}+4}$ (g) $\frac{1}{\sqrt{2}} \frac{s+1}{s^{2}+1}$ (h) $\frac{5 s}{\left(s^{2}+4\right)\left(s^{2}+9\right)}$ (i) $\frac{s^{3}+2 s^{2}+4 s+32}{\left(s^{2}+4\right)\left(s^{2}+16\right)}$
8.1.4 (p. 248) (a) $f(3-)=-1, f(3)=f(3+)=1$ (b) $f(1-)=3, f(1)=4, f(1+)=1$
(c) $f\left(\frac{\pi}{2}-\right)=1, f\left(\frac{\pi}{2}\right)=f\left(\frac{\pi}{2}+\right)=2, f(\pi-)=0, f(\pi)=f(\pi+)=-1$
(d) $f(1-)=1, f(1)=2, f(1+)=1, f(2-)=0, f(2)=3, f(2+)=6$
8.1.5 (p. 248) (a) $\frac{1-e^{-(s+1)}}{s+1}+\frac{e^{-(s+2)}}{s+2}$ (b) $\frac{1}{s}+e^{-4 s}\left(\frac{1}{s^{2}}+\frac{3}{s}\right)$ (c) $\frac{1-e^{-s}}{s^{2}}$ (d) $\frac{1-e^{-(s-1)}}{(s-1)^{2}}$
8.1.7 (p. 248) $\mathcal{L}\left(e^{\lambda t} \cos \omega t\right)=\frac{(s-\lambda)^{2}-\omega^{2}}{\left((s-\lambda)^{2}+\omega^{2}\right)^{2}} \mathcal{L}\left(e^{\lambda t} \sin \omega t\right)=\frac{2 \omega(s-\lambda)}{\left((s-\lambda)^{2}+\omega^{2}\right)^{2}}$
8.1.15 (p. 249) (a) $\tan ^{-1} \frac{\omega}{s}, \quad s>0$ (b) $\frac{1}{2} \ln \frac{s^{2}}{s^{2}+\omega^{2}}, \quad s>0$ (c) $\ln \frac{s-b}{s-a}, \quad s>\max (a, b)$

$$
\text { (d) } \frac{1}{2} \ln \frac{s^{2}}{s^{2}-1}, \quad s>1 \text { (e) } \frac{1}{4} \ln \frac{s^{2}}{s^{2}-4}, \quad s>2
$$

8.1.18 (p. 250) (a) $\frac{1}{s^{2}} \tanh \frac{s}{2}$ (b) $\frac{1}{s} \tanh \frac{s}{4}$ (c) $\frac{1}{s^{2}+1} \operatorname{coth} \frac{\pi s}{2}$ (d) $\frac{1}{\left(s^{2}+1\right)\left(1-e^{-\pi s}\right)}$

Section 8.2 Answers, pp. 256-257
8.2.1 (p. 256) (a) $\frac{t^{3} e^{7 t}}{2}$ (b) $2 e^{2 t} \cos 3 t$ (c) $\frac{e^{-2 t}}{4} \sin 4 t$ (d) $\frac{2}{3} \sin 3 t$ (e) $t \cos t$ (f) $\frac{e^{2 t}}{2} \sinh 2 t$ (g) $\frac{2 t e^{2 t}}{3} \sin 9 t$ (h) $\frac{2 e^{3 t}}{3} \sinh 3 t$ (i) $e^{2 t} t \cos t$
8.2 .2 (p.256) (a) $t^{2} e^{7 t}+\frac{17}{6} t^{3} e^{7 t}$ (b) $e^{2 t}\left(\frac{1}{6} t^{3}+\frac{1}{6} t^{4}+\frac{1}{40} t^{5}\right)$ (c) $e^{-3 t}\left(\cos 3 t+\frac{2}{3} \sin 3 t\right)$
(d) $2 \cos 3 t+\frac{1}{3} \sin 3 t$ (e) $(1-t) e^{-t}$ (f) $\cosh 3 t+\frac{1}{3} \sinh 3 t$ (g) $\left(1-t-t^{2}-\frac{1}{6} t^{3}\right) e^{-t}$
(h) $e^{t}\left(2 \cos 2 t+\frac{5}{2} \sin 2 t\right)$ (i) $1-\cos t$ (j) $3 \cosh t+4 \sinh t$ (k) $3 e^{t}+4 \cos 3 t+\frac{1}{3} \sin 3 t$
(l) $3 t e^{-2 t}-2 \cos 2 t-3 \sin 2 t$
8.2 .3 (p. 256) (a) $\frac{1}{4} e^{2 t}-\frac{1}{4} e^{-2 t}-e^{-t}$ (b) $\frac{1}{5} e^{-4 t}-\frac{41}{5} e^{t}+5 e^{3 t}$ (c) $-\frac{1}{2} e^{2 t}-\frac{13}{10} e^{-2 t}-\frac{1}{5} e^{3 t}$
(d) $-\frac{2}{5} e^{-4 t}-\frac{3}{5} e^{t}$ (e) $\frac{3}{20} e^{2 t}-\frac{37}{12} e^{-2 t}+\frac{1}{3} e^{t}+\frac{8}{5} e^{-3 t}$ (f) $\frac{39}{10} e^{t}+\frac{3}{14} e^{3 t}+\frac{23}{105} e^{-4 t}-\frac{7}{3} e^{2 t}$
8.2.4 (p. 256) (a) $\frac{4}{5} e^{-2 t}-\frac{1}{2} e^{-t}-\frac{3}{10} \cos t+\frac{11}{10} \sin t$ (b) $\frac{2}{5} \sin t+\frac{6}{5} \cos t+\frac{7}{5} e^{-t} \sin t-\frac{6}{5} e^{-t} \cos t$
(c) $\frac{8}{13} e^{2 t}-\frac{8}{13} e^{-t} \cos 2 t+\frac{15}{26} e^{-t} \sin 2 t$ (d) $\frac{1}{2} t e^{t}+\frac{3}{8} e^{t}+e^{-2 t}-\frac{11}{8} e^{-3 t}$
(e) $\frac{2}{3} t e^{t}+\frac{1}{9} e^{t}+t e^{-2 t}-\frac{1}{9} e^{-2 t}$ (f) $-e^{t}+\frac{5}{2} t e^{t}+\cos t-\frac{3}{2} \sin t$
8.2 .5 (p. 256) (a) $\frac{3}{5} \cos 2 t+\frac{1}{5} \sin 2 t-\frac{3}{5} \cos 3 t-\frac{2}{15} \sin 3 t$ (b) $-\frac{4}{15} \cos t+\frac{1}{15} \sin t+\frac{4}{15} \cos 4 t-\frac{1}{60} \sin 4 t$
(c) $\frac{5}{3} \cos t+\sin t-\frac{5}{3} \cos 2 t-\frac{1}{2} \sin 2 t$ (d) $-\frac{1}{3} \cos \frac{t}{2}+\frac{2}{3} \sin \frac{t}{2}+\frac{1}{3} \cos t-\frac{1}{3} \sin t$
(e) $\frac{1}{15} \cos \frac{t}{4}-\frac{8}{15} \sin \frac{t}{4}-\frac{1}{15} \cos 4 t+\frac{1}{30} \sin 4 t$ (f) $\frac{2}{5} \cos \frac{t}{3}-\frac{3}{5} \sin \frac{t}{3}-\frac{2}{5} \cos \frac{t}{2}+\frac{2}{5} \sin \frac{t}{2}$
8.2.6 (p.256) (a) $e^{t}(\cos 2 t+\sin 2 t)-e^{-t}\left(\cos 3 t+\frac{4}{3} \sin 3 t\right)$ (b) $e^{3 t}\left(-\cos 2 t+\frac{3}{2} \sin 2 t\right)+e^{-t}\left(\cos 2 t+\frac{1}{2} \sin 2 t\right)$
(c) $e^{-2 t}\left(\frac{1}{8} \cos t+\frac{1}{4} \sin t\right)-e^{2 t}\left(\frac{1}{8} \cos 3 t-\frac{1}{12} \sin 3 t\right)$ (d) $e^{2 t}\left(\cos t+\frac{1}{2} \sin t\right)-e^{3 t}\left(\cos 2 t-\frac{1}{4} \sin 2 t\right)$
(e) $e^{t}\left(\frac{1}{5} \cos t+\frac{2}{5} \sin t\right)-e^{-t}\left(\frac{1}{5} \cos 2 t+\frac{2}{5} \sin 2 t\right)$ (f) $e^{t / 2}\left(-\cos t+\frac{9}{8} \sin t\right)+e^{-t / 2}\left(\cos t-\frac{1}{8} \sin t\right)$
8.2 .7 (p.257) (a) $1-\cos t$ (b) $\frac{e^{t}}{16}(1-\cos 4 t)$ (c) $\frac{4}{9} e^{2 t}+\frac{5}{9} e^{-t} \sin 3 t-\frac{4}{9} e^{-t} \cos 3 t$ (d) $3 e^{t / 2}-\frac{7}{2} e^{t} \sin 2 t-3 e^{t} \cos 2 t$
(e) $\frac{1}{4} e^{3 t}-\frac{1}{4} e^{-t} \cos 2 t$ (f) $\frac{1}{9} e^{2 t}-\frac{1}{9} e^{-t} \cos 3 t+\frac{5}{9} e^{-t} \sin 3 t$
8.2 .8 (p. 257) (a) $-\frac{3}{10} \sin t+\frac{2}{5} \cos t-\frac{3}{4} e^{t}+\frac{7}{20} e^{3 t}$ (b) $-\frac{3}{5} e^{-t} \sin t+\frac{1}{5} e^{-t} \cos t-\frac{1}{2} e^{-t}+\frac{3}{10} e^{t}$
(c) $-\frac{1}{10} e^{t} \sin t-\frac{7}{10} e^{t} \cos t+\frac{1}{5} e^{-t}+\frac{1}{2} e^{2 t}$ (d) $-\frac{1}{2} e^{t}+\frac{7}{10} e^{-t}-\frac{1}{5} \cos 2 t+\frac{3}{5} \sin 2 t$
(e) $\frac{3}{10}+\frac{1}{10} e^{2 t}+\frac{1}{10} e^{t} \sin 2 t-\frac{2}{5} e^{t} \cos 2 t$ (f) $-\frac{4}{9} e^{2 t} \cos 3 t+\frac{1}{3} e^{2 t} \sin 3 t-\frac{5}{9} e^{2 t}+e^{t}$
8.2 .9 (p. 257) $\frac{1}{a} e^{\frac{b}{a} t} f\left(\frac{t}{a}\right)$

## Section 8.3 Answers, pp. 261-262

8.3 .1 (p. 261) $y=\frac{1}{6} e^{t}-\frac{9}{2} e^{-t}+\frac{16}{3} e^{-2 t} 8.3 .2$ (p. 261) $y=-\frac{1}{3}+\frac{8}{15} e^{3 t}+\frac{4}{5} e^{-2 t}$
8.3 .3 (p. 261) $y=-\frac{23}{15} e^{-2 t}+\frac{1}{3} e^{t}+\frac{1}{5} e^{3 t} 8.3 .4$ (p. 261) $y=-\frac{1}{4} e^{2 t}+\frac{17}{20} e^{-2 t}+\frac{2}{5} e^{3 t}$
8.3 .5 (p.261) $y=\frac{11}{15} e^{-2 t}+\frac{1}{6} e^{t}+\frac{1}{10} e^{3 t} 8.3 .6$ (p. 261) $y=e^{t}+2 e^{-2 t}-2 e^{-t}$
8.3 .7 (p. 261) $y=\frac{5}{3} \sin t-\frac{1}{3} \sin 2 t 8.3 .8$ (p. 261) $y=4 e^{t}-4 e^{2 t}+e^{3 t}$
8.3 .9 (p. 261) $y=-\frac{7}{2} e^{2 t}+\frac{13}{3} e^{t}+\frac{1}{6} e^{4 t} 8.3 .10$ (p.261) $y=\frac{5}{2} e^{t}-4 e^{2 t}+\frac{1}{2} e^{3 t}$
8.3 .11 (p. 262) $y=\frac{1}{3} e^{t}-2 e^{-t}+\frac{5}{3} e^{-2 t} 8.3 .12$ (p. 262) $y=2-e^{-2 t}+e^{t}$
8.3 .13 (p. 262) $y=1-\cos 2 t+\frac{1}{2} \sin 2 t 8.3 .14$ (p. 262) $y=-\frac{1}{3}+\frac{8}{15} e^{3 t}+\frac{4}{5} e^{-2 t}$
8.3 .15 (p. 262) $y=\frac{1}{6} e^{t}-\frac{2}{3} e^{-2 t}+\frac{1}{2} e^{-t} 8.3 .16$ (p. 262) $y=-1+e^{t}+e^{-t}$
8.3.17 (p. 262) $y=\cos 2 t-\sin 2 t+\sin t 8.3 .18$ (p. 262) $y=\frac{7}{3}-\frac{7}{2} e^{-t}+\frac{1}{6} e^{3 t}$
8.3.19 (p. 262) $y=1+\cos t 8.3 .20$ (p. 262) $y=t+\sin t 8.3 .21$ (p. 262) $y=t-6 \sin t+\cos t+\sin 2 t$
8.3.22 (p. 262) $y=e^{-t}+4 e^{-2 t}-4 e^{-3 t} 8.3 .23$ (p. 262) $y=-3 \cos t-2 \sin t+e^{-t}(2+5 t)$
8.3.24 (p. 262) $y=-\sin t-2 \cos t+3 e^{3 t}+e^{-t} 8.3 .25$ (p. 262) $y=(3 t+4) \sin t-(2 t+6) \cos t$
8.3.26 (p. 262) $y=-(2 t+2) \cos 2 t+\sin 2 t+3 \cos t 8.3 .27$ (p. 262) $y=e^{t}(\cos t-3 \sin t)+e^{3 t}$
8.3.28 (p. 262) $y=-1+t+e^{-t}(3 \cos t-5 \sin t) 8.3 .29$ (p. 262) $y=4 \cos t-3 \sin t-e^{t}(3 \cos t-8 \sin t)$
8.3.30 (p. 262) $y=e^{-t}-2 e^{t}+e^{-2 t}(\cos 3 t-11 / 3 \sin 3 t)$
8.3.31 (p. 262) $y=e^{-t}(\sin t-\cos t)+e^{-2 t}(\cos t+4 \sin t)$
8.3 .32 (p. 262) $y=\frac{1}{5} e^{2 t}-\frac{4}{3} e^{t}+\frac{32}{15} e^{-t / 2} 8.3 .33$ (p. 262) $y=\frac{1}{7} e^{2 t}-\frac{2}{5} e^{t / 2}+\frac{9}{35} e^{-t / 3}$
8.3.34 (p. 262) $y=e^{-t / 2}(5 \cos (t / 2)-\sin (t / 2))+2 t-4$
8.3.35 (p. 262) $y=\frac{1}{17}\left(12 \cos t+20 \sin t-3 e^{t / 2}(4 \cos t+\sin t)\right)$.
8.3.36 (p. 262) $y=\frac{e^{-t / 2}}{10}(5 t+26)-\frac{1}{5}(3 \cos t+\sin t) 8.3 .37($ p. 262 $) y=\frac{1}{100}\left(3 e^{3 t}-e^{t / 3}(3+310 t)\right)$

Section 8.4 Answers, pp. 269-271
8.4.1 (p. 269) $1+u(t-4)(t-1) ; \frac{1}{s}+e^{-4 s}\left(\frac{1}{s^{2}}+\frac{3}{s}\right) 8.4 .2\left(\mathbf{p . 2 6 9 )} t+u(t-1)(1-t) ; \frac{1-e^{-s}}{s^{2}}\right.$
8.4 .3 (p. 269) $2 t-1-u(t-2)(t-1) ;\left(\frac{2}{s^{2}}-\frac{1}{s}\right)-e^{-2 s}\left(\frac{1}{s^{2}}+\frac{1}{s}\right)$
8.4 .4 (p. 269) $1+u(t-1)(t+1) ; \frac{1}{s}+e^{-s}\left(\frac{1}{s^{2}}+\frac{2}{s}\right)$
8.4 .5 (p. 269) $t-1+u(t-2)(5-t) ; \frac{1}{s^{2}}-\frac{1}{s}-e^{-2 s}\left(\frac{1}{s^{2}}-\frac{3}{s}\right)$
8.4 .6 (p. 269) $t^{2}(1-u(t-1)) ; \frac{2}{s^{3}}-e^{-s}\left(\frac{2}{s^{3}}+\frac{2}{s^{2}}+\frac{1}{s}\right)$
8.4.7 (p. 269) $u(t-2)\left(t^{2}+3 t\right) ; e^{-2 s}\left(\frac{2}{s^{3}}+\frac{7}{s^{2}}+\frac{10}{s}\right)$
8.4.8 (p. 269) $t^{2}+2+u(t-1)\left(t-t^{2}-2\right) ; \frac{2}{s^{3}}+\frac{2}{s}-e^{-s}\left(\frac{2}{s^{3}}+\frac{1}{s^{2}}+\frac{2}{s}\right)$
$8.4 .9(\mathbf{p} .269) t e^{t}+u(t-1)\left(e^{t}-t e^{t}\right) ; \frac{1-e^{-(s-1)}}{(s-1)^{2}}$
8.4 .10 (p. 269) $e^{-t}+u(t-1)\left(e^{-2 t}-e^{-t}\right) ; \frac{1-e^{-(s+1)}}{s+1}+\frac{e^{-(s+2)}}{s+2}$
8.4.11 (p. 269) $-t+2 u(t-2)(t-2)-u(t-3)(t-5) ;-\frac{1}{s^{2}}+\frac{2 e^{-2 s}}{s^{2}}+e^{-3 s}\left(\frac{2}{s}-\frac{1}{s^{2}}\right)$
8.4 .12 (p. 269) $[u(t-1)-u(t-2)] t ; e^{-s}\left(\frac{1}{s^{2}}+\frac{1}{s}\right)-e^{-2 s}\left(\frac{1}{s^{2}}+\frac{2}{s}\right)$
$8.4 .13\left(\mathbf{p . 2 7 0 )} t+u(t-1)\left(t^{2}-t\right)-u(t-2) t^{2} ; \frac{1}{s^{2}}+e^{-s}\left(\frac{2}{s^{3}}+\frac{1}{s^{2}}\right)-e^{-2 s}\left(\frac{2}{s^{3}}+\frac{4}{s^{2}}+\frac{4}{s}\right)\right.$
$8.4 .14(\mathbf{p . 2 7 0}) t+u(t-1)(2-2 t)+u(t-2)(4+t) ; \frac{1}{s^{2}}-2 \frac{e^{-s}}{s^{2}}+e^{-2 s}\left(\frac{1}{s^{2}}+\frac{6}{s}\right)$
8.4 .15 (p.270) $\sin t+u(t-\pi / 2) \sin t+u(t-\pi)(\cos t-2 \sin t) ; \frac{1+e^{-\frac{\pi}{2} s} s-e^{-\pi s}(s-2)}{s^{2}+1}$
8.4.16 (p. 270) $2-2 u(t-1) t+u(t-3)(5 t-2) ; \frac{2}{s}-e^{-s}\left(\frac{2}{s^{2}}+\frac{2}{s}\right)+e^{-3 s}\left(\frac{5}{s^{2}}+\frac{13}{s}\right)$
8.4.17 (p. 270) $3+u(t-2)(3 t-1)+u(t-4)(t-2) ; \frac{3}{s}+e^{-2 s}\left(\frac{3}{s^{2}}+\frac{5}{s}\right)+e^{-4 s}\left(\frac{1}{s^{2}}+\frac{2}{s}\right)$
8.4.18(p.270) $(t+1)^{2}+u(t-1)(2 t+3) ; \frac{2}{s^{3}}+\frac{2}{s^{2}}+\frac{1}{s}+e^{-s}\left(\frac{2}{s^{2}}+\frac{5}{s}\right)$
8.4 .19 (p. 270) $u(t-2) e^{2(t-2)}=\left\{\begin{array}{cl}0, & 0 \leq t<2, \\ e^{2(t-2)}, & t \geq 2 .\end{array}\right.$
$8.4 .20(\mathbf{p . 2 7 0}) u(t-1)\left(1-e^{-(t-1)}\right)=\left\{\begin{array}{cl}0, & 0 \leq t<1, \\ 1-e^{-(t-1)}, & t \geq 1 .\end{array}\right.$
8.4.21 (p. 270) $u(t-1) \frac{(t-1)^{2}}{2}+u(t-2)(t-2)=\left\{\begin{array}{cl}0, & 0 \leq t<1, \\ \frac{(t-1)^{2}}{2}, & 1 \leq t<2, \\ \frac{t^{2}-3}{2}, & t \geq 2 .\end{array}\right.$
8.4.22 (p. 270) $2+t+u(t-1)(4-t)+u(t-3)(t-2)=\left\{\begin{array}{cl}2+t, & 0 \leq t<1, \\ 6, & 1 \leq t<3, \\ t+4, & t \geq 3 .\end{array}\right.$
8.4 .23 (p. 270) $5-t+u(t-3)(7 t-15)+\frac{3}{2} u(t-6)(t-6)^{2}=\left\{\begin{array}{cc}5-t, & 0 \leq t<3, \\ 6 t-10, & 3 \leq t<6, \\ 44-12 t+\frac{3}{2} t^{2}, & t \geq 6 .\end{array}\right.$
8.4 .24 (p. 270) $u(t-\pi) e^{-2(t-\pi)}(2 \cos t-5 \sin t)=\left\{\begin{array}{cl}0, & 0 \leq t<\pi, \\ e^{-2(t-\pi)}(2 \cos t-5 \sin t), & t \geq \pi .\end{array}\right.$
8.4.25 (p. 270) $1-\cos t+u(t-\pi / 2)(3 \sin t+\cos t)=\left\{\begin{array}{cl}1-\cos t, & 0 \leq t<\frac{\pi}{2}, \\ 1+3 \sin t, & t \geq \frac{\pi}{2} .\end{array}\right.$
8.4 .26 (p. 270) $u(t-2)\left(4 e^{-(t-2)}-4 e^{2(t-2)}+2 e^{(t-2)}\right)=\left\{\begin{array}{cl}0, & 0 \leq t<2, \\ 4 e^{-(t-2)}-4 e^{2(t-2)}+2 e^{(t-2)}, & t \geq 2 .\end{array}\right.$
8.4 .27 (p. 270) $1+t+u(t-1)(2 t+1)+u(t-3)(3 t-5)=\left\{\begin{array}{cl}t+1, & 0 \leq t<1, \\ 3 t+2, & 1 \leq t<3, \\ 6 t-3, & t \geq 3 .\end{array}\right.$
8.4.28 (p. 270) $1-t^{2}+u(t-2)\left(-\frac{t^{2}}{2}+2 t+1\right)+u(t-4)(t-4)=\left\{\begin{array}{cl}1-t^{2}, & 0 \leq t<2 \\ -\frac{3 t^{2}}{2}+2 t+2, & 2 \leq t<4, \\ -\frac{3 t^{2}}{2}+3 t-2, & t \geq 4 .\end{array}\right.$
8.4.29 (p. 270) $\frac{e^{-\tau s}}{s}$ 8.4.30 (p. 270) For each $t$ only finitely many terms are nonzero.
8.4 .33 (p. 271) $1+\sum_{m=1}^{\infty} u(t-m) ; \frac{1}{s\left(1-e^{-s}\right)} 8.4 .34$ (p. 271) $1+2 \sum_{m=1}^{\infty}(-1)^{m} u(t-m) ; \frac{1}{s} ; \frac{1-e^{-s}}{1+e^{-s}}$
8.4.35(p.271) $1+\sum_{m=1}^{\infty}(2 m+1) u(t-m) ; \frac{e^{-s}\left(1+e^{-s}\right)}{s\left(1-e^{-s}\right)^{2}} 8.4 .36\left(\mathbf{p . 2 7 1 )} \sum_{m=1}^{\infty}(-1)^{m}(2 m-1) u(t-m) ; \frac{1}{s} \frac{\left(1-e^{s}\right)}{\left(1+e^{s}\right)^{2}}\right.$

Section 8.5 Answers, pp. 276-279
8.5.1 (p.276) $y=3(1-\cos t)-3 u(t-\pi)(1+\cos t)$
$8.5 .2(\mathbf{p} .277) y=3-2 \cos t+2 u(t-4)(t-4-\sin (t-4)) 8.5 .3(\mathbf{p . 2 7 7}) y=-\frac{15}{2}+\frac{3}{2} e^{2 t}-2 t+\frac{u(t-1)}{2}\left(e^{2(t-1)}-2 t+1\right)$
8.5.4 (p. 277) $y=\frac{1}{2} e^{t}+\frac{13}{6} e^{-t}+\frac{1}{3} e^{2 t}+u(t-2)\left(-1+\frac{1}{2} e^{t-2}+\frac{1}{2} e^{-(t-2)}+\frac{1}{2} e^{t+2}-\frac{1}{6} e^{-(t-6)}-\frac{1}{3} e^{2 t}\right)$
8.5.5 (p. 277) $y=-7 e^{t}+4 e^{2 t}+u(t-1)\left(\frac{1}{2}-e^{t-1}+\frac{1}{2} e^{2(t-1)}\right)-2 u(t-2)\left(\frac{1}{2}-e^{t-2}+\frac{1}{2} e^{2(t-2)}\right)$
8.5.6 (p.277) $y=\frac{1}{3} \sin 2 t-3 \cos 2 t+\frac{1}{3} \sin t-2 u(t-\pi)\left(\frac{1}{3} \sin t+\frac{1}{6} \sin 2 t\right)+u(t-2 \pi)\left(\frac{1}{3} \sin t-\frac{1}{6} \sin 2 t\right)$
8.5.7 (p.277) $y=\frac{1}{4}-\frac{31}{12} e^{4 t}+\frac{16}{3} e^{t}+u(t-1)\left(\frac{2}{3} e^{t-1}-\frac{1}{6} e^{4(t-1)}-\frac{1}{2}\right)+u(t-2)\left(\frac{1}{4}+\frac{1}{12} e^{4(t-2)}-\frac{1}{3} e^{t-2}\right)$
8.5.8 (p. 277) $y=\frac{1}{8}(\cos t-\cos 3 t)-\frac{1}{8} u\left(t-\frac{3 \pi}{2}\right)\left(\sin t-\cos t+\sin 3 t-\frac{1}{3} \cos 3 t\right)$
8.5.9 (p. 277) $y=\frac{t}{4}-\frac{1}{8} \sin 2 t+\frac{1}{8} u\left(t-\frac{\pi}{2}\right)(\pi \cos 2 t-\sin 2 t+2 \pi-2 t)$
8.5.10 (p. 277) $y=t-\sin t-2 u(t-\pi)(t+\sin t+\pi \cos t)$
8.5.11 (p. 277) $y=u(t-2)\left(t-\frac{1}{2}+\frac{e^{2(t-2)}}{2}-2 e^{t-2}\right)$
8.5.12 (p. 277) $y=t+\sin t+\cos t-u(t-2 \pi)(3 t-3 \sin t-6 \pi \cos t)$
8.5.13 (p. 278) $y=\frac{1}{2}+\frac{1}{2} e^{-2 t}-e^{-t}+u(t-2)\left(2 e^{-(t-2)}-e^{-2(t-2)}-1\right)$
8.5 .14 (p. 278) $y=-\frac{1}{3}-\frac{1}{6} e^{3 t}+\frac{1}{2} e^{t}+u(t-1)\left(\frac{2}{3}+\frac{1}{3} e^{3(t-1)}-e^{t-1}\right)$
8.5.15 (p. 278) $y=\frac{1}{4}\left(e^{t}+e^{-t}(11+6 t)\right)+u(t-1)\left(t e^{-(t-1)}-1\right)$
8.5.16 (p. 278) $y=e^{t}-e^{-t}-2 t e^{-t}-u(t-1)\left(e^{t}-e^{-(t-2)}-2(t-1) e^{-(t-2)}\right)$
8.5.17 (p. 278) $y=t e^{-t}+e^{-2 t}+u(t-1)\left(e^{-t}(2-t)-e^{-(2 t-1)}\right)$
8.5.18 (p. 278) $y=y=\frac{t^{2} e^{2 t}}{2}-t e^{2 t}-u(t-2)(t-2)^{2} e^{2 t}$
8.5.19 (p. 278) $y=\frac{t^{4}}{12}+1-\frac{1}{12} u(t-1)\left(t^{4}+2 t^{3}-10 t+7\right)+\frac{1}{6} u(t-2)\left(2 t^{3}+3 t^{2}-36 t+44\right)$
8.5.20 (p. 278) $y=\frac{1}{2} e^{-t}(3 \cos t+\sin t)+\frac{1}{2}$

$$
\begin{aligned}
& -u(t-2 \pi)\left(e^{-(t-2 \pi)}\left((\pi-1) \cos t+\frac{2 \pi-1}{2} \sin t\right)+1-\frac{t}{2}\right) \\
& -\frac{1}{2} u(t-3 \pi)\left(e^{-(t-3 \pi)}(3 \pi \cos t+(3 \pi+1) \sin t)+t\right)
\end{aligned}
$$

8.5 .21 (p. 278) $y=\frac{t^{2}}{2}+\sum_{m=1}^{\infty} u(t-m) \frac{(t-m)^{2}}{2}$
8.5.22 (p. 278) (a) $y=\left\{\begin{array}{cl}2 m+1-\cos t, & 2 m \pi \leq t<(2 m+1) \pi \quad(m=0,1, \ldots) \\ 2 m, & (2 m-1) \pi \leq t<2 m \pi\end{array} \quad(m=1,2, \ldots)\right.$
(b) $y=(m+1)(t-\sin t-m \pi \cos t), 2 m \pi \leq t<(2 m+2) \pi \quad(m=0,1, \ldots)$
(c) $y=(-1)^{m}-(2 m+1) \cos t, m \pi \leq t<(m+1) \pi \quad(m=0,1, \ldots)$
(d) $y=\frac{e^{m+1}-1}{2(e-1)}\left(e^{t-m}+e^{-t}\right)-m-1, \quad m \leq t<m+1(m=0,1 \ldots)$
(e) $y=\left(m+1-\left(\frac{e^{2(m+1) \pi}-1}{e^{2 \pi}-1}\right) e^{-t}\right) \sin t 2 m \pi \leq t<2(m+1) \pi \quad(m=0,1, \ldots)$
(f) $y=\frac{m+1}{2}-e^{t-m} \frac{e^{m+1}-1}{e-1}+\frac{1}{2} e^{2(t-m)} \frac{e^{2 m+2}-1}{e^{2}-1}, m \leq t<m+1 \quad(m=0,1, \ldots)$

## Section 8.6 Answers, pp. 287-290

8.6.1 (p. 287) (a) $\frac{1}{2} \int_{0}^{t} \tau \sin 2(t-\tau) d \tau$ (b) $\int_{0}^{t} e^{-2 \tau} \cos 3(t-\tau) d \tau$
(c) $\frac{1}{2} \int_{0}^{t} \sin 2 \tau \cos 3(t-\tau) d \tau$ or $\frac{1}{3} \int_{0}^{t} \sin 3 \tau \cos 2(t-\tau) d \tau$ (d) $\int_{0}^{t} \cos \tau \sin (t-\tau) d \tau$
(e) $\int_{0}^{t} e^{a \tau} d \tau(\mathbf{f}) e^{-t} \int_{0}^{t} \sin (t-\tau) d \tau(\mathbf{g}) e^{-2 t} \int_{0}^{t} \tau e^{\tau} \sin (t-\tau) d \tau$
(h) $\frac{e^{-2 t}}{2} \int_{0}^{t} \tau^{2}(t-\tau) e^{3 \tau} d \tau(\mathbf{i}) \int_{0}^{t}(t-\tau) e^{\tau} \cos \tau d \tau(\mathbf{j}) \int_{0}^{t} e^{-3 \tau} \cos \tau \cos 2(t-\tau) d \tau$
(k) $\frac{1}{4!5!} \int_{0}^{t} \tau^{4}(t-\tau)^{5} e^{3 \tau} d \tau$ (l) $\frac{1}{4} \int_{0}^{t} \tau^{2} e^{\tau} \sin 2(t-\tau) d \tau$
(m) $\frac{1}{2} \int_{0}^{t} \tau(t-\tau)^{2} e^{2(t-\tau)} d \tau$ (n) $\frac{1}{5!6!} \int_{0}^{t}(t-\tau)^{5} e^{2(t-\tau)} \tau^{6} d \tau$
8.6 .2 (p. 287)
(a) $\frac{a s}{\left(s^{2}+a^{2}\right)\left(s^{2}+b^{2}\right)}$
(b) $\frac{a}{(s-1)\left(s^{2}+a^{2}\right)}$ (c) $\frac{a s}{\left(s^{2}-a^{2}\right)^{2}}$
(d) $\frac{2 \omega s\left(s^{2}-\omega^{2}\right)}{\left(s^{2}+\omega^{2}\right)^{4}}$
(e) $\frac{(s-1) \omega}{\left((s-1)^{2}+\omega^{2}\right)^{2}}$ (f) $\frac{2}{(s-2)^{3}(s-1)^{2}}$ (g) $\frac{s+1}{(s+2)^{2}\left[(s+1)^{2}+\omega^{2}\right]}$
(h) $\frac{1}{(s-3)\left((s-1)^{2}-1\right)}$ (i) $\frac{2}{(s-2)^{2}\left(s^{2}+4\right)}$ (j) $\frac{6}{s^{4}(s-1)}$ (k) $\frac{3 \cdot 6!}{s^{7}\left[(s+1)^{2}+9\right]}$
(l) $\frac{12}{s^{7}}$ (m) $\frac{2 \cdot 7!}{s^{8}\left[(s+1)^{2}+4\right]}$ (n) $\frac{48}{s^{5}\left(s^{2}+4\right)}$
8.6 .3 (p.287) (a) $y=\frac{2}{\sqrt{5}} \int_{0}^{t} f(t-\tau) e^{-3 \tau / 2} \sinh \frac{\sqrt{5} \tau}{2} d \tau$ (b) $y=\frac{1}{2} \int_{0}^{t} f(t-\tau) \sin 2 \tau d \tau$
(c) $y=\int_{0}^{t} \tau e^{-\tau} f(t-\tau) d \tau$ (d) $y(t)=-\frac{1}{k} \sin k t+\cos k t+\frac{1}{k} \int_{0}^{t} f(t-\tau) \sin k \tau d \tau$
(e) $y=-2 t e^{-3 t}+\int_{0}^{t} \tau e^{-3 \tau} f(t-\tau) d \tau$ (f) $y=\frac{3}{2} \sinh 2 t+\frac{1}{2} \int_{0}^{t} f(t-\tau) \sinh 2 \tau d \tau$
(g) $y=e^{3 t}+\int_{0}^{t}\left(e^{3 \tau}-e^{2 \tau}\right) f(t-\tau) d \tau$ (h) $y=\frac{k_{1}}{\omega} \sin \omega t+k_{0} \cos \omega t+\frac{1}{\omega} \int_{0}^{t} f(t-\tau) \sin \omega \tau d \tau$
8.6 .4 (p. 288) (a) $y=\sin t$ (b) $y=t e^{-t}$ (c) $y=1+2 t e^{t}$ (d) $y=t+\frac{t^{2}}{2}$
(e) $y=4+\frac{5}{2} t^{2}+\frac{1}{24} t^{4}$ (f) $y=1-t$
8.6 .5 (p. 288)
$\begin{array}{lll}\text { (a) } \frac{7!8!}{16!} t^{16} & \text { (b) } \frac{13!7!}{21!} t^{21}\end{array}$
(c) $\frac{6!7!}{14!} t^{14}$ (d) $\frac{1}{2}\left(e^{-t}+\sin t-\cos t\right)$
(e) $\frac{1}{3}(\cos t-\cos 2 t)$

## Section 8.7 Answers, pp. 296-297

8.7 .1 (p. 296) $y=\frac{1}{2} e^{2 t}-4 e^{-t}+\frac{11}{2} e^{-2 t}+2 u(t-1)\left(e^{-(t-1)}-e^{-2(t-1)}\right)$
8.7.2 (p. 296) $y=2 e^{-2 t}+5 e^{-t}+\frac{5}{3} u(t-1)\left(e^{(t-1)}-e^{-2(t-1)}\right)$
8.7 .3 (p. 296) $y=\frac{1}{6} e^{2 t}-\frac{2}{3} e^{-t}-\frac{1}{2} e^{-2 t}+\frac{5}{2} u(t-1) \sinh 2(t-1)$
8.7.4 (p. 296) $y=\frac{1}{8}(8 \cos t-5 \sin t-\sin 3 t)-2 u(t-\pi / 2) \cos t$
8.7.5 (p. 296) $y=1-\cos 2 t+\frac{1}{2} \sin 2 t+\frac{1}{2} u(t-3 \pi) \sin 2 t$
8.7.6 (p.296) $y=4 e^{t}+3 e^{-t}-8+2 u(t-2) \sinh (t-2)$
8.7.7 (p. 296) $y=\frac{1}{2} e^{t}-\frac{7}{2} e^{-t}+2+3 u(t-6)\left(1-e^{-(t-6)}\right)$
8.7.8 (p. 296) $y=e^{2 t}+7 \cos 2 t-\sin 2 t-\frac{1}{2} u(t-\pi / 2) \sin 2 t$
8.7 .9 (p. 296) $y=\frac{1}{2}\left(1+e^{-2 t}\right)+u(t-1)\left(e^{-(t-1)}-e^{-2(t-1)}\right)$
8.7.10 (p. 296) $y=\frac{1}{4} e^{t}+\frac{1}{4} e^{-t}(2 t-5)+2 u(t-2)(t-2) e^{-(t-2)}$
8.7.11 (p. 296) $y=\frac{1}{6}(2 \sin t+5 \sin 2 t)-\frac{1}{2} u(t-\pi / 2) \sin 2 t$
8.7.12 (p. 296) $y=e^{-t}(\sin t-\cos t)-e^{-(t-\pi)} \sin t-3 u(t-2 \pi) e^{-(t-2 \pi)} \sin t$
8.7.13(p.296) $y=e^{-2 t}\left(\cos 3 t+\frac{4}{3} \sin 3 t\right)-\frac{1}{3} u(t-\pi / 6) e^{-2(t-\pi / 6)} \cos 3 t-\frac{2}{3} u(t-\pi / 3) e^{-2(t-\pi / 3)} \sin 3 t$
8.7 .14 (p. 296) $y=\frac{7}{10} e^{2 t}-\frac{6}{5} e^{-t / 2}-\frac{1}{2}+\frac{1}{5} u(t-2)\left(e^{2(t-2)}-e^{-(t-2) / 2}\right)$
8.7.15 (p. 296) $y=\frac{1}{17}(12 \cos t+20 \sin t)+\frac{1}{34} e^{t / 2}(10 \cos t-11 \sin t)-u(t-\pi / 2) e^{(2 t-\pi) / 4} \cos t$

$$
+u(t-\pi) e^{(t-\pi) / 2} \sin t
$$

8.7.16 (p. 296) $y=\frac{1}{3}(\cos t-\cos 2 t-3 \sin t)-2 u(t-\pi / 2) \cos t+3 u(t-\pi) \sin t$
8.7.17 (p. 296) $y=e^{t}-e^{-t}(1+2 t)-5 u(t-1) \sinh (t-1)+3 u(t-2) \sinh (t-2)$
8.7.18 (p. 296) $\left.y=\frac{1}{4}\left(e^{t}-e^{-t}(1+6 t)\right)-u(t-1) e^{-(t-1)}+2 u(t-2) e^{-(t-2)}\right)$
8.7.19 (p. 296) $y=\frac{5}{3} \sin t-\frac{1}{3} \sin 2 t+\frac{1}{3} u(t-\pi)(\sin 2 t+2 \sin t)+u(t-2 \pi) \sin t$
8.7.20 (p. 296) $y=\frac{3}{4} \cos 2 t-\frac{1}{2} \sin 2 t+\frac{1}{4}+\frac{1}{4} u(t-\pi / 2)(1+\cos 2 t)+\frac{1}{2} u(t-\pi) \sin 2 t+\frac{3}{2} u(t-3 \pi / 2) \sin 2 t$
8.7 .21 (p. 297) $y=\cos t-\sin t \quad 8.7 .22$ (p. 297) $y=\frac{1}{4}\left(8 e^{3 t}-12 e^{-2 t}\right)$
8.7.23 (p. 297) $y=5\left(e^{-2 t}-e^{-t}\right) \quad 8.7 .24$ (p. 297) $y=e^{-2 t}(1+6 t)$
8.7.25 (p. 297) $y=\frac{1}{4} e^{-t / 2}(4-19 t)$
8.7.29 (p. 297) $y=(-1)^{k} m \omega_{1} R e^{-c \tau / 2 m} \delta(t-\tau)$ if $\omega_{1} \tau-\phi=(2 k+1) \pi / 2(k=$ integer $)$
8.7 .30 (p. 297) (a) $y=\frac{\left(e^{m+1}-1\right)\left(e^{t-m}-e^{-t}\right)}{2(e-1)}, \quad m \leq t<m+1,(m=0,1, \ldots)$
(b) $y=(m+1) \sin t, 2 m \pi \leq t<2(m+1) \pi,(m=0,1, \ldots)$
(c) $y=e^{2(t-m)} \frac{e^{2 m+2}-1}{e^{2}-1}-e^{(t-m)} \frac{e^{m+1}-1}{e-1}, m \leq t<m+1 \quad(m=0,1, \ldots)$
(d) $y=\left\{\begin{array}{cc}0, & 2 m \pi \leq t<(2 m+1) \pi, \\ -\sin t, & (2 m+1) \pi \leq t<(2 m+2) \pi,\end{array} \quad(m=0,1, \ldots)\right.$

## Section 10.1 Answers, pp. 307-308

10.1.1(p.307) $\begin{array}{lll}Q_{1}^{\prime} & =2-\frac{1}{10} Q_{1}+\frac{1}{25} Q_{2} & 10.1 .2(\mathbf{p . 3 0 7 )} \\ Q_{2}^{\prime} & =6+\frac{3}{50} Q_{1}-\frac{1}{20} Q_{2} . & \begin{array}{l}Q_{1}^{\prime} \\ \end{array} \quad \begin{array}{l}12-\frac{5}{100+2 t} Q_{1}+\frac{1}{100+3 t} Q_{2} \\ Q_{2}^{\prime}\end{array} \\ =5+\frac{1}{50+t} Q_{1}-\frac{4}{100+3 t} Q_{2} .\end{array}$
10.1.3 (p. 307) $m_{1} y_{1}^{\prime \prime}=-\left(c_{1}+c_{2}\right) y_{1}^{\prime}+c_{2} y_{2}^{\prime}-\left(k_{1}+k_{2}\right) y_{1}+k_{2} y_{2}+F_{1}$
$m_{2} y_{2}^{\prime \prime}=\left(c_{2}-c_{3}\right) y_{1}^{\prime}-\left(c_{2}+c_{3}\right) y_{2}^{\prime}+c_{3} y_{3}^{\prime}+\left(k_{2}-k_{3}\right) y_{1}-\left(k_{2}+k_{3}\right) y_{2}+k_{3} y_{3}+F_{2}$
$m_{3} y_{3}^{\prime \prime}=c_{3} y_{1}^{\prime}+c_{3} y_{2}^{\prime}-c_{3} y_{3}^{\prime}+k_{3} y_{1}+k_{3} y_{2}-k_{3} y_{3}+F_{3}$
10.1.4(p. 307) $x^{\prime \prime}=-\frac{\alpha}{m} x^{\prime}+\frac{g R^{2} x}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}} y^{\prime \prime}=-\frac{\alpha}{m} y^{\prime}+\frac{g R^{2} y}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}$
$z^{\prime \prime}=-\frac{\alpha}{m} z^{\prime}+\frac{g R^{2} z}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}$

(c) $\begin{aligned} y_{1}^{\prime} & =y_{2} \\ y_{2}^{\prime} & =y_{3} \\ y_{3}^{\prime} & =f\left(t, y_{1}, y_{2}, y_{3}\right)\end{aligned}$
(d) $\begin{aligned} y_{1}^{\prime} & =y_{2} \\ y_{2}^{\prime} & =y_{3} \\ y_{3}^{\prime} & =y_{4} \\ y_{4}^{\prime} & =f\left(t, y_{1}\right)\end{aligned}$ $x_{1}^{\prime}=x_{2}$
(e) $\begin{aligned} x_{2}^{\prime} & =f\left(t, x_{1}, y_{1}\right) \\ y_{1}^{\prime} & =y_{2}\end{aligned}$ $y_{2}^{\prime}=g\left(t, x_{1}, y_{1}\right)$

$$
x^{\prime}=x_{1}
$$

$$
x_{1}^{\prime}=-\frac{g R^{2} x}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}
$$

$$
y_{1}^{\prime}=-\frac{g R^{2} y}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}
$$

$$
z_{1}^{\prime}=-\frac{g R^{2} z}{\left(x^{2}+y^{2}+z^{2}\right)^{3 / 2}}
$$

## Section 10.2 Answers, pp. 310-313

10.2.1 (p. 310) (a) $\mathbf{y}^{\prime}=\left[\begin{array}{ll}2 & 4 \\ 4 & 2\end{array}\right] \mathbf{y} \quad$ (b) $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-2 & -2 \\ -5 & 1\end{array}\right] \mathbf{y}$ (c) $\mathbf{y}^{\prime}=\left[\begin{array}{rr}-4 & -10 \\ 3 & 7\end{array}\right] \mathbf{y}$ (d) $\mathbf{y}^{\prime}=\left[\begin{array}{ll}2 & 1 \\ 1 & 2\end{array}\right] \mathbf{y}$
$\mathbf{1 0 . 2 . 2}$ (p. 310) (a) $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-1 & 2 & 3 \\ 0 & 1 & 6 \\ 0 & 0 & -2\end{array}\right] \mathbf{y} \quad$ (b) $\mathbf{y}^{\prime}=\left[\begin{array}{lll}0 & 2 & 2 \\ 2 & 0 & 2 \\ 2 & 2 & 0\end{array}\right] \mathbf{y}$ (c) $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}-1 & 2 & 2 \\ 2 & -1 & 2 \\ 2 & 2 & -1\end{array}\right] \mathbf{y} \quad$ (d) $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}3 & -1 & -1 \\ -2 & 3 & 2 \\ 4 & -1 & -2\end{array}\right] \mathbf{y}$
$10.2 .3(\mathbf{p . 3 1 1}) \quad$ (a) $\mathbf{y}^{\prime}=\left[\begin{array}{rr}1 & 1 \\ -2 & 4\end{array}\right] \mathbf{y}, \mathbf{y}(0)=\left[\begin{array}{l}1 \\ 0\end{array}\right] \quad$ (b) $\mathbf{y}^{\prime}=\left[\begin{array}{rl}5 & 3 \\ -1 & 1\end{array}\right] \mathbf{y}, \mathbf{y}(0)=\left[\begin{array}{l}9 \\ -5\end{array}\right]$
$\mathbf{1 0 . 2} \mathbf{4}$ (p. 311) (a) $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}6 & 4 & 4 \\ -7 & -2 & -1 \\ 7 & 4 & 3\end{array}\right] \mathbf{y}, \mathbf{y}(0)=\left[\begin{array}{r}3 \\ -6 \\ 4\end{array}\right]$
(b) $\mathbf{y}^{\prime}=\left[\begin{array}{rrr}8 & 7 & 7 \\ -5 & -6 & -9 \\ 5 & 7 & 10\end{array}\right] \mathbf{y}, \mathbf{y}(0)=\left[\begin{array}{r}2 \\ -4 \\ 3\end{array}\right]$
10.2.5 (p. 311) (a) $\mathbf{y}^{\prime}=\left[\begin{array}{ll}-3 & 2 \\ -5 & 3\end{array}\right]+\left[\begin{array}{l}3-2 t \\ 6-3 t\end{array}\right] \quad$ (b) $\mathbf{y}^{\prime}=\left[\begin{array}{rr}3 & 1 \\ -1 & 1\end{array}\right] \mathbf{y}+\left[\begin{array}{c}-5 e^{t} \\ e^{t}\end{array}\right]$
10.2.10 (p. 313) (a) $\frac{d}{d t} Y^{2}=Y^{\prime} Y+Y Y^{\prime}$
(b) $\frac{d}{d t} Y^{n}=Y^{\prime} Y^{n-1}+Y Y^{\prime} Y^{n-2}+Y^{2} Y^{\prime} Y^{n-3}+\cdots+Y^{n-1} Y^{\prime}=\sum_{r=0}^{n-1} Y^{r} Y^{\prime} Y^{n-r-1}$
10.2.13 (p. 313) $B=\left(P^{\prime}+P A\right) P^{-1}$.

Section 10.3 Answers, pp. 317-320
10.3.2 (p.317) $\mathbf{y}^{\prime}=\left[\begin{array}{cc}0 & 1 \\ -\frac{P_{2}(x)}{P_{0}(x)} & -\frac{P_{1}(x)}{P_{0}(x)}\end{array}\right] \mathbf{y} \quad 10.3 .3$ (p.317) $\mathbf{y}^{\prime}=\left[\begin{array}{cccc}0 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 \\ -\frac{P_{n}(x)}{P_{0}(x)} & -\frac{P_{n-1}(x)}{P_{0}(x)} & \cdots & -\frac{P_{1}(x)}{P_{0}(x)}\end{array}\right] \mathbf{y}$
10.3.7 (p. 318) (b) $\mathbf{y}=\left[\begin{array}{l}3 e^{6 t}-6 e^{-2 t} \\ 3 e^{6 t}+6 e^{-2 t}\end{array}\right]$
(c) $\mathbf{y}=\frac{1}{2}\left[\begin{array}{ll}e^{6 t}+e^{-2 t} & e^{6 t}-e^{-2 t} \\ e^{6 t}-e^{-2 t} & e^{6 t}+e^{-2 t}\end{array}\right] \mathbf{k}$
10.3 .8 (p. 319) (b) $\mathbf{y}=\left[\begin{array}{c}6 e^{-4 t}+4 e^{3 t} \\ 6 e^{-4 t}-10 e^{3 t}\end{array}\right]$
(c) $\mathbf{y}=\frac{1}{7}\left[\begin{array}{ll}5 e^{-4 t}+2 e^{3 t} & 2 e^{-4 t}-2 e^{3 t} \\ 5 e^{-4 t}-5 e^{3 t} & 2 e^{-4 t}+5 e^{3 t}\end{array}\right] \mathbf{k}$
10.3 .9 (p. 319) (b) $\mathbf{y}=\left[\begin{array}{c}-15 e^{2 t}-4 e^{t} \\ 9 e^{2 t}+2 e^{t}\end{array}\right]$
(c) $\mathbf{y}=\left[\begin{array}{cc}-5 e^{2 t}+6 e^{t} & -10 e^{2 t}+10 e^{t} \\ 3 e^{2 t}-3 e^{t} & 6 e^{2 t}-5 e^{t}\end{array}\right] \mathbf{k}$
10.3 .10 (p. 319) (b) $\mathbf{y}=\left[\begin{array}{l}5 e^{3 t}-3 e^{t} \\ 5 e^{3 t}+3 e^{t}\end{array}\right]$
(c) $\mathbf{y}=\frac{1}{2}\left[\begin{array}{cc}e^{3 t}+e^{t} & e^{3 t}-e^{t} \\ e^{3 t}-e^{t} & e^{3 t}+e^{t}\end{array}\right] \mathbf{k}$
10.3.11 (p.319) (b) $\mathbf{y}=\left[\begin{array}{c}e^{2 t}-2 e^{3 t}+3 e^{-t} \\ 2 e^{3 t}-9 e^{-t} \\ e^{2 t}-2 e^{3 t}+21 e^{-t}\end{array}\right] \quad$ (c) $\mathbf{y}=\frac{1}{6}\left[\begin{array}{ccc}4 e^{2 t}+3 e^{3 t}-e^{-t} & 6 e^{2 t}-6 e^{3 t} & 2 e^{2 t}-3 e^{3 t}+e^{-t} \\ -3 e^{3 t}+3 e^{-t} & 6 e^{3 t} & 3 e^{3 t}-3 e^{-t} \\ 4 e^{2 t}+3 e^{3 t}-7 e^{-t} & 6 e^{2 t}-6 e^{3 t} & 2 e^{2 t}-3 e^{3 t}+7 e^{-t}\end{array}\right] \mathbf{k}$
10.3.12 (p.319) (b) $\mathbf{y}=\frac{1}{3}\left[\begin{array}{c}-e^{-2 t}+e^{4 t} \\ -10 e^{-2 t}+e^{4 t} \\ 11 e^{-2 t}+e^{4 t}\end{array}\right]$
(c) $\mathbf{y}=\frac{1}{3}\left[\begin{array}{ccc}2 e^{-2 t}+e^{4 t} & -e^{-2 t}+e^{4 t} & -e^{-2 t}+e^{4 t} \\ -e^{-2 t}+e^{4 t} & 2 e^{-2 t}+e^{4 t} & -e^{-2 t}+e^{4 t} \\ -e^{-2 t}+e^{4 t} & -e^{-2 t}+e^{4 t} & 2 e^{-2 t}+e^{4 t}\end{array}\right] \mathbf{k}$
$\mathbf{1 0 . 3 . 1 3}$ (p. 319) (b) $\mathbf{y}=\left[\begin{array}{c}3 e^{t}+3 e^{-t}-e^{-2 t} \\ 3 e^{t}+2 e^{-2 t} \\ -e^{-2 t}\end{array}\right]$
(c) $\mathbf{y}=\left[\begin{array}{ccc}e^{-t} & e^{t}-e^{-t} & 2 e^{t}-3 e^{-t}+e^{-2 t} \\ 0 & e^{t} & 2 e^{t}-2 e^{-2 t} \\ 0 & 0 & e^{-2 t}\end{array}\right] \mathbf{k}$
10.3.14 (p. 319) $Y Z^{-1}$ and $Z Y^{-1}$

Section 10.4 Answers, pp. 328-330
10.4.1 (p. 328) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{3 t}+c_{2}\left[\begin{array}{r}1 \\ -1\end{array}\right] e^{-t} \quad 10.4 .2(\mathbf{p . 3 2 8}) \mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{-t / 2}+c_{2}\left[\begin{array}{r}-1 \\ 1\end{array}\right] e^{-2 t}$ 10.4.3 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}-3 \\ 1\end{array}\right] e^{-t}+c_{2}\left[\begin{array}{r}-1 \\ 2\end{array}\right] e^{-2 t} \mathbf{1 0 . 4 . 4}$ (p.329) $\mathbf{y}=c_{1}\left[\begin{array}{l}2 \\ 1\end{array}\right] e^{-3 t}+c_{2}\left[\begin{array}{r}-2 \\ 1\end{array}\right] e^{t}$ 10.4.5 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{-2 t}+c_{1}\left[\begin{array}{c}-4 \\ 1\end{array}\right] e^{3 t} \quad 10.4 .6(\mathbf{p . 3 2 9}) \mathbf{y}=c_{1}\left[\begin{array}{l}3 \\ 2\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{l}1 \\ 1\end{array}\right] e^{t}$ 10.4.7 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}-3 \\ 1\end{array}\right] e^{-5 t}+c_{2}\left[\begin{array}{r}-1 \\ 1\end{array}\right] e^{-3 t}$
10.4.8 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 2 \\ 1\end{array}\right] e^{-3 t}+c_{2}\left[\begin{array}{c}-1 \\ -4 \\ 1\end{array}\right] e^{-t}+c_{3}\left[\begin{array}{c}-1 \\ -1 \\ 1\end{array}\right] e^{2 t}$
10.4.9 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{l}2 \\ 1 \\ 2\end{array}\right] e^{-16 t}+c_{2}\left[\begin{array}{r}-1 \\ 2 \\ 0\end{array}\right] e^{2 t}+c_{3}\left[\begin{array}{r}-1 \\ 0 \\ 1\end{array}\right] e^{2 t}$
$\mathbf{1 0 . 4 . 1 0}$ (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}-2 \\ -4 \\ 3\end{array}\right] e^{t}+c_{2}\left[\begin{array}{r}-1 \\ 1 \\ 0\end{array}\right] e^{-2 t}+c_{3}\left[\begin{array}{r}-7 \\ -5 \\ 4\end{array}\right] e^{2 t}$
10.4.11 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ -1 \\ 1\end{array}\right] e^{-2 t}+c_{2}\left[\begin{array}{r}-1 \\ -2 \\ 1\end{array}\right] e^{-3 t}+c_{3}\left[\begin{array}{r}-2 \\ -6 \\ 3\end{array}\right] e^{-5 t}$
10.4.12 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}11 \\ 7 \\ 1\end{array}\right] e^{3 t}+c_{2}\left[\begin{array}{l}1 \\ 2 \\ 1\end{array}\right] e^{-2 t}+c_{3}\left[\begin{array}{l}1 \\ 1 \\ 1\end{array}\right] e^{-t}$
10.4.13 (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}4 \\ -1 \\ 1\end{array}\right] e^{-4 t}+c_{2}\left[\begin{array}{r}-1 \\ -1 \\ 1\end{array}\right] e^{6 t}+c_{3}\left[\begin{array}{r}-1 \\ 0 \\ 1\end{array}\right] e^{4 t}$
10.4.14 (p.329) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1 \\ 5\end{array}\right] e^{-5 t}+c_{2}\left[\begin{array}{r}-1 \\ 0 \\ 1\end{array}\right] e^{5 t}+c_{3}\left[\begin{array}{l}1 \\ 1 \\ 0\end{array}\right] e^{5 t}$
$\mathbf{1 0 . 4 . 1 5}$ (p. 329) $\mathbf{y}=c_{1}\left[\begin{array}{r}1 \\ -1 \\ 2\end{array}\right]+c_{2}\left[\begin{array}{r}-1 \\ 0 \\ 3\end{array}\right] e^{6 t}+c_{3}\left[\begin{array}{l}1 \\ 3 \\ 0\end{array}\right] e^{6 t}$
10.4.16 (p. 329) $\mathbf{y}=-\left[\begin{array}{l}2 \\ 6\end{array}\right] e^{5 t}+\left[\begin{array}{l}4 \\ 2\end{array}\right] e^{-5 t} 10.4 .17(\mathbf{p} .329) \mathbf{y}=\left[\begin{array}{l}2 \\ -4\end{array}\right] e^{t / 2}+\left[\begin{array}{l}-2 \\ 1\end{array}\right] e^{t}$
10.4.18 (p. 329) $\mathbf{y}=\left[\begin{array}{l}7 \\ 7\end{array}\right] e^{9 t}-\left[\begin{array}{l}2 \\ 4\end{array}\right] e^{-3 t} \mathbf{1 0 . 4 . 1 9}$ (p. 329) $\mathbf{y}=\left[\begin{array}{l}3 \\ 9\end{array}\right] e^{5 t}-\left[\begin{array}{l}4 \\ 2\end{array}\right] e^{-5 t}$
$\mathbf{1 0 . 4 . 2 0} \mathbf{( p . 3 2 9}) \mathbf{y}=\left[\begin{array}{l}5 \\ 5 \\ 0\end{array}\right] e^{t / 2}+\left[\begin{array}{l}0 \\ 0 \\ 1\end{array}\right] e^{t / 2}+\left[\begin{array}{r}-1 \\ 2 \\ 0\end{array}\right] e^{-t / 2} \quad \mathbf{1 0 . 4 . 2 1}(\mathbf{p . 3 2 9}) \mathbf{y}=\left[\begin{array}{l}3 \\ 3 \\ 3\end{array}\right] e^{t}+\left[\begin{array}{r}-2 \\ -2 \\ 2\end{array}\right] e^{-t}$
10.4.22 (p. 329) $\mathbf{y}=\left[\begin{array}{r}2 \\ -2 \\ 2\end{array}\right] e^{t}-\left[\begin{array}{l}3 \\ 0 \\ 3\end{array}\right] e^{-2 t}+\left[\begin{array}{l}1 \\ 1 \\ 0\end{array}\right] e^{3 t}$
10.4.23 (p. 329) $\mathbf{y}=-\left[\begin{array}{l}1 \\ 2 \\ 1\end{array}\right] e^{t}+\left[\begin{array}{l}4 \\ 2 \\ 4\end{array}\right] e^{-t}+\left[\begin{array}{l}1 \\ 1 \\ 0\end{array}\right] e^{2 t}$
10.4.24 (p. 329) $\mathbf{y}=\left[\begin{array}{r}-2 \\ -2 \\ 2\end{array}\right] e^{2 t}-\left[\begin{array}{l}0 \\ 3 \\ 0\end{array}\right] e^{-2 t}+\left[\begin{array}{r}4 \\ 12 \\ 4\end{array}\right] e^{4 t}$
10.4.25 (p. 329) $\mathbf{y}=\left[\begin{array}{r}-1 \\ -1 \\ 1\end{array}\right] e^{-6 t}+\left[\begin{array}{r}2 \\ -2 \\ 2\end{array}\right] e^{2 t}+\left[\begin{array}{r}7 \\ -7 \\ -7\end{array}\right] e^{4 t}$
10.4.26 (p.330) $\mathbf{y}=\left[\begin{array}{l}1 \\ 4 \\ 4\end{array}\right] e^{-t}+\left[\begin{array}{r}6 \\ 6 \\ -2\end{array}\right] e^{2 t} \mathbf{1 0 . 4 . 2 7}(\mathbf{p} .330) \mathbf{y}=\left[\begin{array}{r}4 \\ -2 \\ 2\end{array}\right]+\left[\begin{array}{r}3 \\ -9 \\ 6\end{array}\right] e^{4 t}+\left[\begin{array}{r}-1 \\ 1 \\ -1\end{array}\right] e^{2 t}$
10.4.29 (p. 330) Half lines of $L_{1}: y_{2}=y_{1}$ and $L_{2}: y_{2}=-y_{1}$ are trajectories other trajectories are asymptotically tangent to $L_{1}$ as $t \rightarrow-\infty$ and asymptotically tangent to $L_{2}$ as $t \rightarrow \infty$.
10.4.30 (p. 330) Half lines of $L_{1}: y_{2}=-2 y_{1}$ and $L_{2}: y_{2}=-y_{1} / 3$ are trajectories other trajectories are asymptotically parallel to $L_{1}$ as $t \rightarrow-\infty$ and asymptotically tangent to $L_{2}$ as $t \rightarrow \infty$.
10.4.31 (p. 330) Half lines of $L_{1}: y_{2}=y_{1} / 3$ and $L_{2}: y_{2}=-y_{1}$ are trajectories other trajectories are asymptotically tangent to $L_{1}$ as $t \rightarrow-\infty$ and asymptotically parallel to $L_{2}$ as $t \rightarrow \infty$.
10.4.32 (p. 330) Half lines of $L_{1}: y_{2}=y_{1} / 2$ and $L_{2}: y_{2}=-y_{1}$ are trajectories other trajectories are asymptotically tangent to $L_{1}$ as $t \rightarrow-\infty$ and asymptotically tangent to $L_{2}$ as $t \rightarrow \infty$.
10.4.33(p.330) Half lines of $L_{1}: y_{2}=-y_{1} / 4$ and $L_{2}: y_{2}=-y_{1}$ are trajectories other trajectories are asymptotically tangent to $L_{1}$ as $t \rightarrow-\infty$ and asymptotically parallel to $L_{2}$ as $t \rightarrow \infty$.
10.4.34 (p. 330) Half lines of $L_{1}: y_{2}=-y_{1}$ and $L_{2}: y_{2}=3 y_{1}$ are trajectories other trajectories are asymptotically parallel to $L_{1}$ as $t \rightarrow-\infty$ and asymptotically tangent to $L_{2}$ as $t \rightarrow \infty$.
10.4.36 (p. 330) Points on $L_{2}: y_{2}=y_{1}$ are trajectories of constant solutions. The trajectories of nonconstant solutions are half-lines on either side of $L_{1}$, parallel to $\left[\begin{array}{r}1 \\ -1\end{array}\right]$, traversed toward $L_{1}$.
10.4.37 (p. 330) Points on $L_{1}: y_{2}=-y_{1} / 3$ are trajectories of constant solutions. The trajectories of nonconstant solutions are half-lines on either side of $L_{1}$, parallel to $\left[\begin{array}{r}-1 \\ 2\end{array}\right]$, traversed away from $L_{1}$.
10.4.38 (p. 330) Points on $L_{1}: y_{2}=y_{1} / 3$ are trajectories of constant solutions. The trajectories of nonconstant solutions are half-lines on either side of $L_{1}$, parallel to $\left[\begin{array}{r}1 \\ -1\end{array}\right],\left[\begin{array}{l}1 \\ 1\end{array}\right]-1$, traversed away from $L_{1}$.
10.4.39 (p. 330) Points on $L_{1}: y_{2}=y_{1} / 2$ are trajectories of constant solutions. The trajectories of nonconstant solutions are half-lines on either side of $L_{1}$, parallel to $\left[\begin{array}{r}1 \\ -1\end{array}\right], L_{1}$.
10.4.40 (p. 330) Points on $L_{2}: y_{2}=-y_{1}$ are trajectories of constant solutions. The trajectories of nonconstant solutions are half-lines on either side of $L_{2}$, parallel to $\left[\begin{array}{r}-4 \\ 1\end{array}\right]$, traversed toward $L_{1}$.
10.4.41 (p. 330) Points on $L_{1}: y_{2}=3 y_{1}$ are trajectories of constant solutions. The trajectories of nonconstant solutions are half-lines on either side of $L_{1}$, parallel to $\left[\begin{array}{r}1 \\ -1\end{array}\right]$, traversed away from $L_{1}$.

## Section 10.5 Answers, pp. 342-344

10.5.1 (p. 342) $\mathbf{y}=c_{1}\left[\begin{array}{l}2 \\ 1\end{array}\right] e^{5 t}+c_{2}\left(\left[\begin{array}{r}-1 \\ 0\end{array}\right] e^{5 t}+\left[\begin{array}{l}2 \\ 1\end{array}\right] t e^{5 t}\right)$.

$$
\begin{aligned}
& \text { 10.5.2 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{l}
1 \\
1
\end{array}\right] e^{-t}+c_{2}\left(\left[\begin{array}{l}
1 \\
0
\end{array}\right] e^{-t}+\left[\begin{array}{l}
1 \\
1
\end{array}\right] t e^{-t}\right) \\
& \text { 10.5.3 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{r}
-2 \\
1
\end{array}\right] e^{-9 t}+c_{2}\left(\left[\begin{array}{r}
-1 \\
0
\end{array}\right] e^{-9 t}+\left[\begin{array}{r}
-2 \\
1
\end{array}\right] t e^{-9 t}\right) \\
& \text { 10.5.4 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{r}
-1 \\
1
\end{array}\right] e^{2 t}+c_{2}\left(\left[\begin{array}{r}
-1 \\
0
\end{array}\right] e^{2 t}+\left[\begin{array}{r}
-1 \\
1
\end{array}\right] t e^{2 t}\right) \\
& \text { 10.5.5 (p. 342) } c_{1}\left[\begin{array}{r}
-2 \\
1
\end{array}\right]+c_{2}\left(\left[\begin{array}{r}
-1 \\
0
\end{array}\right] \frac{e^{-2 t}}{3}+\left[\begin{array}{r}
-2 \\
1
\end{array}\right] t e^{-2 t}\right) \\
& \text { 10.5.6 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{l}
3 \\
2
\end{array}\right] e^{-4 t}+c_{2}\left(\left[\begin{array}{r}
-1 \\
0
\end{array}\right] \frac{e^{-4 t}}{2}+\left[\begin{array}{l}
3 \\
2
\end{array}\right] t e^{-4 t}\right) \\
& \text { 10.5.7 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{l}
4 \\
3
\end{array}\right] e^{-t}+c_{2}\left(\left[\begin{array}{r}
-1 \\
0
\end{array}\right] \frac{e^{-t}}{3}+\left[\begin{array}{l}
4 \\
3
\end{array}\right] t e^{-t}\right) \\
& \text { 10.5.8 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{r}
-1 \\
-1 \\
2
\end{array}\right]+c_{2}\left[\begin{array}{l}
1 \\
1 \\
2
\end{array}\right] e^{4 t}+c_{3}\left(\left[\begin{array}{l}
0 \\
1 \\
0
\end{array}\right] \frac{e^{4 t}}{2}+\left[\begin{array}{l}
1 \\
1 \\
2
\end{array}\right] t e^{4 t}\right) \\
& \text { 10.5.9 (p.342) } \mathbf{y}=c_{1}\left[\begin{array}{r}
-1 \\
1 \\
1
\end{array}\right] e^{t}+c_{2}\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] e^{-t}+c_{3}\left(\left[\begin{array}{l}
0 \\
3 \\
0
\end{array}\right] e^{-t}+\left[\begin{array}{r}
1 \\
-1 \\
1
\end{array}\right] t e^{-t}\right) \text {. } \\
& \text { 10.5.10 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{l}
0 \\
1 \\
1
\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] e^{-2 t}+c_{3}\left(\left[\begin{array}{l}
1 \\
1 \\
0
\end{array}\right] \frac{e^{-2 t}}{2}+\left[\begin{array}{l}
1 \\
0 \\
1
\end{array}\right] t e^{-2 t}\right) \\
& \text { 10.5.11 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{r}
-2 \\
-3 \\
1
\end{array}\right] e^{2 t}+c_{2}\left[\begin{array}{r}
0 \\
-1 \\
1
\end{array}\right] e^{4 t}+c_{3}\left(\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right] \frac{e^{4 t}}{2}+\left[\begin{array}{r}
0 \\
-1 \\
1
\end{array}\right] t e^{4 t}\right) \\
& \text { 10.5.12 (p. 342) } \mathbf{y}=c_{1}\left[\begin{array}{r}
-1 \\
-1 \\
1
\end{array}\right] e^{-2 t}+c_{2}\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right] e^{4 t}+c_{3}\left(\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right] \frac{e^{4 t}}{2}+\left[\begin{array}{l}
1 \\
1 \\
1
\end{array}\right] t e^{4 t}\right) \text {. } \\
& \text { 10.5.13 (p. 342) } \mathbf{y}=\left[\begin{array}{l}
6 \\
2
\end{array}\right] e^{-7 t}-\left[\begin{array}{l}
8 \\
4
\end{array}\right] t e^{-7 t} 10.5 .14(\mathbf{p . 3 4 2}) \mathbf{y}=\left[\begin{array}{l}
5 \\
8
\end{array}\right] e^{3 t}-\left[\begin{array}{l}
12 \\
16
\end{array}\right] t e^{3 t} \\
& \text { 10.5.15 (p. 342) } \mathbf{y}=\left[\begin{array}{l}
2 \\
3
\end{array}\right] e^{-5 t}-\left[\begin{array}{l}
8 \\
4
\end{array}\right] t e^{-5 t} \mathbf{1 0 . 5 . 1 6}(\mathbf{p} .342) \mathbf{y}=\left[\begin{array}{l}
3 \\
1
\end{array}\right] e^{5 t}-\left[\begin{array}{l}
12 \\
6
\end{array}\right] t e^{5 t} \\
& \text { 10.5.17 (p. 342) } \mathbf{y}=\left[\begin{array}{l}
0 \\
2
\end{array}\right] e^{-4 t}+\left[\begin{array}{l}
6 \\
6
\end{array}\right] t e^{-4 t} \\
& \text { 10.5.18 (p. 342) } \mathbf{y}=\left[\begin{array}{r}
4 \\
8 \\
-6
\end{array}\right] e^{t}+\left[\begin{array}{r}
2 \\
-3 \\
-1
\end{array}\right] e^{-2 t}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] t e^{-2 t} \\
& \text { 10.5.19 (p. 342) } \mathbf{y}=\left[\begin{array}{l}
3 \\
3 \\
6
\end{array}\right] e^{2 t}-\left[\begin{array}{l}
9 \\
5 \\
6
\end{array}\right]+\left[\begin{array}{l}
2 \\
2 \\
0
\end{array}\right] t \\
& \mathbf{1 0 . 5 . 2 0} \text { (p. 342) } \mathbf{y}=-\left[\begin{array}{l}
2 \\
0 \\
2
\end{array}\right] e^{-3 t}+\left[\begin{array}{r}
-4 \\
9 \\
1
\end{array}\right] e^{t}-\left[\begin{array}{l}
0 \\
4 \\
4
\end{array}\right] t e^{t} \\
& \text { 10.5.21 (p. 342) } \mathbf{y}=\left[\begin{array}{r}
-2 \\
2 \\
2
\end{array}\right] e^{4 t}+\left[\begin{array}{r}
0 \\
-1 \\
1
\end{array}\right] e^{2 t}+\left[\begin{array}{r}
3 \\
-3 \\
3
\end{array}\right] t e^{2 t} \\
& \text { 10.5.22 (p. 342) } \mathbf{y}=-\left[\begin{array}{l}
1 \\
1 \\
0
\end{array}\right] e^{-4 t}+\left[\begin{array}{r}
-3 \\
2 \\
-3
\end{array}\right] e^{8 t}+\left[\begin{array}{r}
8 \\
0 \\
-8
\end{array}\right] t e^{8 t} \\
& \text { 10.5.23 (p. 343) } \mathbf{y}=\left[\begin{array}{l}
3 \\
6 \\
3
\end{array}\right] e^{4 t}-\left[\begin{array}{l}
3 \\
4 \\
1
\end{array}\right]+\left[\begin{array}{l}
8 \\
4 \\
4
\end{array}\right] t
\end{aligned}
$$

$$
\begin{aligned}
& 10.5 .24\left(\mathbf{p . 3 4 3 )} \mathbf{y}=c_{1}\left[\begin{array}{l}
0 \\
1 \\
1
\end{array}\right] e^{6 t}+c_{2}\left(\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{e^{6 t}}{4}+\left[\begin{array}{l}
0 \\
1 \\
1
\end{array}\right] t e^{6 t}\right)\right. \\
& +c_{3}\left(\left[\begin{array}{l}
1 \\
1 \\
0
\end{array}\right] \frac{e^{6 t}}{8}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{t e^{6 t}}{4}+\left[\begin{array}{l}
0 \\
1 \\
1
\end{array}\right] \frac{t^{2} e^{6 t}}{2}\right) \\
& \begin{aligned}
\mathbf{1 0 . 5 . 2 5}\left(\mathbf{p . ~ 3 4 3 )} \mathbf{y}=c_{1}\right.
\end{aligned} \\
& \left.+\begin{array}{r}
-1 \\
1 \\
1
\end{array}\right] e^{3 t}+c_{2}\left(\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right] \frac{e^{3 t}}{2}+\left[\begin{array}{r}
-1 \\
1 \\
1
\end{array}\right] t e^{3 t}\right) \\
& \\
& +c_{3}\left(\left[\begin{array}{l}
1 \\
2 \\
0
\end{array}\right] \frac{e^{3 t}}{36}+\left[\begin{array}{l}
1 \\
0 \\
0
\end{array}\right] \frac{t e^{3 t}}{2}+\left[\begin{array}{r}
-1 \\
1 \\
1
\end{array}\right] \frac{t^{2} e^{3 t}}{2}\right)
\end{aligned}
$$

10.5.26 (p. 343) $\mathbf{y}=c_{1}\left[\begin{array}{r}0 \\ -1 \\ 1\end{array}\right] e^{-2 t}+c_{2}\left(\left[\begin{array}{r}-1 \\ 1 \\ 0\end{array}\right] e^{-2 t}+\left[\begin{array}{r}0 \\ -1 \\ 1\end{array}\right] t e^{-2 t}\right)$

$$
+c_{3}\left(\left[\begin{array}{r}
3 \\
-2 \\
0
\end{array}\right] \frac{e^{-2 t}}{4}+\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] t e^{-2 t}+\left[\begin{array}{r}
0 \\
-1 \\
1
\end{array}\right] \frac{t^{2} e^{-2 t}}{2}\right)
$$

10.5.27 (p. 343) $\mathbf{y}=c_{1}\left[\begin{array}{l}0 \\ 1 \\ 1\end{array}\right] e^{2 t}+c_{2}\left(\left[\begin{array}{l}1 \\ 1 \\ 0\end{array}\right] \frac{e^{2 t}}{2}+\left[\begin{array}{l}0 \\ 1 \\ 1\end{array}\right] t e^{2 t}\right)$

$$
+c_{3}\left(\left[\begin{array}{r}
-1 \\
1 \\
0
\end{array}\right] \frac{e^{2 t}}{8}+\left[\begin{array}{l}
1 \\
1 \\
0
\end{array}\right] \frac{t e^{2 t}}{2}+\left[\begin{array}{l}
0 \\
1 \\
1
\end{array}\right] \frac{t^{2} e^{2 t}}{2}\right)
$$

10.5.28 (p. 343) $\mathbf{y}=c_{1}\left[\begin{array}{r}-2 \\ 1 \\ 2\end{array}\right] e^{-6 t}+c_{2}\left(-\left[\begin{array}{l}6 \\ 1 \\ 0\end{array}\right] \frac{e^{-6 t}}{6}+\left[\begin{array}{r}-2 \\ 1 \\ 2\end{array}\right] t e^{-6 t}\right)$

$$
+c_{3}\left(-\left[\begin{array}{r}
12 \\
1 \\
0
\end{array}\right] \frac{e^{-6 t}}{36}-\left[\begin{array}{l}
6 \\
1 \\
0
\end{array}\right] \frac{t e^{-6 t}}{6}+\left[\begin{array}{r}
-2 \\
1 \\
2
\end{array}\right] \frac{t^{2} e^{-6 t}}{2}\right)
$$

10.5.29 (p. 343) $\mathbf{y}=c_{1}\left[\begin{array}{r}-4 \\ 0 \\ 1\end{array}\right] e^{-3 t}+c_{2}\left[\begin{array}{l}6 \\ 1 \\ 0\end{array}\right] e^{-3 t}+c_{3}\left(\left[\begin{array}{l}1 \\ 0 \\ 0\end{array}\right] e^{-3 t}+\left[\begin{array}{l}2 \\ 1 \\ 1\end{array}\right] t e^{-3 t}\right)$
$10.5 .30(\mathbf{p . 3 4 3}) \mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ 0 \\ 1\end{array}\right] e^{-3 t}+c_{2}\left[\begin{array}{l}0 \\ 1 \\ 0\end{array}\right] e^{-3 t}+c_{3}\left(\left[\begin{array}{l}1 \\ 0 \\ 0\end{array}\right] e^{-3 t}+\left[\begin{array}{r}-1 \\ -1 \\ 1\end{array}\right] t e^{-3 t}\right)$
10.5.31 (p. 343) $\mathbf{y}=c_{1}\left[\begin{array}{l}2 \\ 0 \\ 1\end{array}\right] e^{-t}+c_{2}\left[\begin{array}{r}-3 \\ 2 \\ 0\end{array}\right] e^{-t}+c_{3}\left(\left[\begin{array}{l}1 \\ 0 \\ 0\end{array}\right] \frac{e^{-t}}{2}+\left[\begin{array}{r}-1 \\ 2 \\ 1\end{array}\right] t e^{-t}\right)$
10.5.32 (p. 343) $\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ 1 \\ 0\end{array}\right] e^{-2 t}+c_{2}\left[\begin{array}{l}0 \\ 0 \\ 1\end{array}\right] e^{-2 t}+c_{3}\left(\left[\begin{array}{r}-1 \\ 0 \\ 0\end{array}\right] e^{-2 t}+\left[\begin{array}{r}1 \\ -1 \\ 1\end{array}\right] t e^{-2 t}\right)$

Section 10.6 Answers, pp. 352-353
10.6.1 (p. 352) $\mathbf{y}=c_{1} e^{2 t}\left[\begin{array}{c}3 \cos t+\sin t \\ 5 \cos t\end{array}\right]+c_{2} e^{2 t}\left[\begin{array}{c}3 \sin t-\cos t \\ 5 \sin t\end{array}\right]$.
10.6.2 (p. 352) $\mathbf{y}=c_{1} e^{-t}\left[\begin{array}{c}5 \cos 2 t+\sin 2 t \\ 13 \cos 2 t\end{array}\right]+c_{2} e^{-t}\left[\begin{array}{c}5 \sin 2 t-\cos 2 t \\ 13 \sin 2 t\end{array}\right]$.
10.6.3 (p. 352) $\mathbf{y}=c_{1} e^{3 t}\left[\begin{array}{c}\cos 2 t+\sin 2 t \\ 2 \cos 2 t\end{array}\right]+c_{2} e^{3 t}\left[\begin{array}{c}\sin 2 t-\cos 2 t \\ 2 \sin 2 t\end{array}\right]$.
10.6.4 (p. 352) $\mathbf{y}=c_{1} e^{2 t}\left[\begin{array}{c}\cos 3 t-\sin 3 t \\ \cos 3 t\end{array}\right]+c_{2} e^{2 t}\left[\begin{array}{c}\sin 3 t+\cos 3 t \\ \sin 3 t\end{array}\right]$.
10.6.5 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ -1 \\ 2\end{array}\right] e^{-2 t}+c_{2} e^{4 t}\left[\begin{array}{c}\cos 2 t-\sin 2 t \\ \cos 2 t+\sin 2 t \\ 2 \cos 2 t\end{array}\right]+c_{3} e^{4 t}\left[\begin{array}{c}\sin 2 t+\cos 2 t \\ \sin 2 t-\cos 2 t \\ 2 \sin 2 t\end{array}\right]$.
10.6.6 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{c}-1 \\ -1 \\ 1\end{array}\right] e^{-t}+c_{2} e^{-2 t}\left[\begin{array}{c}\cos 2 t-\sin 2 t \\ -\cos 2 t-\sin 2 t \\ 2 \cos 2 t\end{array}\right]+c_{3} e^{-2 t}\left[\begin{array}{c}\sin 2 t+\cos 2 t \\ -\sin 2 t+\cos 2 t \\ 2 \sin 2 t\end{array}\right]$
10.6.7 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1 \\ 1\end{array}\right] e^{2 t}+c_{2} e^{t}\left[\begin{array}{r}-\sin t \\ \sin t \\ \cos t\end{array}\right]+c_{3} e^{t}\left[\begin{array}{r}\cos t \\ -\cos t \\ \sin t\end{array}\right]$
10.6.8 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ 1 \\ 1\end{array}\right] e^{t}+c_{2} e^{-t}\left[\begin{array}{c}-\sin 2 t-\cos 2 t \\ 2 \cos 2 t \\ 2 \cos 2 t\end{array}\right]+c_{3} e^{-t}\left[\begin{array}{c}\cos 2 t-\sin 2 t \\ 2 \sin 2 t \\ 2 \sin 2 t\end{array}\right]$
10.6.9 (p. 352) $\mathbf{y}=c_{1} e^{3 t}\left[\begin{array}{c}\cos 6 t-3 \sin 6 t \\ 5 \cos 6 t\end{array}\right]+c_{2} e^{3 t}\left[\begin{array}{c}\sin 6 t+3 \cos 6 t \\ 5 \sin 6 t\end{array}\right]$
10.6.10 (p. 352) $\mathbf{y}=c_{1} e^{2 t}\left[\begin{array}{c}\cos t-3 \sin t \\ 2 \cos t\end{array}\right]+c_{2} e^{2 t}\left[\begin{array}{c}\sin t+3 \cos t \\ 2 \sin t\end{array}\right]$
10.6.11 (p. 352) $\mathbf{y}=c_{1} e^{2 t}\left[\begin{array}{c}3 \sin 3 t-\cos 3 t \\ 5 \cos 3 t\end{array}\right]+c_{2} e^{2 t}\left[\begin{array}{c}-3 \cos 3 t-\sin 3 t \\ 5 \sin 3 t\end{array}\right]$
10.6.12 (p. 352) $\mathbf{y}=c_{1} e^{2 t}\left[\begin{array}{c}\sin 4 t-8 \cos 4 t \\ 5 \cos 4 t\end{array}\right]+c_{2} e^{2 t}\left[\begin{array}{c}-\cos 4 t-8 \sin 4 t \\ 5 \sin 4 t\end{array}\right]$
10.6.13 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{r}-1 \\ 1 \\ 1\end{array}\right] e^{-2 t}+c_{2} e^{t}\left[\begin{array}{r}\sin t \\ -\cos t \\ \cos t\end{array}\right]+c_{3} e^{t}\left[\begin{array}{c}-\cos t \\ -\sin t \\ \sin t\end{array}\right]$
10.6.14 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{l}2 \\ 2 \\ 1\end{array}\right] e^{-2 t}+c_{2} e^{2 t}\left[\begin{array}{c}-\cos 3 t-\sin 3 t \\ -\sin 3 t \\ \cos 3 t\end{array}\right]+c_{3} e^{2 t}\left[\begin{array}{c}-\sin 3 t+\cos 3 t \\ \cos 3 t \\ \sin 3 t\end{array}\right]$
10.6.15 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 2 \\ 1\end{array}\right] e^{3 t}+c_{2} e^{6 t}\left[\begin{array}{r}-\sin 3 t \\ \sin 3 t \\ \cos 3 t\end{array}\right]+c_{3} e^{6 t}\left[\begin{array}{r}\cos 3 t \\ -\cos 3 t \\ \sin 3 t\end{array}\right]$
10.6.16 (p. 352) $\mathbf{y}=c_{1}\left[\begin{array}{l}1 \\ 1 \\ 1\end{array}\right] e^{t}+c_{2} e^{t}\left[\begin{array}{c}2 \cos t-2 \sin t \\ \cos t-\sin t \\ 2 \cos t\end{array}\right]+c_{3} e^{t}\left[\begin{array}{c}2 \sin t+2 \cos t \\ \cos t+\sin t \\ 2 \sin t\end{array}\right]$
$\mathbf{1 0 . 6 . 1 7}$ (p. 352) $\mathbf{y}=e^{t}\left[\begin{array}{c}5 \cos 3 t+\sin 3 t \\ 2 \cos 3 t+3 \sin 3 t\end{array}\right] \mathbf{1 0 . 6 . 1 8} \mathbf{( p . 3 5 2 )} \mathbf{y}=e^{4 t}\left[\begin{array}{c}5 \cos 6 t+5 \sin 6 t \\ \cos 6 t-3 \sin 6 t\end{array}\right]$
10.6.19 (p. 352) $\mathbf{y}=e^{t}\left[\begin{array}{c}17 \cos 3 t-\sin 3 t \\ 7 \cos 3 t+3 \sin 3 t\end{array}\right] \mathbf{1 0 . 6 . 2 0}$ (p. 352) $\mathbf{y}=e^{t / 2}\left[\begin{array}{c}\cos (t / 2)+\sin (t / 2) \\ -\cos (t / 2)+2 \sin (t / 2)\end{array}\right]$
10.6.21 (p. 352) $\mathbf{y}=\left[\begin{array}{r}1 \\ -1 \\ 2\end{array}\right] e^{t}+e^{4 t}\left[\begin{array}{c}3 \cos t+\sin t \\ \cos t-3 \sin t \\ 4 \cos t-2 \sin t\end{array}\right]$
10.6.22 (p. 352) $\mathbf{y}=\left[\begin{array}{l}4 \\ 4 \\ 2\end{array}\right] e^{8 t}+e^{2 t}\left[\begin{array}{c}4 \cos 2 t+8 \sin 2 t \\ -6 \sin 2 t+2 \cos 2 t \\ 3 \cos 2 t+\sin 2 t\end{array}\right]$
10.6.23 (p. 352) $\mathbf{y}=\left[\begin{array}{l}0 \\ 3 \\ 3\end{array}\right] e^{-4 t}+e^{4 t}\left[\begin{array}{c}15 \cos 6 t+10 \sin 6 t \\ 14 \cos 6 t-8 \sin 6 t \\ 7 \cos 6 t-4 \sin 6 t\end{array}\right]$
$\mathbf{1 0 . 6 . 2 4}$ (p. 352) $\mathbf{y}=\left[\begin{array}{r}6 \\ -3 \\ 3\end{array}\right] e^{8 t}+\left[\begin{array}{c}10 \cos 4 t-4 \sin 4 t \\ 17 \cos 4 t-\sin 4 t \\ 3 \cos 4 t-7 \sin 4 t\end{array}\right]$
10.6.29 (p. 353) $\mathbf{U}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}-1 \\ 1\end{array}\right], \mathbf{V}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}1 \\ 1\end{array}\right]$
10.6.30 (p. 353) $\mathbf{U} \approx\left[\begin{array}{l}.5257 \\ .8507\end{array}\right], \mathbf{V} \approx\left[\begin{array}{r}-.8507 \\ .5257\end{array}\right]$
10.6.31 (p. 353) $\mathrm{U} \approx\left[\begin{array}{l}.8507 \\ .5257\end{array}\right]$,
$\mathbf{V} \approx\left[\begin{array}{r}-.5257 \\ .8507\end{array}\right] \quad \mathbf{1 0 . 6 . 3 2}(\mathbf{p . 3 5 3}) \mathbf{U} \approx\left[\begin{array}{r}-.9732 \\ .2298\end{array}\right], \mathbf{V} \approx\left[\begin{array}{l}.2298 \\ .9732\end{array}\right]$
10.6.33 (p. 353) $\mathbf{U} \approx\left[\begin{array}{c}.5257 \\ .8507\end{array}\right], \mathbf{V} \approx\left[\begin{array}{l}-.8507 \\ .5257\end{array}\right]$
10.6.34 (p. 353) $\mathbf{U} \approx\left[\begin{array}{r}-.5257 \\ .8507\end{array}\right], \mathbf{V} \approx\left[\begin{array}{l}.8507 \\ .5257\end{array}\right]$
10.6.35 (p. 353) $\mathbf{U} \approx\left[\begin{array}{r}-.8817 \\ .4719\end{array}\right], \mathbf{V} \approx\left[\begin{array}{l}.4719 \\ .8817\end{array}\right]$
10.6.36 (p. 353) $\mathbf{U} \approx\left[\begin{array}{l}.8817 \\ .4719\end{array}\right], \mathbf{V} \approx\left[\begin{array}{r}-.4719 \\ .8817\end{array}\right]$
10.6.37(p.353) $\mathrm{U}=\left[\begin{array}{l}0 \\ 1\end{array}\right], \mathrm{V}=\left[\begin{array}{l}-1 \\ 0\end{array}\right] \quad \quad 10.6 .38(\mathbf{p . 3 5 3}) \mathrm{U}=\left[\begin{array}{l}0 \\ 1\end{array}\right], \mathrm{V}=\left[\begin{array}{l}1 \\ 0\end{array}\right]$
$10.6 .39(\mathbf{p} .353) \mathrm{U}=\frac{1}{\sqrt{2}}\left[\begin{array}{l}1 \\ 1\end{array}\right], \mathrm{V}=\frac{1}{\sqrt{2}}\left[\begin{array}{r}-1 \\ 1\end{array}\right] \quad \mathbf{1 0 . 6 . 4 0}(\mathbf{p} .353) \mathbf{U} \approx\left[\begin{array}{l}.5257 \\ .8507\end{array}\right], \mathbf{V} \approx\left[\begin{array}{r}-.8507 \\ .5257\end{array}\right]$
Section 10.7 Answers, pp. 360-361
10.7.1 (p.360) $\left[\begin{array}{c}5 e^{4 t}+e^{-3 t}(2+8 t) \\ -e^{4 t}-e^{-3 t}(1-4 t)\end{array}\right] \mathbf{1 0 . 7 . 2}(\mathbf{p . 3 6 0})\left[\begin{array}{c}13 e^{3 t}+3 e^{-3 t} \\ -e^{3 t}-11 e^{-3 t}\end{array}\right] \mathbf{1 0 . 7 . 3}(\mathbf{p} .360) \frac{1}{9}\left[\begin{array}{c}7-6 t \\ -11+3 t\end{array}\right]$
$\mathbf{1 0 . 7 . 4}$ (p. 360) $\left[\begin{array}{c}5-3 e^{t} \\ -6+5 e^{t}\end{array}\right]$
$\mathbf{1 0 . 7 . 5}$ (p. 360) $\left[\begin{array}{c}e^{-5 t}(3+6 t)+e^{-3 t}(3-2 t) \\ -e^{-5 t}(3+2 t)-e^{-3 t}(1-2 t)\end{array}\right] \mathbf{1 0 . 7 . 6}$ (p.360) $\left[\begin{array}{c}t \\ 0\end{array}\right] \mathbf{1 0 . 7 . 7}(\mathbf{p . 3 6 0})-\frac{1}{6}\left[\begin{array}{c}2-6 t \\ 7+6 t \\ 1-12 t\end{array}\right]$
10.7.8 (p. 360) $-\frac{1}{6}\left[\begin{array}{c}3 e^{t}+4 \\ 6 e^{t}-4 \\ 10\end{array}\right]$
10.7.9(p.360) $\frac{1}{18}\left[\begin{array}{c}e^{t}(1+12 t)-e^{-5 t}(1+6 t) \\ -2 e^{t}(1-6 t)-e^{-5 t}(1-12 t) \\ e^{t}(1+12 t)-e^{-5 t}(1+6 t)\end{array}\right] \quad \mathbf{1 0 . 7 . 1 0}\left(\right.$ p. 360) $\frac{1}{3}\left[\begin{array}{c}2 e^{t} \\ e^{t} \\ 2 e^{t}\end{array}\right] \mathbf{1 0 . 7 . 1 1 ( \mathbf { p . 3 6 0 } )}\left[\begin{array}{c}t \sin t \\ 0\end{array}\right]$ 10.7.12 (p. 360) $-\left[\begin{array}{l}t^{2} \\ 2 t\end{array}\right]$
$10.7 .13(\mathbf{p . 3 6 0})(t-1)(\ln |t-1|+t)\left[\begin{array}{r}1 \\ -1\end{array}\right] \mathbf{1 0 . 7 . 1 4}(\mathbf{p . 3 6 0}) \frac{1}{9}\left[\begin{array}{c}5 e^{2 t}-e^{-3 t} \\ e^{3 t}-5 e^{-2 t}\end{array}\right] \mathbf{1 0 . 7 . 1 5 ( \mathbf { p } . 3 6 0 )} \frac{1}{4 t}\left[\begin{array}{c}2 t^{3} \ln |t|+t^{3}(t+2) \\ 2 \ln |t|+3 t-2\end{array}\right]$
10.7.16 (p. 360) $\frac{1}{2}\left[\begin{array}{c}t e^{-t}(t+2)+\left(t^{3}-2\right) \\ t e^{t}(t-2)+\left(t^{3}+2\right)\end{array}\right] \mathbf{1 0 . 7 . 1 7}$ (p. 360) $-\left[\begin{array}{c}t \\ t \\ t\end{array}\right] \mathbf{1 0 . 7 . 1 8}$ (p.360) $\frac{1}{4}\left[\begin{array}{c}-3 e^{t} \\ 1 \\ e^{-t}\end{array}\right]$
$\mathbf{1 0 . 7 . 1 9}$ (p. 360) $\left[\begin{array}{c}2 t^{2}+t \\ t \\ -t\end{array}\right] \mathbf{1 0 . 7 . 2 0}$ (p.360) $\frac{e^{t}}{4 t}\left[\begin{array}{l}2 t+1 \\ 2 t-1 \\ 2 t+1\end{array}\right]$
10.7.22 (p. 361) (a) $\mathbf{y}^{\prime}=\left[\begin{array}{cccc}0 & 1 & \cdots & 0 \\ 0 & 0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 1 \\ -P_{n}(t) / P_{0}(t) & -P_{n-1} / P_{0}(t) & \cdots & -P_{1}(t) / P_{0}(t)\end{array}\right] \mathbf{y}+\left[\begin{array}{c}0 \\ 0 \\ \vdots \\ F(t) / P_{0}(t)\end{array}\right]$.
(b) $\left[\begin{array}{cccc}y_{1} & y_{2} & \cdots & y_{n} \\ y_{1}^{\prime} & y_{2}^{\prime} & \cdots & y_{n}^{\prime} \\ \vdots & \vdots & \ddots & \vdots \\ y_{1}^{(n-1)} & y_{2}^{(n-1)} & \cdots & y_{n}^{(n-1)}\end{array}\right]$

## Index

## A

Abel's formula, 136-139, ??
Accelerated payment, ??
Acceleration due to gravity, 106
Airy's equation, 219
Amplitude,
of oscillation, 191
time-varying, 198
Amplitude-phase form, 191
Aphelion distance, ??
Apogee, ??
Applications,
of first order equations, ??-??
autonomous second order equations, 115-??
cooling problems, 96-97
curves, ??-??
elementary mechanics, 105-129
growth and decay, ??-??
mixing problems, 98-105
of linear second order equations, 188-??
motion under a central force, ??-??
motion under inverse square law force, ??-??
RLC circuit, 205-??
spring-mass systems, 188-205
Autonomous second order equations, 115-??
conversion to first order equations, 115
damped 125-130
pendulum 126
spring-mass system, 125
Newton's second law of motion and, 116
undamped 117-124
pendulum 125-121
spring-mass system, 117-125
stability and instability conditions for, 122-130

## B

Beat, 194
Bernoulli's equation, ??-??
Bessel functions of order $\nu$,??
Bessel's equation, 141 287, ??
of order $\nu, ? ?$
of order zero, ??
ordinary point of, 219
singular point of, 219,??
Bifurcation value, 47, 128
Birth rate, 2
Laplace equation,

Capacitance, ??
Capacitor, ??
Carbon dating, ??
Central force,
motion under a, ??-??
in terms of polar coordinates,
Characteristic equation, 146
with complex conjugate roots, 149-151
with disinct real roots, 146-151
with repeated real root, 147,151
Characteristic polynomial, 146, 237, ??
Charge, ??
steady state, ??
Chebyshev polynomials, 221
Chebshev's equation, 221
Circuit, RLC. See RLC circuit
Closed Circuit, 205
Coefficient(s) See also Constant coefficient equations
computing recursively, 221
in Frobenius solutions, ??-??
undetermined, method of, 162-178, ??-??
principle of superposition and, 166
Coefficient matrix, 308, 309
Competition, species, 6, 330
Complementary equation, 31 , ??
Complementary system, 354
Compound interest, continuous, ??, ??
Constant,
damping, 125
decay, ??
spring, 188
temperature decay, 96
Constant coefficient equations, 146, ??
homogeneous, 146-155
with complex conjugate roots, 149-151
with distinct real roots, 146, 151
higher order. See Higher order constant coef-
ficient homogeneous equations
with repeated real roots, 147,151
with impulses, 290-296
nonhomogeneous, 162-178
with piecewise continuous forcing functions, 272279
Constant coefficient homogeneous linear systems of differential equations, 320-353
geometric properties of solutions,
when $n=2$, 326-328, 339-342, 349-352
with complex eigenvalue of constant matrix, 344352
with defective constant matrix, 345-344
with linearly independent eigenvetors, 320-330
Constant solutions of separable first order equations, 43-46
Converge absolutely, 208
Convergence,
of improper integral, 240
open interval of, 208
radius of, 208
Convergent power series, 208
Convolution, 280-290
convolution integral, 284-288
defined, 280
theorem, 281
transfer functions, 285-287
Volterra integral equation, 284
Cooling, Newton's law of, 3, 96
Cooling problems, 96-97, 103-103
Cosine series, Fourier,
Critically damped motion, 199-199
oscillation, ??-??
Critical point, 116
Current, 205
steady state, ??
transient, ??
Curves, ??-??
equipotential, ??
geometric problems, ??
isothermal, ??
one-parameter famlies of, ??-?? subsubitem defined, ??
differential equation for, ??
orthogonal trajectories, ??-??, ??
finding, ??-??

## D

Damped autonomous second order equations, 124-??
for pendulum, 126
for spring-mass system, 125
Damped motion, 188
Damping,
$R L C$ circuit in forced oscllation with, ??
spring-mass systems with, 125, 188, 197-205
critically damped motion, 199-201
forced vibrations, 201-204
free vibrations, 197-201
overdamped motion, 198
underdamped motion, 198
spring-mass systems without, 189-196
forced oscillation, 193-196
Damping constant, 125
Damping forces, 116, 188
Dashpot, 188
Dating, carbon, ??-??
Death rate 2
Decay, See Exponential growth and decay,
Decay constant, ??
Derivatives, Laplace transform of, 257-258
Differential equations,
defined, 7
order of, 7
ordinary, 7
partial, 7
solutions of, 8-10

Differentiation of power series, 210
Dirac, Paul A. M., 290
Dirac delta function, 290
Direction fields for first order equations, 14-24
Discontinuity,
jump, 244
removable, 253
Distributions, theory of, 291
Divergence of improper integral, 240
Divergent power series, 208

## E

Eccentricity of orbit, ??
Elliptic orbit, ??
Epidemics 5-47
Equidimensional equation, ??
Equilibrium, 116
spring-mass system, 188
Equilibrium position, 188
Equipotentials, ??
Error(s),
in applying numerical methods, 74
in Euler's method, 75-79
at the $i$-th step, 74
truncation, 74
global, 79, 86, ??
local, 77
local, numerical methods with $O\left(h^{3}\right), 90-92$
Escape velocity, 112
Euler's equation, ??-??, 161
Euler's identity, 74-84
Euler's method, 74-84
error in, 75-77
truncation, 77-79
improved, 85-90
semilinear, 79-82
step size and accuracy of, 75
Exact first order equations, 55-63
implicit solutions of, 55-56
procedurs for solving, 58
Exactness condition, 56
Existence of solutions of nonlinear first order equations, 49-54
Existence theorem, 35, 49
Exponential growth and decay, ??-??
carbon dating, ??
interest compounded continuously, ??
mixed growth and decay, ??
radioactive decay, ??
savings program, ??
Exponential order, function of, 246
F
First order equations, 28-71
applications of See under Applications.
autonomous second order equation converted to, 115
direction fields for, 14-18
exact, 55-63
implicit solution of, 55
procedurs for solving, 58
linear, 27-39
homogeneous, 83-31
nonhomogeneous, 31-36
solutions of, 27
nonlinear, 36, 45, 49-??
existence and uniqueness of solutions of, 4954
transformation into separables, ??-??
numerical methods for solving. See Numerical method
separable, 39-48, ??-??
constant solutions of, 43-44
implicit solutions of, 42-43
First order systems of equations,
higher order systems written as, 304
scalar differential equations written as, 305
First shifting theorem, 243
Force(s)
damping, 116, 188
gravitational, 105, 112
impulsive, 291
lines of, ??
motion under central, ??-??
motion under inverse square law, ??-??
Forced motion, 188
oscillation
damped, ??-??
undamped, 193-196
vibrations, 201-204
Forcing function, 132
without exponential factors, 308,309 ??-??
with exponential factors, 175-174
piecewise continuous constant equations with, 272-279
Free fall under constant gravity, 11
Free motion, 188
oscillation, $R L C$ circuit in, ??-??
vibrations, 197-201
Frequency, 198
of simple harmonic motion, ??
Frobenius solutions, ??-??
indicial equation with distinct real roots differing by an integer, ??-??
indicial equation with distinct real roots not differing by an integer, ??-??
indicial equation with repeated root, ??-??
power series in, ??
recurrence relationship in, ??
two term, ??-??
verifying, ??
Fundamental matrix, 315
Fundamental set of solutions, of higher order constant coefficient homogeneous equations, ??-??
of homogeneous linear second order equations, 135, 139
of homogeneous linear systems of differential equations, 313, 315
of linear higher order equations, ??

## G

Gamma function, 249
Generalized Riccati equation, ??, ??
General solution
of higher order constant coefficient homogeneous equations, ??-??
of homogeneous linear second order equations, 135
of homogeneous linear systems of differential equations, 313,315
of linear higher order equations, ??, ??
of nonhomogeneous linear first order equations, 27, 35
of nonhomogeneous linear second order equations, 155, ??-??
Geometric problems, 166-??
Global truncation error in Euler's method, 79
Glucose absorption by the body, 5
Gravitation, Newton's law of, 105, 128, ??, 303, 307
Gravity, acceleration due to, 105
Grid, rectangular, 14
Growth and decay,
carbon dating, ??
exponential, ??-??
interest compounded continuously, ??-??
mixed growth and decay, ??
radioactive decay, ??
savings program, ??

## H

Half-life, ??
Half-line, 326
Half-plane, 340
Harmonic conjugate function, 63
Harmonic function, 63
Harmonic motion, simple, 118, 191 ??, ??
amplitude of oscillation, 191
natural frequancy of, 191
phase angle of, 191
nonhomogeneous problems,
Heat flow lines, ??
Heaviside's method, 252, 256
Hermite's equation, 221
Heun's method, 92
Higher order constant coefficient homogeneous equations, ??-??
characteristic polynomial of, ??-??
fundamental sets of solutions of, ??
general solution of, ??-??
Homogeneous linear first order equations, 27-31
general solutions of, 30
separation of variables in, 31
Homogeneous linear higher order equations, ??
Homogeneous linear second order equations, 132-155
constant coefficient, 146-155
with complex conjugate roots, 149-151
with distinct real roots, 146-151
with repeated real roots, 146-149, 151
solutions of, 132, 135
the Wronskian and Abel's formula, 136-139
Homogeneous linear systems of differential equations, 308
basic theory of, 313-319
constant coefficient, 320-353
with complex eigenvalues of coefficient matrix, 344-353
with defective coefficient matrix, 331-339
geometric properties of solutions when $n=$ $2,320-328,339-342,349-352$
with linearly independent eigenvectors, 320328 subitem fundamental set of solutions of, 313, 315
general solution of, 313, 315
trivial and nontrivial solution of, 313
Wronskian of solution set of, 315
Homogeneous nonlinear equations
defined, ??
transformation into separable equations, ??-??
Hooke's law, 188-189

## I

Imaginary part, 150
Implicit function theorem, 42
Implicit solution(s) 55-56
of exact first order equations, 55-56
of initial value problems, 42
of separable first order equations, 42-43
Impressed voltage, 47
Improper integral, 240
Improved Euler method, 84-88 ??-??
semilinear, 88-90
Impulse function, 290
Impulse response, 287, 292
Impulses, constant coefficient equations with, 290297
Independence, linear
of $n$ function, ??
of two functions, 136
of vector functions, 317
Indicial equation, ??, ??
with distinct real roots differing by an integer, ??-??
with distinct real roots not differing by an integer, ??-? ?
with repeated root, ??-??
Indicial polynomial, ??, ??
Inductance, ??
Initial conditions, 10
Initial value problems, 10-12
implicit solution of, 42
Laplace transforms to solve, 257-262
formula for, 282-283
second order equations, 259-262
Integral curves, 8-8, 257-24,
Integrals,
convolution, 284-284
improper, 240
Integrating factors, 63-71
finding, 64-71
Interest compounded continuously, ??-??
Interval of validity, 11
Inverse Laplace transforms, 250-257
defined, 250
linearity property of, 250
of rational functions, 251-257
Inverse square law force, motion under, ??-??
Irregular singular point, ??
Isothermal curves, ??

Jump discontinuity, 244

## K

Kepler's second law, ??
Kepler's third law, ??
Kirchoff's Law, ??

L
Laguerre's equation, ??
boundary conditions,
formal solutions of,
Laplace transforms, 240-297
computation of simple, 240-242
of constant coefficient equations
with impulses 290-297
with piecewise continuous forcing functions, 272-279
convolution, 280-290
convolution integral, 284
defined, 280
theorem, 281
transfer functions, 285-287
definition of, 240
existence of, 244
First shifting theorem, 243
inverse, 250
defined, 248
linearity property of, 250
of rational functons, 251-256
linearity of, 243
of piecewise continuous functions, 264-271
unit step function and, 263-271
Second shifting theorem, 267
to solve initial value problems, 257-262
derivatives, 257-258
formula for, 282-283
second order equations, 259
tables of, 242
Legendre's equation, 141, 219
ordinary points of, 219
singular points of, 219,??
Limit, 244
Limit cycle, 128
Linear combination(s), 135, ??, 313
of power series, 214-216
Linear difference equations, second order homogeneous, 237
Linear first order equations, 27-39 homogeneous, 27-31
general solution of, 30
separation of variables, 31
nonhomogeneous, 27, 31-36
general solution of, 32-36
solutions in integral form, 34-35
variation of parameters to solve, 31,34
solutions of, 27-28
Linear higher order equations, ??-??
fundamental set of solutions of, ??, ??
general solution of, ??, ??
higher order constant coefficient homogeneous
equations, ??-?? characteristic polyomial of ??-??
fundamental sets of solutions of, ??-??
general solution of, ??-??
homogeneous, ??
nonhomogeneous, ??, ??
trivial and nontrivial solutions of, ??
undetermined coefficients for, ??-??
variation of parameters for, ??-??
derivation of method, ??-??
fourth order equations, ??-??
third order equations, ??
Wronskian of solutions of ??-??
Linear independence 136
of $n$ functions, ??
of two functions, 136
of vector functions, 313-315
Linearity,
of inverse Laplace transform, 250
of Laplace transform, 243
Linear second order equations, 132-186
applications of. See under Applications
defined, 132
homogeneous, 132-155
constant coefficient, 146-138
solutions of, 132-135
the Wronskian and Abel's formula, 136-139
nonhomnogeneous, 132, 155-186, ??, ??
comparison of methods for solving, 140
complementary equation for, 155
constant coefficient, 162-??
general solution of, 155-159
particular solution of, 155, 159-160
reduction of order to find general solution of, ??-??
superposition principle and, 159-160
undetermined coefficients method for, 162 ??
variation of parameters to find particular solution of, 178-186
series solutions of, 208-??
Euler's equation, ??-??
Frobenius solutions, ??-??
near an ordinary point, 219-236
with regular singular points, ??-??
Linear systems of differential equations, 308-361
defined, 308
homogeneous, 308
basic theory of, 313-320
constant coefficient, 320-353
fundamental set of solutions of, 313-315
general solution of, 313, 315
linear indeopendence of, 313, 315
trivial and nontrivial solution of, 313
Wronskian of solution set of, 315
nonhomogeneous, 308
variation of parameters for, 354-361
solutions to initial value problem, 308-310
Lines of force, ??
local truncation error, 77-79
numerical methods with $O\left(h^{3}\right)$, 90-92
Logistic equation 3
M
Maclaurin series, 209

Magnitude of acceleration due to gravity at Earth's surface, 106
Malthusian model, 2
Mathematical models, 2
validity of, ??, 96, 104
Matrix/matrices, 308-310
coefficient,
complex eigenvalue of, 344-353
defective, 331
fundamental, 315
Mechanics, elementary, 105-130
escape velocity, 112-113 115, 115
motion through resisting medium under constant gravitational force, 106-110
Newton's second law of motion, 105-106
pendulum motion
damped, 126-128
undamped, 125-121
spring-mass system
damped, 125-126, 189, 197-205
undamped, 117-125,188
units used in, 105
Midpoint method, 85
Mixed Fourier cosine series,
Mixed Fourier sine series,
Mixed growth and decay, ??
Mixing problems, 98-103
Models, mathematical, 2-2
validity of, ??, 96, 104
Motion,
damped, 188
critically, 199
overdamped, 198-199
underdamped, 197
elementary, See Mechanics, elementary
equation of, 189
forced, 189
free, 189
Newton's second law of, 6, 105-106, 116, 118, 125, 125-126, 188, ??, 304
autonomous second order equations and, 116
simple harmonic, 118, 189-192
amplitude of oscillation, 191
frequency of, 191
phase angle of, 191
through resisting medium under constant gravitational force, 106-110
under a central force, ??-??
under inverse square law force, ??-??
undamped, 188
Multiplicity, ??

## N

Natural frequency, 191
Natural length of spring, 188
Negative half plane, 340
Newton's law of cooling, 3, 96-96, 103-105
Newton's law of gravitation, 105, 128, ??, 303, 311
Newton's second law of motion, 105-106, 116, 118, 125, 128, 188, ??, 302, 303
autonomous second order equations and, 116
Nonhomogeneous linear second order equations, 27, 31, 36
general solution of, 32-34 35-36
solutions in integral form, 34
variation of parameters to solve, 31, 34
Nonhomogeneous linear second order equations, 132, 155-186
comparison of methods for solving, 184
complementary equation for, 155,155
constant coefficient, 162-??
general solution of, 155-159
particular solution of, 155, 155-159, 162-166, 178-184
reduction of order to find general solution of, ??-??
superposition principle and, 159-156
undetermined coefficients method for, 162-??
forcing functions with exponential factors, 172174
forcing functions without exponential factors, 169-172
superposition principle and, 166
variation of parameters to find particular solution of, 178-186
Nonhomogeneous linear systems of differential equations, 308
variation of parameters for, 354-361
Nonlinear first order equations, 45 49-??
existence and uniqueness of solutions of, 49-??
transformation into separable equations, ??-??
Nonoscillatory solution, ??
Nontrivial solutions
of homogeneous linear first order equations, 27
of homogeneous linear higher order equations, ??
of homogeneous linear second order equations, 132
of homogeneous linear systems of differential equations, 313
Numerical methods, 74-??, 307
with $O\left(h^{3}\right)$ local truncation, 90-92
error in, 74
Euler's method, 74-84
error in, 75-79
semilinear, 79-82
step size and accuracy of, 75
truncation error in, 76-79
Heun's method, 91
semilinear, 82
improved Euler method, 82, 85-88
semilinear, 88
midpoint, 92
Runge-Kutta method, 76, 82 ??-??, 306-307
for cases where $x_{0}$ isn't the left endpoint, ????
semilinear, 82, ??
for systems of differential equations, 307
Numerical quadrature, 94, ??

0
One-parameter families of curves, ??-??
defined, ??
differential equation for, ??
One-parameter families of functions, 27
Open interval of convergence, 208

Open rectangle, 49
Orbit, ??
eccentricity of, ??
elliptic, ??
period of, ??
Order of differential equation, 7
Ordinary differential equation, defined, 7
Ordinary point, series solutions of linear second order equations near, 219-238
Orthogonal trajectories, ??-??, finding, ??
Orthogonal with respect to a weighting function, 229, 229
Oscillation
amplitude of, 191
critically damped, ??
overdamped, ??
$R L C$ circuit in forced, with damping, ??-??
RLC circuit in free, ??-??
undamped forced, 193-196
underdamped, ??
Oscillatory solutions, 164-128, ??
Overdamped motion, 197-198

## P

Partial differential equations defined, 7
Partial fraction expansions, software packages to find, 256
Particular solutions of nonhomogeneous higher equations, ??, ??-??
Particular solutions of nonhomogeneous linear second order equations, 155, 159-159, 162166, 178-183
Particular solutions of nonhomogeneous linear systems equations, 354-361
Pendulum
damped, 126-128
undamped, 125-121
Perigee, ??
Perihelion distance, ??
Periodic functions, 249
Period of orbit, ??
Phase angle of simple harmonic motion, 191-191
Phase plane equivalent, 116
Piecewise continuous functions, 246
forcing, constant coeffocient equations with, 272279
Laplace transforms of 244-247, 264-271
unit step functions and, 263-271
Plucked string, wave equation applied to,
Poinccaré, Henri, 116
Polar coordinates
central force in terms of, ??-??
in amplitude-phase form, 191
Polynomial(s)
characteristic, 146, 237, ??
of higher order constant coefficient homogeneous equations, ??-??
Chebyshev, 221
indicial, ??, ??
Taylor, 211
trigonometric, ??
Polynomial operator, ??
Population growth and decay, 2
Positive half-plane, 339
Power series, 208-218
convergent, 208-209
defined, 208
differentiation of, 210-210
divergent, 208
linear combinations of, 214-216
radius of convergence of, 208, 208
shifting summation index in, 211-213
solutions of linear second order equations, represented by, 219-238
Taylor polynomials, 210
Taylor series, 209
uniqueness of 210-211

## Q

Quasi-period, 198

## R

Radioactive decay, ??-??
Radius of convergence of power series, 208, 208
Rational functions, inverse Laplace transforms of, 251257
Rayleigh, Lord, 122
Rayleigh's equation, 129
Real part, 150
Rectangle, open, 49
Rectangular grid, 14
Recurrence relations, 221
in Frobenius solutions, ??
two term, ??-??
Reduction of order, 148, ??-??
Regular singular points, ??-??
at $x_{0}=0, ? ?-? ?$
Removable discontinuity, 244
Resistance, ??
Resistor, ??
Resonance, 196
Ricatti, Jacopo Francesco, ??
Ricatti equation, ??
RLC circuit, 205-??
closed, 205
in forced oscillation with danping, ??
in free oscillation, ??-??
Roundoff errors, 74
Runge-Kutta method, 76, ??-??, 307
for cases where $x_{0}$ isn't the left endpoint, ??
for linear systems of differential equations, 307
semilinear, 82 , ??

## S

Savings program, growth of, ??
Scalar differential equations, 305
Second order differential equation, 6
autonomous, 115-130
conversion to first order equation, 115
damped, 125-130
Newton's second law of mation and, 116
undamped, 117-121
Laplace transform to solve, 259-261
linear, See linear second equations
two-point boundary value problems for,
Second order homogeneous linear difference equa-
tion, 237
Second shifting Theorem, 267-269
Semilinear Euler method, 79
Semilinear improved Euler method, 82, 88
Semilinear Runge-Kutta method, 84, ??
Separable first order equations, 39-48
constant solutions of, 43-44
implicit solutions, 42
transfomations of nonlinear equations to, ??-??
Bernoulli's equation, ??-??
homogeneous nonlinear equations, ??-??
other equations, ??
Separation of variables, 31, 39
to solve Laplace's equation,
Separatrix, 123, 122
Series, power. See Power series
Series solution of linear second order equations, 208??
Frobenius solutions, ??-??
near an ordinary point, 219
Shadow trajectory, 350-352
Shifting theorem
first, 243
second, 267-269
Simple harmonic motion, 189-193
amplitude of oscillation, 191
natural frequency of, 192
phase angle of, 192
Simpson's rule, ??
Singular point, 219
irregular, ??
regular, ??-??
Solution(s), 8-9 See also Frobenius solutions Nontrivial solutions Series solutions of linear second order equations Trivial solution
nonoscillatory, ??
oscillatory, ??
Solution curve, 8-8
Species, interacting, 6, 329
Spring, natural length of, 188, 189
Spring constant, 188
Spring-mass systems, 188-205
damped, 124, 189, 204-205
critically damped motion, 204-201
forced vibrations, 201-204
free vibrations, 201-202
overdamped motion, 197
underdamped motion, 197
in equilibrium, 188
simple harmonic motion, 189-193
amplitude of oscillation, 191
natural frequency of, 191
phase angle of, 191
undamped, 117-118, 189-196
forced oscillation, 193-204
Stability of equilibrium and critical point, 116-117
Steady state, ??
Steady state charge, ??
Steady state component, 203, 286
Steady state current, ??

String motion, wave equation applied to,
Summation index in power series, 212-213
Superposition, principle of, $38,159,166$, ??
method of undetermine coefficients and, 166
Systems of differential equations, 301-310 See also Linear systems of differential equations
first order
higher order systems rewritten as, 221-305
scalar differential equations rewritten as, 305
numerical solutions of, 307
two first order equations in two unknowns, 301303

## T

Tangent lines, ??
Taylor polynomials, 210
Taylor Series, 209
Temperature, Newton's law of cooling, 396-97, 103104
Temperature decay constant of the medium, 96
Terminal velocity, 106
Time-varying amplitude, 198
Total impulse, 290
Trajectory(ies),
of autonomous second order equations, 116
orthogonal, ??-??
finding, ??-174
shadow, 350
of $2 \times 2$ systems, 326-328, 339-342, 349-352
Transfer functions, 285
Transformation of nonlinear equations to separable first order, equations, ??-61
Bernoulli's equation, ??
homogeneous nonlinear equations, ??-??
other equations, ??
Transform pair, 240
Transient current, ??
Transient components, 203, 286
Transient solutions, ??
Trapezoid rule, 94
Trivial solution,
of homogeneous linear first order equations, 27
of homogeneous linear second order equations, 132
of homogeneous linear systems of differential equations, 313
of linear higher order differential equations, ??
Truncation error(s), 74
in Euler's method, 77
global, 79, 85
local, 77
numerical methods with $O\left(h^{3}\right), 90-92$
Two-point boundary value problems,

## U

Undamped autonomous second order equations, 117123
pendulum, 125-121
spring-mass system, 117-118
stability and instabilty conditions for, 122-123
Undamped motion, 188
Underdamped motion, 197
Underdamped oscillation, ??

Undetermined coefficients
for linear higher order equations, ??-?? forcing functions, ??-??
for linear second order equations, 162-178
principle of superposition, 166
Uniqueness of solutions of nonlinear first equations, 49-54
Uniqueness theorem, 35, 49, 132, ??, 309
Unit step function, 264-271
V
Validity, interval of, 11
Vandermonde, ??
Vandermonde determinant, ??
van der Pol's equation, 128
Variables, separation of, 31, 39
Variation of parameters
for linear first order equations, 31
for linear higher order equations, ??-??
derivation of method, ??-??
fourth order equations, ??-??
third order equations, ??
for linear higher second order equations, 178
for nonhomogeneous linear systems of differential equations, 354-361
Velocity
escape, 112-105
terminal, 106-109
Verhulst, Pierre, 3
Verhulst model, 3, 24, ??
Vibrations
forced, 201-204
free, 197-201
Voltage, impressed, 205
Voltage drop, ??
Volterra, Vito 284
Volterra integral equation, 284
W
plucked string,
assumptions, ??
Wave, traveling, ??
229
Wronskian
of solutions of homogeneous linear systems of differential equations, 315
of solutions of homogeneous second differential equations, 136-138
of solutions of homogeneous linear higher order differential equations, ??-??


[^0]:    IN THIS CHAPTER we study applications of linear second order equations.
    SECTIONS 6.1 AND 6.2 is about spring-mass systems.
    SECTION 6.2 is about $R L C$ circuits, the electrical analogs of spring-mass systems.
    SECTION 6.3 is about motion of an object under a central force, which is particularly relevant in the space age, since, for example, a satellite moving in orbit subject only to Earth's gravity is experiencing motion under a central force.

