# Spectral Theory of Partial Differential Equations 

## Lecture Notes

University of Illinois at Urbana-Champaign

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## Preface

A textbook presents far more material than any professor can cover in class. These lecture notes present only somewhat more than I covered during the half-semester course Spectral Theory of Partial Differential Equations (Math 595 STP) at the University of Illinois, Urbana-Champaign, in Fall 2011.

I make no claims of originality for the material presented, other than some originality of emphasis - I emphasize computable examples before developing the general theory. This approach leads to occasional redundancy, and sometimes we use ideas before they are properly defined, but I think students gain a better understanding of the purpose of a theory after they are first well grounded in specific examples.

Please email me with corrections, and suggested improvements.

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## Prerequisites and notation

We assume familiarity with elementary Hilbert space theory: inner product, norm, Cauchy-Schwarz, orthogonal complement, Riesz Representation Theorem, orthonormal basis (ONB), bounded operators, and compact operators.

All functions are assumed to be measurable. We use the function spaces

$$
\begin{aligned}
\mathrm{L}^{1} & =\text { integrable functions } \\
\mathrm{L}^{2} & =\text { square integrable functions } \\
\mathrm{L}^{\infty} & =\text { bounded functions }
\end{aligned}
$$

but we have no need of general $L^{p}$ spaces.
Sometimes we employ the $\mathrm{L}^{2}$-theory of the Fourier transform,

$$
\widehat{\mathrm{f}}(\xi)=\int_{\mathbb{R}^{\mathrm{d}}} f(x) e^{-2 \pi \mathrm{i} \xi \cdot x} d x
$$

Only the basic facts are needed, such as that the Fourier transform preserves the $L^{2}$ norm and maps derivatives in the spatial domain to multipliers in the frequency domain.

We use the language of Sobolev spaces throughout. Readers unfamiliar with this language can proceed unharmed: we mainly need only that

$$
\begin{aligned}
& \mathrm{H}^{1}=\mathrm{W}^{1,2}=\left\{\mathrm{L}^{2} \text {-functions with } 1 \text { derivative in } \mathrm{L}^{2}\right\} \\
& \mathrm{H}^{2}=\mathrm{W}^{2,2}=\left\{\mathrm{L}^{2} \text {-functions with } 2 \text { derivatives in } \mathrm{L}^{2}\right\}
\end{aligned}
$$

and

$$
\mathrm{H}_{0}^{1}=\mathrm{W}_{0}^{1,2}=\left\{\mathrm{H}^{1} \text {-functions that equal zero on the boundary }\right\} \text {. }
$$

(These characterizations are not mathematically precise, but they are good enough for our purposes.) Later we will recall the standard inner products that make these spaces into Hilbert spaces.

For more on Sobolev space theory, and related concepts of weak solutions and elliptic regularity, see [Evans].

## Introduction

Spectral methods permeate the theory of partial differential equations. One solves linear PDEs by separation of variables, getting eigenvalues when the spectrum is discrete and continuous spectrum when it is not. Linearized stability of a steady state or traveling wave of a nonlinear PDE depends on the sign of the first eigenvalue, or on the location of the continuous spectrum in the complex plane.

This minicourse aims at highlights of spectral theory for selfadjoint partial differential operators, with a heavy emphasis on problems with discrete spectrum.

Style of the course. Research work differs from standard course work. Research often starts with questions motivated by analogy, or by trying to generalize special cases. Normally we find answers in a nonlinear fashion, slowly developing a coherent theory by linking up and extending our scraps of known information. We cannot predict what we will need to know in order to succeed, and we certainly do not have enough time to study all relevant background material. To succeed in research, we must develop a rough mental map of the surrounding mathematical landscape, so that we know the key concepts and canonical examples (without necessarily knowing the proofs). Then when we need to learn more about a topic, we know where to begin.

This course aims to develop your mental map of spectral theory in partial differential equations. We will emphasize computable examples, and will be neither complete in our coverage nor completely rigorous in our approach. Yet you will finish the course having a much better appreciation of the main issues and techniques in the subject.

Closing thoughts. If the course were longer, then we could treat topics such as nodal patterns, geometric bounds for the first eigenvalue and the spectral gap, majorization techniques (passing from eigenvalue sums to spectral zeta functions and heat traces), and inverse spectral problems. And we could investigate more deeply the spectral and scattering theory of operators with continuous spectrum, giving applications to stability of traveling waves and similarity solutions. These fascinating topics must await another course. . .

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## Part I

## Discrete Spectrum

## Chapter 1

## ODE preview

## Goal

To review the role of eigenvalues and eigenvectors in solving 1st and 2nd order systems of linear ODEs; to interpret eigenvalues as decay rates, frequencies, and stability indices; and to observe formal analogies with PDEs.

## Notational convention

Eigenvalues are written with multiplicity, and are listed in increasing order (when real-valued):

$$
\lambda_{1} \leq \lambda_{2} \leq \lambda_{3} \leq \cdots
$$

## Spectrum of a real symmetric matrix

If $A$ is a real symmetric $d \times d$ matrix (e.g. $A=\left[\begin{array}{ll}a & b \\ b & c\end{array}\right]$ when $d=2$ ) or Hermitian matrix then its spectrum is the collection of eigenvalues:

$$
\operatorname{spec}(A)=\left\{\lambda_{1}, \ldots, \lambda_{d}\right\} \subset \mathbb{R}
$$

(see the figure). Recall that

$$
A v_{j}=\lambda_{j} v_{j}
$$

where the eigenvectors $\left\{v_{1}, \ldots, v_{d}\right\}$ can be chosen to form an ONB for $\mathbb{R}^{d}$.


Observe $A: \mathbb{R}^{\mathrm{d}} \rightarrow \mathbb{R}^{\mathrm{d}}$ is diagonal with respect to the eigenbasis:

$$
\begin{array}{r}
A\left(\sum c_{j} v_{j}\right)=\sum \lambda_{j} c_{j} v_{j} \\
{\left[\begin{array}{lll}
\lambda_{1} & & 0 \\
& \ddots & \\
0 & & \lambda_{d}
\end{array}\right]\left[\begin{array}{c}
c_{1} \\
\vdots \\
c_{d}
\end{array}\right]=\left[\begin{array}{c}
\lambda_{1} c_{1} \\
\vdots \\
\lambda_{d} c_{d}
\end{array}\right]}
\end{array}
$$

## What does the spectrum tell us about linear ODEs?

Example 1.1 (1st order). The equation

$$
\begin{aligned}
\frac{d v}{d t} & =-A v \\
v(0) & =\sum c_{j} v_{j}
\end{aligned}
$$

has solution

$$
v(\mathrm{t})=\mathrm{e}^{-\lambda \mathrm{t}} v(0) \stackrel{\text { def }}{=} \sum e^{-\lambda_{j} t} c_{j} v_{j} .
$$

Notice $\lambda_{j}=$ decay rate of the solution in direction $v_{j}$ if $\lambda_{j}>0$, or growth rate (if $\lambda_{j}<0$ ).

Long-time behavior: the solution is dominated by the first mode, with

$$
v(t) \sim e^{-\lambda_{1} t} c_{1} v_{1} \quad \text { for large } t
$$

assuming $\lambda_{1}<\lambda_{2}$ (so that the second mode decays faster than the first). The rate of collapse onto the first mode is governed by the spectral gap $\lambda_{2}-\lambda_{1}$
since

$$
\begin{aligned}
v(t) & =e^{-\lambda_{1} t}\left(c_{1} v_{1}+\sum_{j=2}^{d} e^{-\left(\lambda_{j}-\lambda_{1}\right) t} c_{j} v_{j}\right) \\
& \sim e^{-\lambda_{1} t}\left(c_{1} v_{1}+O\left(e^{-\left(\lambda_{2}-\lambda_{1}\right) t}\right)\right.
\end{aligned}
$$

Example 1.2 (2nd order). Assume $\lambda_{1}>0$, so that all the eigenvalues are positive. Then

$$
\begin{aligned}
\frac{\mathrm{d}^{2} v}{\mathrm{dt}^{2}} & =-\mathrm{A} v \\
v(0) & =\sum c_{j} v_{j} \\
v^{\prime}(0) & =\sum c_{j}^{\prime} v_{j}
\end{aligned}
$$

has solution

$$
\begin{aligned}
v(t) & =\cos (\sqrt{A} t) v(0)+\frac{1}{\sqrt{A}} \sin (\sqrt{A} t) v^{\prime}(0) \\
& \xlongequal{\text { def }} \sum \cos \left(\sqrt{\lambda_{j}} t\right) c_{j} v_{j}+\sum \frac{1}{\sqrt{\lambda_{j}}} \sin \left(\sqrt{\lambda_{j}} t\right) c_{j}^{\prime} v_{j} .
\end{aligned}
$$

Notice $\sqrt{\lambda_{j}}=$ frequency of the solution in direction $v_{j}$.
Example 1.3 (1st order imaginary). The equation

$$
\begin{aligned}
i \frac{d v}{d t} & =A v \\
v(0) & =\sum c_{j} v_{j}
\end{aligned}
$$

has complex-valued solution

$$
v(t)=e^{-i A t} v(0) \stackrel{\text { def }}{=} \sum e^{-i \lambda_{j} t} c_{j} v_{j} .
$$

This time $\lambda_{j}=$ frequency of the solution in direction $v_{j}$.

## What does the spectrum tell us about nonlinear ODEs? In/stability!

Example 1.4 (1st order nonlinear). Suppose

$$
\frac{\mathrm{d} v}{\mathrm{dt}}=\mathrm{F}(v)
$$

where the vector field $F$ satisfies $F(0)=0$, with first order Taylor expansion

$$
\mathrm{F}(v)=\mathrm{B} v+\mathrm{O}\left(|v|^{2}\right)
$$

for some matrix B having d linearly independent eigenvectors $\nu_{1}, \ldots, v_{\mathrm{d}}$ and corresponding eigenvalues $\lambda_{1}, \ldots, \lambda_{d} \in \mathbb{C}$. (The eigenvalues come in complex conjugate pairs, since $B$ is real.)

Clearly $v(\mathrm{t}) \equiv 0$ is an equilibrium solution. Is it stable? To investigate, we linearize the ODE around the equilibrium to get $\frac{\mathrm{d} v}{\mathrm{dt}}=\mathrm{B} v$, which has solution

$$
v(t)=e^{B t} v(0)=\sum e^{\lambda_{j} t} c_{j} v_{j} .
$$

Notice $v(\mathrm{t}) \rightarrow 0$ as $\mathrm{t} \rightarrow \infty$ if $\operatorname{Re}\left(\lambda_{\mathrm{j}}\right)<0$ for all $\mathfrak{j}$, whereas $|v(\mathrm{t})| \rightarrow \infty$ if $\operatorname{Re}\left(\lambda_{j}\right)>0$ for some $\mathfrak{j}$ (provided the corresponding coefficient $c_{j}$ is nonzero, and so on). Hence the equilibrium solution $v(t) \equiv 0$ is:

- linearly asymptotically stable if $\operatorname{spec}(\mathrm{B}) \subset$ LHP,

- linearly unstable if $\operatorname{spec}(B) \cap \operatorname{RHP} \neq \emptyset$.

The Linearization Theorem guarantees that the nonlinear ODE indeed behaves like the linearized ODE near the equilibrium solution, in the stable and unstable cases.


The nonlinear ODE's behavior requires further investigation in the neutrally stable case where the spectrum lies in the closed left half plane and intersects the imaginary axis $\left(\operatorname{Re}\left(\lambda_{j}\right) \leq 0\right.$ for all $j$ and $\operatorname{Re}\left(\lambda_{j}\right)=0$ for some j).


For example, if $B=\left[\begin{array}{cc}0 & -1 \\ 1 & 0\end{array}\right]$ (which has eigenvalues $\pm i$ ), then the phase portrait for $\frac{d v}{d t}=\mathrm{B} v$ consists of circles centered at the origin, but the phase portrait for the nonlinear system $\frac{\mathrm{d} v}{\mathrm{dt}}=\mathrm{F}(v)$ might spiral in towards the origin (stability) or out towards infinity (instability), or could display even more complicated behavior.

## Looking ahead to PDEs

Now suppose $A$ is an elliptic operator on a domain $\Omega \subset \mathbb{R}^{d}$. For simplicity, take $A=-\Delta$. Assume boundary conditions that make the operator self-adjoint (we will say more about boundary conditions later). Then the eigenvalues $\lambda_{j}$ and eigenfunctions $v_{j}(x)$ of the Laplacian satisfy

$$
-\Delta v_{j}=\lambda_{j} v_{j} \quad \text { in } \Omega
$$

and the spectrum increases to infinity:

$$
\lambda_{1} \leq \lambda_{2} \leq \lambda_{3} \leq \cdots \rightarrow \infty
$$

The eigenfunctions form an ONB for $\mathrm{L}^{2}(\Omega)$.
Substituting $A=-\Delta$ into the ODE Examples 1.1-1.3 transforms them into famous partial differential equations for the function $v(x, t)$. We solve these PDEs formally by separation of variables:

- Example 1.1 - diffusion equation $v_{\mathrm{t}}=\Delta v$. Separation of variables gives the solution

$$
v=e^{\Delta t} v(\cdot, 0) \stackrel{\text { def }}{=} \sum e^{-\lambda_{j} t} c_{j} v_{j}
$$

where the initial value is $v(\cdot, 0)=\sum c_{j} v_{j}$. Here $\lambda_{j}=$ decay rate.

- Example 1.2 - wave equation $\nu_{\mathrm{tt}}=\Delta \nu$. Separation of variables gives

$$
\begin{aligned}
v & =\cos (\sqrt{-\Delta} t) v(\cdot, 0)+\frac{1}{\sqrt{-\Delta}} \sin (\sqrt{-\Delta} t) v_{t}(\cdot, 0) \\
& \stackrel{\text { def }}{=} \sum \cos \left(\sqrt{\lambda_{j}} t\right) c_{j} v_{j}+\sum \frac{1}{\sqrt{\lambda_{j}}} \sin \left(\sqrt{\lambda_{j}} t\right) c_{j}^{\prime} v_{j}
\end{aligned}
$$

So $\sqrt{\lambda_{j}}=$ frequency and $v_{j}=$ mode of vibration.

- Example 1.3 - Schrödinger equation $\mathfrak{i} v_{\mathrm{t}}=-\Delta v$. Separation of variables gives

$$
v=e^{\mathrm{i} \Delta t} v(\cdot, 0) \stackrel{\text { def }}{=} \sum e^{-\mathrm{i} \lambda_{j} \mathrm{t}} \mathrm{c}_{\mathrm{j}} v_{j} .
$$

Here $\lambda_{j}=$ frequency or energy level, and $v_{j}=$ quantum state.
We aim in what follows to analyze not just the Laplacian, but a whole family of related operators including:

$$
\begin{array}{ll}
A=-\Delta & \text { Laplacian } \\
A=-\Delta+V(x) & \text { Schrödinger operator, } \\
A=(i \nabla+\vec{A})^{2} & \text { magnetic Laplacian } \\
A=(-\Delta)^{2}=\Delta \Delta & \text { biLaplace operator. }
\end{array}
$$

The spectral theory of these operators will help us understand the stability of the different kinds of "equilibria" for evolution equations: steady states, standing waves, traveling waves, and similarity solutions.

## Chapter 2

## Laplacian - computable spectra

## Goal

To develop a library of explicitly computable spectra, which we use later to motivate and understand the general theory.

References [Strauss] Chapters 4, 10; [Farlow] Lesson 30.

## Notation

Let $\Omega$ be a bounded domain in $\mathbb{R}^{d}, d \geq 1$. Fix $L>0$.
Abbreviate "boundary condition" as "BC":

- Dirichlet BC means $u=0$ on $\partial \Omega$,
- Robin BC means $\frac{\partial u}{\partial n}+\sigma u=0$ on $\partial \Omega$ (where $\sigma \in \mathbb{R}$ is the Robin constant),
- Neumann BC means $\frac{\partial u}{\partial n}=0$ on $\partial \Omega$.


## Spectra of the Laplacian

$$
\Delta=\nabla \cdot \nabla=\left(\frac{\partial}{\partial x_{1}}\right)^{2}+\cdots+\left(\frac{\partial}{\partial x_{d}}\right)^{2}
$$

Eigenfunctions satisfy $-\Delta \mathfrak{u}=\lambda \mathfrak{u}$, and we order the eigenvalues in increasing order as

$$
\lambda_{1} \leq \lambda_{2} \leq \lambda_{3} \leq \cdots \rightarrow \infty
$$

To get an ONB one should normalize the eigenfunctions in $L^{2}$, but we will not normalize the following examples.

One dimension $-\mathfrak{u}^{\prime \prime}=\lambda u$

1. Circle $\mathbb{T}=\mathbb{R} / 2 \pi \mathbb{Z}$, periodic BC: $\mathfrak{u}(-\pi)=\mathfrak{u}(\pi), \mathfrak{u}^{\prime}(-\pi)=\mathfrak{u}^{\prime}(\pi)$.

Eigenfunctions $e^{i j x}$ for $\mathfrak{j} \in \mathbb{Z}$, or equivalently $1, \cos (j x), \sin (j x)$ for $j \geq 1$.
Eigenvalues $\lambda_{j}=\mathfrak{j}^{2}$ for $\mathfrak{j} \in \mathbb{Z}$, or $\lambda=0^{2}, 1^{2}, 1^{2}, 2^{2}, 2^{2}, \ldots$
2. Interval ( $0, \mathrm{~L}$ )
(a) Dirichlet BC: $u(0)=u(L)=0$.

Eigenfunctions $u_{j}(x)=\sin (j \pi x / L)$ for $\mathfrak{j} \geq 1$.
Eigenvalues $\lambda_{j}=(j \pi / L)^{2}$ for $\mathfrak{j} \geq 1$, e.g. $L=\pi \Rightarrow \lambda=1^{2}, 2^{2}, 3^{2}, \ldots$

(b) Robin BC: $-u^{\prime}(0)+\sigma u(0)=u^{\prime}(L)+\sigma u(L)=0$.

Eigenfunctions $u_{j}(x)=\sqrt{\rho_{j}} \cos \left(\sqrt{\rho_{j}} x\right)+\sigma \sin \left(\sqrt{\rho_{j}} x\right)$.
Eigenvalues $\rho_{j}=j$ th positive root of $\tan (\sqrt{\rho} L)=\frac{2 \sigma \sqrt{\rho}}{\rho-\sigma^{2}}$ for $\mathfrak{j} \geq 1$.
(c) Neumann BC: $u^{\prime}(0)=u^{\prime}(L)=0$.

Eigenfunctions $\mathfrak{u}_{j}(x)=\cos (j \pi x / L)$ for $\mathfrak{j} \geq 0\left(\right.$ note $\left.u_{0} \equiv 1\right)$.
Eigenvalues $\mu_{j}=(j \pi / L)^{2}$ for $\mathfrak{j} \geq 0$, e.g. $L=\pi \Rightarrow \lambda=0^{2}, 1^{2}, 2^{2}, 3^{2}, \ldots$


## Spectral features in 1 dim

i. Scaling: eigenvalue must balance $\mathrm{d}^{2} / \mathrm{d} x^{2}$, and so $\lambda \sim(\text { length scale })^{-2}$.

Precisely, $\lambda_{j}([0, t L])=\lambda_{j}([0, L]) / t^{2}$.
ii. Asymptotic: eigenvalues grow at a regular rate, $\lambda_{j} \sim($ const. $) j^{2}$
iii. Robin spectrum lies between Neumann and Dirichlet:
as one sees formally by letting $\sigma$ approach 0 or $\infty$ in the Robin $\mathrm{BC} \frac{\partial u}{\partial n}+\sigma u=$ 0 .

Two dimensions $-\Delta u=\lambda u$

1. Rectangle $\Omega=(0, \mathrm{~L}) \times(0, \mathrm{M})$ (product of intervals).

Separate variables using rectangular coordinates $x_{1}, x_{2}$. See the figures at the end of the chapter!
(Note that every rectangle can be reduced to a rectangle with sides parallel to the coordinate axes because the Laplacian, and hence its spectrum, is rotationally and translationally invariant.)
(a) Dirichlet BC: $\mathfrak{u}=0$

Eigenfunctions $\mathfrak{u}_{\mathfrak{j k}}(x)=\sin \left(j \pi x_{1} / L\right) \sin \left(k \pi x_{2} / M\right)$ for $\mathfrak{j}, k \geq 1$.

Eigenvalues $\lambda_{j k}=(j \pi / L)^{2}+(k \pi / M)^{2}$ for $\mathfrak{j}, k \geq 1$,
e.g. $L=M=\pi \Rightarrow \lambda=2,5,5,8,10,10, \ldots$
(b) Neumann BC: $\frac{\partial u}{\partial n}=0$

Eigenfunctions $\mathfrak{u}_{j k}(x)=\cos \left(j \pi x_{1} / L\right) \cos \left(k \pi x_{2} / M\right)$ for $\mathfrak{j}, k \geq 0$.
Eigenvalues $\mu_{j k}=(j \pi / L)^{2}+(k \pi / M)^{2}$ for $j, k \geq 0$,
e.g. $\mathrm{L}=\mathrm{M}=\pi \Rightarrow \lambda=0,1,1,2,4,4, \ldots$
2. Disk $\Omega=\left\{x \in \mathbb{R}^{2}:|x|<R\right\}$.

Separate variables using polar coordinates $r, \theta$.
(a) Dirichlet BC: $u=0$

Eigenfunctions

$$
\begin{aligned}
& J_{0}\left(r j_{0, m} / R\right) \text { for } m \geq 1, \\
& J_{n}\left(r j_{n, m} r / R\right) \cos (n \theta) \text { and } J_{n}\left(r j_{n, m} r / R\right) \sin (n \theta) \text { for } n \geq 1, m \geq 1 \text {. }
\end{aligned}
$$

Notice the modes with $n=0$ are purely radial, whereas when $n \geq 1$ the modes have angular dependence.

Eigenvalues $\lambda=\left(j_{n, m} / R\right)^{2}$ for $n \geq 0, m \geq 1$, where

$$
\begin{aligned}
\mathrm{J}_{\mathrm{n}} & =\text { Bessel function of order } \mathfrak{n} \text {, and } \\
\mathfrak{j}_{n, m} & =m \text {-th positive root of } \mathrm{J}_{\mathrm{n}}(\mathrm{r})=0 .
\end{aligned}
$$

The eigenvalue $\lambda_{n, m}$ has multiplicity 2 when $n \geq 1$, associated to both cosine and sine modes.

From the graphs of the Bessel functions $\mathrm{J}_{0}, \mathrm{~J}_{1}, \mathrm{~J}_{2}$ we can read off the first 4 roots:

$$
j_{0,1} \simeq 2.40, \quad j_{1,1} \simeq 3.83 \quad j_{2,1} \simeq 5.13 \quad j_{1,2} \simeq 5.52
$$

These roots generate the first 6 eigenvalues (remembering the eigenvalues are double when $n \geq 1$ ).
(b) Neumann BC: $\frac{\partial u}{\partial n}=0$

Use roots of $J_{\mathfrak{n}}^{\prime}(r)=0$. See [Bandle, Chapter III].
3. Equilateral triangle of sidelength L.

Separation of variables fails, but one may reflect repeatedly to a hexagonal lattice whose eigenfunctions are trigonometric.

Dirichlet eigenvalues $\lambda_{j k}=\frac{16 \pi^{2}}{9 L^{2}}\left(j^{2}+j k+k^{2}\right)$ for $j, k \geq 1$.
Neumann eigenvalues $\mu_{j k}=\frac{16 \pi^{2}}{9 L^{2}}\left(j^{2}+j k+k^{2}\right)$ for $j, k \geq 0$. See [Mathews \& Walker, McCartin].


## Spectral features in 2 dim

i. Scaling: eigenvalue must balance $\Delta$, and so $\lambda \sim(\text { length scale })^{-2}$.

Precisely, $\lambda_{j}(\mathrm{t} \Omega)=\lambda_{j}(\Omega) / \mathrm{t}^{2}$.
ii. Dirichlet and Neumann spectra behave quite differently when the domain degenerates. Consider the rectangle, for example. Fix one side length L , and let the other side length M tend to 0 . Then the first positive Dirichlet eigenvalue blows up: taking $\mathfrak{j}=\mathrm{k}=1$ gives eigenvalue $(\pi / L)^{2}+(\pi / M)^{2} \rightarrow$ $\infty$. The first positive Neumann eigenvalue is constant (independent of $M$ ): taking $\mathfrak{j}=1, k=0$, gives eigenvalue $(\pi / L)^{2}$.
iii. Asymptotic: eigenvalues of the rectangle grow at a regular rate.

Proposition 2.1. (Weyl's law for rectangles) The rectangle ( $0, \mathrm{~L}$ ) $\times(0, \mathrm{M})$ has

$$
\lambda_{\mathrm{j}} \sim \mu_{\mathrm{j}} \sim \frac{4 \pi \mathfrak{j}}{\text { Area }} \quad \text { as } \mathfrak{j} \rightarrow \infty
$$

where Area $=\mathrm{LM}$ is the area of the rectangle and $\lambda_{1}, \lambda_{2}, \lambda_{3}, \ldots$ and $\mu_{1}, \mu_{2}, \mu_{3}, \ldots$ are the Dirichlet and Neumann eigenvalues respectively, in increasing order.

Proof. We give the proof for Dirichlet eigenvalues. (The Neumann case is similar.) Define for $\alpha>0$ the eigenvalue counting function

$$
\begin{aligned}
N(\alpha) & =\#\{\text { eigenvalues } \leq \alpha\} \\
& =\#\left\{j, k \geq 1: \frac{j^{2}}{\alpha L^{2} / \pi^{2}}+\frac{k^{2}}{\alpha M^{2} / \pi^{2}} \leq 1\right\} \\
& =\#\{(j, k) \in \mathbb{N} \times \mathbb{N}:(j, k) \in E\}
\end{aligned}
$$

where $E$ is the ellipse $(x / a)^{2}+(y / b)^{2} \leq 1$ and $a=\sqrt{\alpha} L / \pi, b=\sqrt{\alpha} M / \pi$.


We associate each lattice point $(\mathfrak{j}, \mathrm{k}) \in \mathrm{E}$ with the square

$$
S(j, k)=[j-1, j] \times[k-1, k]
$$

whose upper right corner lies at ( $\mathfrak{j}, \mathrm{k}$ ). These squares all lie within $E$, and so by comparing areas we find

$$
\mathrm{N}(\alpha) \leq(\text { area of } \mathrm{E} \text { in first quadrant })=\frac{1}{4} \pi \mathrm{ab}=\frac{\text { Area }}{4 \pi} \alpha
$$

On the other hand, a little thought shows that the union of the squares covers a copy of $E$ shifted down and left by one unit:

$$
\cup_{(j, k) \in E} S(j, k) \supset(E-(1,1)) \cap(\text { first quadrant }) .
$$

Comparing areas shows that

$$
\begin{aligned}
\mathrm{N}(\alpha) & \geq \frac{1}{4} \pi a b-a-b \\
& =\frac{\mathrm{LM}}{4 \pi} \alpha-\frac{\mathrm{L}+\mathrm{M}}{\pi} \sqrt{\alpha} \\
& =\frac{\text { Area }}{4 \pi} \alpha-\frac{\text { Perimeter }}{2 \pi} \sqrt{\alpha} .
\end{aligned}
$$

Combining our upper and lower estimates shows that

$$
\mathrm{N}(\alpha) \sim \frac{\text { Area }}{4 \pi} \alpha
$$

as $\alpha \rightarrow \infty$. To complete the proof we simply invert this last asymptotic, with the help of the lemma below.

Lemma 2.2. (Inversion of asymptotics) Fix $\mathrm{c}>0$. Then:

$$
N(\alpha) \sim \frac{\alpha}{c} \quad \Longrightarrow \quad \lambda_{j} \sim c j .
$$

Proof. Formally substituting $\alpha=\lambda_{j}$ and $N(\alpha)=j$ takes us from the first asymptotic to the second. The difficulty with making this substitution rigorous is that if $\lambda_{j}$ is a multiple eigenvalue, then $N\left(\lambda_{j}\right)$ can exceed $j$.

To circumvent the problem, we argue as follows. Given $\varepsilon>0$ we know from $N(\alpha) \sim \alpha / c$ that

$$
(1-\varepsilon) \frac{\alpha}{c}<N(\alpha)<(1+\varepsilon) \frac{\alpha}{c}
$$

for all large $\alpha$. Substituting $\alpha=\lambda_{j}$ into the right hand inequality implies that

$$
j<(1+\varepsilon) \frac{\lambda_{j}}{c}
$$

for all large $\mathfrak{j}$. Substituting $\alpha=\lambda_{j}-\delta$ into the left hand inequality implies that

$$
(1-\varepsilon) \frac{\lambda_{j}-\delta}{c}<j
$$

for each large j and $0<\delta<1$, and hence (by letting $\delta \rightarrow 0$ ) that

$$
(1-\varepsilon) \frac{\lambda_{j}}{c} \leq j
$$

We conclude that

$$
\frac{1}{1+\varepsilon}<\frac{\lambda_{j}}{c j} \leq \frac{1}{1-\varepsilon}
$$

for all large $\mathfrak{j}$, so that

$$
\lim _{j \rightarrow \infty} \frac{\lambda_{j}}{c j}=1
$$

as desired.
Later, in Chapter 11, we will prove Weyl's Asymptotic Law that

$$
\lambda_{j} \sim 4 \pi j / \text { Area }
$$

for all bounded domains in 2 dimensions, regardless of shape or boundary conditions.

Question to ask yourself What does a "typical" eigenfunction look like, in each of the examples above? See the following figures.


Neumann square (L=M=ת)
$x * \times * * * * * *$


## Chapter 3

## Schrödinger - computable spectra

## Goal

To study the classic examples of the harmonic oscillator (1 dim) and hydrogen atom (3 dim).

References [Strauss] Sections 9.4, 9.5, 10.7; [GustafsonSigal] Section 7.5, 7.7

## Harmonic oscillator in 1 dimension $-u^{\prime \prime}+x^{2} u=E u$

Boundary condition: $\mathfrak{u}(x) \rightarrow 0$ as $x \rightarrow \pm \infty$. (Later we give a deeper perspective, in terms of a weighted $\mathrm{L}^{2}$-space.)

Eigenfunctions $u_{k}(x)=H_{k}(x) e^{-x^{2} / 2}$ for $k \geq 0$, where $H_{k}=k$-th Hermite polynomial.

Eigenvalues $E_{k}=2 k+1$ for $k \geq 0$, or $E=1,3,5,7, \ldots$
Examples. $\mathrm{H}_{0}(x)=1, \mathrm{H}_{1}(x)=2 x, \mathrm{H}_{2}(x)=4 x^{2}-2, \mathrm{H}_{\mathrm{k}}(\mathrm{x})=(-1)^{\mathrm{k}} e^{x^{2}}\left(\frac{\mathrm{~d}}{\mathrm{~d} x}\right)^{k} e^{-x^{2}}$
Ground state: $u_{0}(x)=e^{-x^{2} / 2}=$ Gaussian. (Check: $-u_{0}^{\prime \prime}+x^{2} u_{0}=u_{0}$ )

## Quantum mechanical interpretation

If $\mathfrak{u}(x, t)$ solves the time-dependent Schrödinger equation

$$
\mathfrak{i} \mathfrak{u}_{\mathrm{t}}=-\mathfrak{u}^{\prime \prime}+\mathfrak{x}^{2} \mathfrak{u}
$$


with potential $\mathrm{V}(\mathrm{x})=\mathrm{x}^{2}$ and $\mathfrak{u}$ has $\mathrm{L}^{2}$ norm equal to 1 , then $|\mathfrak{u}|^{2}$ represents the probability density for the location of a particle in a quadratic potential well.

The $k$-th eigenfunction $\mathfrak{u}_{\mathrm{k}}(\mathrm{x})$ is called the k -th excited state, because it gives a "standing wave" solution

$$
u(x, t)=e^{-i E_{k} t} u_{k}(x)
$$

to the time-dependent equation. The higher the frequency or "energy" $E_{k}$ of the excited state, the more it can spread out in the confining potential well, as the solution plots show.

## Harmonic oscillator investigations

Method 1: ODEs Since $u_{0}(x)=e^{-x^{2} / 2}$ is an eigenfunction, we guess that all eigenfunctions decay like $e^{-x^{2} / 2}$. So we try the change of variable $u=w e^{-x^{2} / 2}$. The eigenfunction equation becomes

$$
w^{\prime \prime}-2 x w^{\prime}+(\mathrm{E}-1) w=0,
$$

which we recognize as the Hermite equation. Solving by power series, we find that the only appropriate solutions have terminating power series: they are the Hermite polynomials. (All other solutions grow like $e^{x^{2}}$ at infinity, violating the boundary condition on $\mathbf{u}$.)
Method 2: Raising and lowering Define

$$
\begin{aligned}
& \mathrm{h}^{+}=-\frac{\mathrm{d}}{\mathrm{dx}}+\mathrm{x} \quad \text { (raising or creation operator) } \\
& \mathrm{h}^{-}=\frac{\mathrm{d}}{\mathrm{dx}}+\mathrm{x} \quad \text { (lowering or annihilation operator). }
\end{aligned}
$$

Write $\mathrm{H}=-\frac{\mathrm{d}^{2}}{\mathrm{~d} x^{2}}+\mathrm{x}^{2}$ for the harmonic oscillator operator. Then one computes that

$$
\begin{aligned}
\mathrm{H} & =\mathrm{h}^{+} h^{-}+1 \\
& =h^{-} h^{+}-1
\end{aligned}
$$

Claim. If $\mathfrak{u}$ is an eigenfunction with eigenvalue $E$ then $h^{ \pm} \mathfrak{u}$ is an eigenfunction with eigenvalue $\mathrm{E} \pm 2$. (In other words, $h^{+}$"raises" the energy, and $h^{-}$ "lowers" the energy.)
Proof.

$$
\begin{aligned}
H\left(h^{+} u\right) & =\left(h^{+} h^{-}+1\right)\left(h^{+} u\right) \\
& =h^{+}\left(h^{-} h^{+}+1\right) u \\
& =h^{+}(H+2) u \\
& =h^{+}(E+2) u \\
& =(E+2) h^{+} u
\end{aligned}
$$

and similarly $H\left(h^{-} \mathfrak{u}\right)=(E-2) h^{-} \mathfrak{u}$ (exercise).
The only exception to the Claim is that $h^{-} u$ will not be an eigenfunction if $h^{-} u \equiv 0$, which occurs precisely when $u=u_{0}=e^{-x^{2} / 2}$. Thus the lowering operator annihilates the ground state.

## Relation to classical harmonic oscillator

Consider a classical oscillator with mass $m=2$, spring constant $k=2$, and displacement $x(t)$, so that $2 \ddot{x}=-2 x$. The total energy is

$$
\dot{x}^{2}+x^{2}=\text { const. }=\mathrm{E}
$$

To describe a quantum oscillator, we formally replace the momentum $\dot{\chi}$ with the "momentum operator" $-\mathfrak{i} \frac{\mathrm{d}}{\mathrm{dx}}$ and let the equation act on a function $\mathfrak{u}$ :

$$
\left[\left(-i \frac{d}{d x}\right)^{2}+x^{2}\right] u=E u
$$

This is exactly the eigenfunction equation $-\mathfrak{u}^{\prime \prime}+x^{2} u=E u$.

## Harmonic oscillator in higher dimensions $-\Delta u+|x|^{2} u=E u$

Here $|x|^{2}=x_{1}^{2}+\cdots+x_{d}^{2}$. The operator separates into a sum of 1 dimensional operators, and hence has product type eigenfunctions

$$
u=u_{k_{1}}\left(x_{1}\right) \cdots u_{k_{d}}\left(x_{d}\right), \quad E=\left(2 k_{1}+1\right)+\cdots+\left(2 k_{d}+1\right) .
$$

## Hydrogen atom in 3 dimensions $-\Delta \mathfrak{u}-\frac{2}{|x|} \mathfrak{u}=\mathrm{Eu}$

Here $\mathrm{V}(\mathrm{x})=-2 /|x|$ is an attractive electrostatic ("Coulomb") potential created by the proton in the hydrogen nucleus. (Notice the gradient of this potential gives the correct $|x|^{-2}$ inverse square law for electrostatic force.)

Boundary conditions: $\mathfrak{u}(x) \rightarrow 0$ as $|x| \rightarrow \infty$ (we will say more later about the precise formulation of the eigenvalue problem).

Eigenvalues: $E=-1,-\frac{1}{4},-\frac{1}{9}, \ldots$ with multiplicities $1,4,9, \ldots$ That is, the eigenvalue $E=-1 / n^{2}$ has multiplicity $n^{2}$.

Eigenfunctions: $e^{-r / n} L_{n}^{\ell}(r) Y_{\ell}^{m}(\theta, \phi)$ for $0 \leq|m| \leq n-1$, where $Y_{\ell}^{m}$ is a spherical harmonic and $L_{n}^{\ell}$ equals $r^{\ell}$ times a Laguerre polynomial.
(Recall the spherical harmonics are eigenfunctions of the spherical Laplacian in 3 dimensions, with $-\Delta_{\text {sphere }} Y_{\ell}^{m}=\ell(\ell+1) Y_{\ell}^{m}$. In 2 dimensions the spherical harmonics have the form $Y=\cos (k \theta)$ and $Y=\sin (k \theta)$, which satisfy $-\frac{d^{2}}{d \theta^{2}} Y=k^{2} Y$.)
Examples. The first three purely radial eigenfunctions $(\ell=\mathrm{m}=0, \mathrm{n}=$ $1,2,3)$ are $e^{-r}, e^{-r / 2}\left(1-\frac{r}{2}\right), e^{-r / 3}\left(1-\frac{2}{3} r+\frac{2}{27} r^{2}\right)$.


The corner in the graph of the eigenfunction at $\mathrm{r}=0$ is caused by the singularity of the Coulomb potential.


Continuous spectrum Eigenfunctions with positive energy E $>0$ do exist, but they oscillate as $|x| \rightarrow \infty$, and thus do not satisfy our boundary conditions. They represent "free electrons" that are not bound to the nucleus. See our later discussion of continuous spectrum, in Chapter 18.

## Chapter 4

## Discrete spectral theorem

## Goal

To state the spectral theorem for an elliptic sesquilinear form on a dense, compactly imbedded Hilbert space; and to apply this discrete spectral theorem to the Dirichlet, Robin and Neumann Laplacians.

References [Showalter] Section III. 7

## PDE preview - weak eigenfunctions

Consider the eigenfunction equation $-\Delta \mathfrak{u}=\lambda \mathfrak{u}$ for the Laplacian, in a domain $\Omega$. Multiply by a function $v \in \mathrm{H}_{0}^{1}(\Omega)$, so that $v$ equals 0 on $\partial \Omega$, and integrate to obtain

$$
-\int_{\Omega} v \Delta u \mathfrak{u d x}=\lambda \int_{\Omega} \mathrm{u} v \mathrm{dx}
$$

Green's theorem and the boundary condition on $v$ imply

$$
\int_{\Omega} \nabla \mathfrak{u} \cdot \nabla v \mathrm{dx}=\lambda\langle\mathfrak{u}, v\rangle_{\mathrm{L}^{2}(\Omega)}, \quad \forall v \in \mathrm{H}_{0}^{1}(\Omega)
$$

We call this condition the "weak form" of the eigenfunction equation. To prove existence of ONBs of such weak eigenfunctions, we first generalize to a Hilbert space problem.

## Hypotheses

Consider two infinite dimensional Hilbert spaces $\mathcal{H}$ and $\mathcal{K}$ over $\mathbb{R}$ (or $\mathbb{C}$ ).
$\mathcal{H}$ : inner product $\langle u, v\rangle_{\mathcal{H}}$, norm $\|u\|_{\mathcal{H}}$
$\mathcal{K}$ : inner product $\langle\mathfrak{u}, v\rangle_{\mathcal{K}}$, norm $\|u\|_{\mathcal{K}}$
Assume:

1. $\mathcal{K}$ is continuously and densely imbedded in $\mathcal{H}$, meaning there exists a continuous linear injection $\iota: \mathcal{K} \rightarrow \mathcal{H}$ with $\iota(\mathcal{K})$ dense in $\mathcal{H}$.
2. The imbedding $\mathcal{K} \hookrightarrow \mathcal{H}$ is compact, meaning if $B$ is a bounded subset of $\mathcal{K}$ then $B$ is precompact when considered as a subset of $\mathcal{H}$. (Equivalently, every bounded sequence in $\mathcal{K}$ has a subsequence that converges in $\mathcal{H}$.)
3. We have a map $a: \mathcal{K} \times \mathcal{K} \rightarrow \mathbb{R}$ (or $\mathbb{C}$ ) that is sesquilinear, continuous, and symmetric, meaning

$$
\begin{aligned}
& u \mapsto a(u, v) \text { is linear, for each fixed } v, \\
& v \mapsto a(u, v) \text { is linear (or conjugate linear), for each fixed } u, \\
& \qquad|a(u, v)| \leq\left(\text { const.) }\|\mathfrak{u}\|_{\mathcal{K}}\|v\|_{\mathcal{K}}\right. \\
& \qquad \mathfrak{a}(v, u)=a(u, v) \quad(\text { or } \overline{\mathfrak{a}(u, v)})
\end{aligned}
$$

4. a is elliptic on $\mathcal{K}$, meaning

$$
a(u, u) \geq c\|u\|_{\mathcal{K}}^{2} \quad \forall u \in \mathcal{K}
$$

for some $c>0$. Hence $a(u, u) \asymp\|u\|_{\mathcal{K}}^{2}$.
Consequence of symmetry and ellipticity:
$a(u, v)$ defines an inner product whose norm is equivalent to the $\|\cdot\|_{\mathcal{K}}$-norm.

## Spectral theorem

Theorem 4.1. Under the assumptions above, there exist vectors $\mathfrak{u}_{1}, \mathfrak{u}_{2}, \mathfrak{u}_{3}, \ldots \in$ $\mathcal{K}$ and numbers

$$
0<\gamma_{1} \leq \gamma_{2} \leq \gamma_{3} \leq \cdots \rightarrow \infty
$$

such that:

- $\mathfrak{u}_{\mathfrak{j}}$ is an eigenvector of $\mathfrak{a}(\cdot, \cdot)$ with eigenvalue $\gamma_{j}$, meaning

$$
\begin{equation*}
a\left(u_{j}, v\right)=\gamma_{j}\left\langle u_{j}, v\right\rangle_{\mathcal{H}} \quad \forall v \in \mathcal{K}, \tag{4.1}
\end{equation*}
$$

- $\left\{\mathbf{u}_{j}\right\}$ is an ONB for $\mathcal{H}$,
- $\left\{\mathrm{u}_{\mathrm{j}} / \sqrt{\gamma_{\mathrm{j}}}\right\}$ is an ONB for $\mathcal{K}$ with respect to the $\mathbf{a}$-inner product.

The idea is to show that a certain "inverse" operator associated with a is compact and selfadjoint on $\mathcal{H}$. This approach makes sense in terms of differential equations, where a would correspond to a differential operator such as $-\Delta$ (which is unbounded) and the inverse would correspond to an integral operator $(-\Delta)^{-1}$ (which is bounded, and in fact compact, on suitable domains). Indeed, we will begin by solving the analogue of $-\Delta \mathfrak{u}=f$ weakly, in our Hilbert space setting.

Proof. We first claim that for each $\mathrm{f} \in \mathcal{H}$ there exists a unique $\boldsymbol{u} \in \mathcal{K}$ such that

$$
\begin{equation*}
\mathrm{a}(\mathrm{u}, v)=\langle\mathrm{f}, v\rangle_{\mathcal{H}} \quad \forall v \in \mathcal{K} . \tag{4.2}
\end{equation*}
$$

Furthermore, the map

$$
\begin{gathered}
\mathrm{B}: \mathcal{H} \rightarrow \mathcal{K} \\
\mathrm{f} \mapsto \mathrm{u}
\end{gathered}
$$

is linear and bounded. To prove this claim, fix $\mathrm{f} \in \mathcal{H}$ and define a bounded linear functional $\mathrm{F}(v)=\langle v, \mathrm{f}\rangle_{\mathcal{H}}$ on $\mathcal{K}$, noting for the boundedness that

$$
\begin{aligned}
|\mathrm{F}(v)| & \leq\|v\|_{\mathcal{H}}\|\mathrm{f}\|_{\mathcal{H}} \\
& \leq \text { (const.) }\|v\|_{\mathcal{K}}\|\mathrm{f}\|_{\mathcal{H}} \quad \text { since } \mathcal{K} \text { is imbedded in } \mathcal{H} \\
& \leq \text { (const.) } \mathfrak{a}(v, v)^{1 / 2}\|f\|_{\mathcal{H}}
\end{aligned}
$$

by ellipticity. Hence by the Riesz Representation Theorem on $\mathcal{K}$ (with respect to the $a$-inner product and norm on $\mathcal{K}$ ), there exists a unique $u \in \mathcal{K}$ such that $\mathrm{F}(v)=\mathrm{a}(v, u)$ for all $v \in \mathcal{K}$. That is,

$$
\langle v, f\rangle_{\mathcal{H}}=\mathrm{a}(v, u) \quad \forall v \in \mathcal{K}
$$

as desired for (4.2). Thus the map B: $\mathfrak{f} \mapsto \boldsymbol{u}$ is well defined. Clearly it is linear. And

$$
a(u, u)=|F(u)| \leq(\text { const. }) a(u, u)^{1 / 2}\|f\|_{\mathcal{H}}
$$

Hence $\mathfrak{a}(u, u)^{1 / 2} \leq($ const. $)\|f\|_{\mathcal{H}}$, so that $B$ is bounded from $\mathcal{H}$ to $\mathcal{K}$, which proves our initial claim.

Next, $\mathrm{B}: \mathcal{H} \rightarrow \mathcal{K} \rightarrow \mathcal{H}$ is compact, since $\mathcal{K}$ imbeds compactly into $\mathcal{H}$. Further, B is selfadjoint on $\mathcal{H}$, since for all $\mathrm{f}, \mathrm{g} \in \mathcal{H}$ we have

$$
\begin{aligned}
\langle\mathrm{Bf}, \mathrm{~g}\rangle_{\mathcal{H}} & =\overline{\langle\mathrm{g}, \mathrm{Bf}\rangle_{\mathcal{H}}} & & \\
& =\overline{\mathrm{a}(\mathrm{Bg}, \mathrm{Bf})} & & \text { by definition of } \mathrm{B}, \\
& =\mathrm{a}(\mathrm{Bf}, \mathrm{Bg}) & & \text { by symmetry of } \mathrm{a}, \\
& =\langle\mathrm{f}, \mathrm{Bg}\rangle_{\mathcal{H}} & & \text { by definition of } \mathrm{B},
\end{aligned}
$$

which implies $B^{*}=B$.
Hence the spectral theorem for compact, self-adjoint operators [Evans, App. D] provides an ONB for $\mathcal{H}$ consisting of eigenvectors of B , with

$$
B u_{j}=\widetilde{\gamma}_{j} u_{j}
$$

for some eigenvalues $\widetilde{\gamma}_{j} \rightarrow 0$.
The eigenvalues are all nonzero, because B is injective: $\mathrm{Bf}=0$ would imply $\langle\mathrm{f}, v\rangle_{\mathcal{H}}=0$ for all $v \in \mathcal{K}$ by (4.2), so that $\mathrm{f}=0$ (using density of $\mathcal{K}$ in $\mathcal{H})$.

Since we may divide by the eigenvalue, we deduce that $u_{j}=B\left(u_{j} / \widetilde{\gamma}_{j}\right)$ belongs to the range of $B$, and hence $u_{j} \in \mathcal{K}$.

The eigenvalues are all positive, in fact, since

$$
\widetilde{\gamma}_{j} a\left(u_{j}, v\right)=a\left(B u_{j}, v\right)=\left\langle u_{j}, v\right\rangle_{\mathcal{H}} \quad \forall v \in \mathcal{K}
$$

and choosing $v=u_{j} \in \mathcal{K}$ and using ellipticity shows that $\widetilde{\gamma}_{j}>0$. Thus we see that the reciprocal numbers $0<\gamma_{j} \stackrel{\text { def }}{=} 1 / \widetilde{\gamma}_{j} \rightarrow \infty$ satisfy

$$
a\left(u_{j}, v\right)=\gamma_{j}\left\langle u_{j}, v\right\rangle_{\mathcal{H}} \quad \forall v \in \mathcal{K}
$$

which is (4.1).
Finally, we have $\alpha$-orthonormality of the set $\left\{\mathbf{u}_{j} / \sqrt{\gamma_{j}}\right\}$ :

$$
\begin{aligned}
a\left(u_{j}, u_{k}\right) & =\gamma_{j}\left\langle u_{j}, u_{k}\right\rangle_{\mathcal{H}} \\
& =\gamma_{j} \delta_{j k} \\
& =\sqrt{\gamma_{j}} \sqrt{\gamma_{k}} \delta_{j k} .
\end{aligned}
$$

This orthonormal set is complete in $\mathcal{K}$, because if $a\left(u_{j}, v\right)=0$ for all $j$ then $\left\langle u_{j}, v\right\rangle_{\mathcal{H}}=0$ for all $\mathfrak{j}$, by (4.1), so that $v=0$.

Remark. Eigenvectors corresponding to distinct eigenvalues are automatically orthogonal, since

$$
\begin{aligned}
\left(\gamma_{j}-\gamma_{k}\right)\left\langle u_{j}, u_{k}\right\rangle_{\mathcal{H}} & =\gamma_{j}\left\langle u_{j}, u_{k}\right\rangle_{\mathcal{H}}-\overline{\gamma_{k}\left\langle u_{k}, u_{j}\right\rangle_{\mathcal{H}}} \\
& =a\left(u_{j}, u_{k}\right)-\overline{a\left(u_{k}, u_{j}\right)} \\
& =0
\end{aligned}
$$

by symmetry of $a$.

## Chapter 5

## Application: ONBs of Laplace eigenfunctions

## Goal

To apply the spectral theorem from the previous chapter to the Dirichlet, Robin and Neumann Laplacians, and to the fourth order biLaplacian.

## Laplacian

## Dirichlet Laplacian

$$
\begin{aligned}
-\Delta \mathfrak{u} & =\lambda \mathfrak{u} & & \text { in } \Omega \\
\mathfrak{u} & =0 & & \text { on } \partial \Omega
\end{aligned}
$$

$\Omega=$ bounded domain in $\mathbb{R}^{\mathrm{d}}$.
$\mathcal{H}=\mathrm{L}^{2}(\Omega)$, inner product $\langle u, v\rangle_{\mathrm{L}^{2}}=\int_{\Omega} u v \mathrm{dx}$.
$\mathcal{K}=\mathrm{H}_{0}^{1}(\Omega)=$ Sobolev space, which is the completion of $\mathrm{C}_{0}^{\infty}(\Omega)$ (smooth functions equalling zero on a neighborhood of the boundary) under the inner product

$$
\langle\mathfrak{u}, v\rangle_{\mathrm{H}^{1}}=\int_{\Omega}[\nabla \mathbf{u} \cdot \nabla v+\mathfrak{u v}] \mathrm{dx}
$$

Density: $\mathrm{H}_{0}^{1}$ contains $\mathrm{C}_{0}^{\infty}$, which is dense in $\mathrm{L}^{2}$.

Continuous imbedding $\mathrm{H}_{0}^{1} \hookrightarrow \mathrm{~L}^{2}$ is trivial:

$$
\begin{aligned}
\|u\|_{L^{2}} & =\left(\int_{\Omega} u^{2} d x\right)^{1 / 2} \\
& \leq\left(\int_{\Omega}\left[|\nabla u|^{2}+u^{2}\right] d x\right)^{1 / 2} \\
& =\|u\|_{H^{1}}
\end{aligned}
$$

Compact imbedding: $\mathrm{H}_{0}^{1} \hookrightarrow \mathrm{~L}^{2}$ compactly by the Rellich-Kondrachov Theorem [GilbargTrudinger, Theorem 7.22].

Sesquilinear form: define

$$
\mathrm{a}(\mathrm{u}, v)=\int_{\Omega} \nabla \mathfrak{u} \cdot \nabla v \mathrm{dx}+\int_{\Omega} \mathfrak{u v d x}=\langle\mathfrak{u}, v\rangle_{\mathrm{H}^{1}}, \quad \mathfrak{u}, v \in \mathrm{H}_{0}^{1}(\Omega)
$$

Clearly a is symmetric and continuous on $\mathrm{H}_{0}^{1}(\Omega)$.
Ellipticity: $\mathfrak{a}(u, u)=\|u\|_{\mathrm{H}^{1}}^{2}$
The Spectral Theorem 4.1 gives an $\operatorname{ONB}\left\{\mathfrak{u}_{j}\right\}$ for $L^{2}(\Omega)$ and corresponding eigenvalues which we denote $\gamma_{j}=\lambda_{j}+1>0$ satisfying

$$
\left\langle u_{j}, v\right\rangle_{\mathrm{H}^{1}}=\left(\lambda_{\mathrm{j}}+1\right)\left\langle u_{\mathrm{j}}, v\right\rangle_{\mathrm{L}^{2}} \quad \forall v \in \mathrm{H}_{0}^{1}(\Omega) .
$$

Equivalently,

$$
\int_{\Omega} \nabla \mathfrak{u}_{\mathrm{j}} \cdot \nabla v \mathrm{~d} x=\lambda_{\mathrm{j}} \int_{\Omega} \mathrm{u}_{\mathrm{j}} v \mathrm{~d} x \quad \forall v \in \mathrm{H}_{0}^{1}(\Omega)
$$

That is,

$$
-\Delta \mathfrak{u}_{j}=\lambda_{j} u_{j}
$$

weakly, so that $\mathfrak{u}_{j}$ is a weak eigenfunction of the Laplacian with eigenvalue $\lambda_{j}$. Elliptic regularity theory gives that $\mathfrak{u}_{\mathrm{j}}$ is $\mathrm{C}^{\infty}$-smooth in $\Omega$ [GilbargTrudinger, Corollary 8.11], and hence satisfies the eigenfunction equation classically. The boundary condition $\mathfrak{u}_{j}=0$ is satisfied in the sense of Sobolev spaces (since $\mathrm{H}_{0}^{1}$ is the closure of $\mathrm{C}_{0}^{\infty}$ ), and is satisfied classically on any smooth portion of $\partial \Omega$, again by elliptic regularity.

The eigenvalues are nonnegative, with

$$
\lambda_{j}=\frac{\int_{\Omega}\left|\nabla u_{j}\right|^{2} d x}{\int_{\Omega} u_{j}^{2} d x} \geq 0
$$

as we see by choosing $v=\mathfrak{u}_{\mathrm{j}}$ in the weak formulation.
Further, $\lambda_{j}>0$ because: if $\lambda_{j}=0$ then $\left|\nabla \mathfrak{u}_{j}\right| \equiv 0$ by the last formula, so that $u_{j} \equiv 0$ by the Sobolev inequality for $\mathrm{H}_{0}^{1}$ [GilbargTrudinger, Theorem 7.10], but $\mathfrak{u}_{\mathrm{j}}$ cannot vanish identically because it has $\mathrm{L}^{2}$-norm equal to 1 . Hence

$$
0<\lambda_{1} \leq \lambda_{2} \leq \lambda_{3} \leq \cdots \rightarrow \infty
$$

Aside. The Sobolev inequality we used is easily proved: for $u \in H_{0}^{1}(\Omega)$,

$$
\begin{aligned}
\|u\|_{L^{2}}^{2} & =\int_{\Omega} u^{2} d x \\
& =-\int_{\Omega} 2 x_{i} u \frac{\partial u}{\partial x_{i}} d x \quad \text { by parts } \\
& \leq 2\left(\max _{x \in \bar{\Omega}}|x|\right)\|u\|_{L^{2}}\left\|\partial u / \partial x_{i}\right\|_{L^{2}} \\
& \leq(\text { const. })\|u\|_{L^{2}}\|\nabla u\|_{L^{2}}
\end{aligned}
$$

so that we have a Sobolev inequality

$$
\|\mathfrak{u}\|_{\mathrm{L}^{2}} \leq \text { (const.) }\|\nabla \mathfrak{u}\|_{\mathrm{L}^{2}} \quad \forall \mathfrak{u} \in \mathrm{H}_{0}^{1}(\Omega)
$$

where the constant depends on the domain $\Omega$. Incidentally, this Sobolev inequality provides another proof that $\lambda_{j}>0$ for the Dirichlet Laplacian.

## Neumann Laplacian

$$
\begin{aligned}
-\Delta \mathfrak{u} & =\mu u & & \text { in } \Omega \\
\frac{\partial \mathfrak{u}}{\partial \mathrm{n}} & =0 & & \text { on } \partial \Omega
\end{aligned}
$$

$\Omega=$ bounded domain in $\mathbb{R}^{\mathrm{d}}$ with Lipschitz boundary.
$\mathcal{H}=\mathrm{L}^{2}(\Omega)$
$\mathcal{K}=\mathrm{H}^{1}(\Omega)=$ Sobolev space, which is the completion of $\mathrm{C}^{\infty}(\bar{\Omega})$ under the inner product $\langle\boldsymbol{u}, v\rangle_{\mathrm{H}^{1}}$ (see [GilbargTrudinger, p. 174]).

Argue as for the Dirichlet Laplacian. The compact imbedding is provided by the Rellich-Kondrachov Theorem [GilbargTrudinger, Theorem 7.26], which relies on Lipschitz smoothness of the boundary.

One writes the eigenvalues in the Spectral Theorem 4.1 as $\gamma_{j}=\mu_{j}+1>0$ and finds

$$
\begin{equation*}
\int_{\Omega} \nabla \mathrm{u}_{\mathrm{j}} \cdot \nabla v \mathrm{~d} x=\mu_{\mathrm{j}} \int_{\Omega} \mathrm{u}_{\mathrm{j}} v \mathrm{~d} x \quad \forall v \in \mathrm{H}^{1}(\Omega) \tag{5.1}
\end{equation*}
$$

which implies that

$$
-\Delta \mathfrak{u}_{j}=\mu_{j} \mathfrak{u}_{j}
$$

weakly (and hence classically). In fact (5.1) says a little more, because it holds for all $v \in H^{1}(\Omega)$, not just for $v \in H_{0}^{1}(\Omega)$ as needed for a weak solution. We will use this additional information in the next chapter to show that eigenfunctions automatically satisfy the Neumann boundary condition (even though we never imposed it)!

Choosing $v=u_{j}$ proves $\mu_{j} \geq 0$. The first Neumann eigenvalue is zero: $\mu_{1}=0$, with a constant eigenfunction $u_{1} \equiv$ const. $\neq 0$. (This constant function belongs to $\mathrm{H}^{1}(\Omega)$, although not to $\mathrm{H}_{0}^{1}(\Omega)$.) Hence

$$
0=\mu_{1} \leq \mu_{2} \leq \mu_{3} \leq \cdots \rightarrow \infty
$$

## Robin Laplacian

$$
\begin{aligned}
-\Delta u & =\rho u & & \text { in } \Omega \\
\frac{\partial u}{\partial \mathfrak{n}}+\sigma u & =0 & & \text { on } \partial \Omega
\end{aligned}
$$

$\Omega=$ bounded domain in $\mathbb{R}^{\mathrm{d}}$ with Lipschitz boundary.
$\mathcal{H}=\mathrm{L}^{2}(\Omega)$
$\mathcal{K}=\mathrm{H}^{1}(\Omega)$
$\sigma>0$ is the Robin constant.
The density and compact imbedding conditions are as in the Neumann case above.

Before defining the sesquilinear form, we need to make sense of the boundary values of $\mathbf{u}$. Sobolev functions do have well defined boundary values. More precisely, there is a bounded linear operator (called the trace operator) $\mathrm{T}: \mathrm{H}^{1}(\Omega) \rightarrow \mathrm{L}^{2}(\partial \Omega)$ such that

$$
\begin{equation*}
\|T u\|_{L^{2}(\partial \Omega)} \leq \tau\|u\|_{H^{1}(\Omega)} \tag{5.2}
\end{equation*}
$$

for some $\tau>0$, and with the property that if $u$ extends to a continuous function on $\partial \Omega$, then $T u=u$ on $\partial \Omega$. (Thus the trace operator truly captures the boundary values of $\mathfrak{u}$.) Further, if $u \in H_{0}^{1}(\Omega)$ then $T u=0$, meaning that functions in $\mathrm{H}_{0}^{1}$ "equal zero on the boundary". For these trace results, see [Evans, Section 5.5] for domains with $C^{1}$ boundary, or [EvansGariepy, §4.3] for the slightly rougher case of Lipschitz boundary.

Sesquilinear form:

$$
\mathrm{a}(\mathrm{u}, v)=\int_{\Omega} \nabla \mathfrak{u} \cdot \nabla v \mathrm{~d} x+\sigma \int_{\partial \Omega} \mathfrak{u v d S}(x)+\int_{\Omega} \mathfrak{u v d x}
$$

(where $u$ and $v$ on the boundary should be interpreted as the trace values Tu and $T v$ ). Clearly $a$ is symmetric and continuous on $\mathrm{H}^{1}(\Omega)$.

Ellipticity: $\mathfrak{a}(u, u) \geq\|u\|_{\mathbf{H}^{1}}^{2}$, since $\sigma>0$.
One writes the eigenvalues in the Spectral Theorem 4.1 as $\gamma_{j}=\rho_{j}+1>0$ and finds

$$
\int_{\Omega} \nabla \mathfrak{u}_{\mathrm{j}} \cdot \nabla v \mathrm{~d} x+\sigma \int_{\partial \Omega} \mathfrak{u}_{\mathrm{j}} v \mathrm{dS}(\mathrm{x})=\rho_{\mathrm{j}} \int_{\Omega} \mathrm{u}_{\mathrm{j}} v \mathrm{~d} x \quad \forall v \in \mathrm{H}^{1}(\Omega),
$$

which implies that

$$
-\Delta u_{j}=\rho_{j} u_{j}
$$

weakly and hence classically. For the weak solution here we need (by definition) only to use trial functions $v \in \mathrm{H}_{0}^{1}(\Omega)$ (functions equalling zero on the boundary). In the next chapter we use the full class $v \in \mathrm{H}^{1}(\Omega)$ to show that the eigenfunctions satisfy the Robin boundary condition.

Choosing $v=u_{j}$ proves

$$
\rho_{j}=\frac{\int_{\Omega}\left|\nabla u_{j}\right|^{2} d x+\sigma \int_{\partial \Omega} u_{j}^{2} d S(x)}{\int_{\Omega} u_{j}^{2} d x} \geq 0
$$

using again that $\sigma>0$. Further, $\rho_{j}>0$ because: if $\rho_{j}=0$ then $\left|\nabla u_{j}\right| \equiv 0$ so that $u_{j} \equiv$ const., and this constant must equal zero because $\int_{\partial \Omega} u_{j}^{2} d S(x)=0$; but $u_{j}$ cannot vanish identically because it has $L^{2}$-norm equal to 1 . Hence when $\sigma>0$ we have

$$
0<\rho_{1} \leq \rho_{2} \leq \rho_{3} \leq \cdots \rightarrow \infty
$$

Negative Robin constant: $\sigma<0$. Ellipticity more difficult to prove when $\sigma<0$. We start by controlling the boundary values in terms of the gradient and $\mathrm{L}^{2}$ norm. We have

$$
\int_{\partial \Omega} u^{2} \mathrm{dS}(\mathrm{x}) \leq \text { (const.) } \int_{\Omega}|\nabla \mathfrak{u}||\mathfrak{u}| d x+(\text { const. }) \int_{\Omega} u^{2} d x
$$

as one sees by inspecting the proof of the trace theorem ([Evans, §5.5] or [EvansGariepy, §4.3]). An application of Cauchy-with- $\varepsilon$ gives

$$
\int_{\partial \Omega} u^{2} \mathrm{dS}(\mathrm{x}) \leq \varepsilon\|\nabla u\|_{\mathrm{L}^{2}}^{2}+\mathrm{C}\|u\|_{\mathrm{L}^{2}}^{2}
$$

for some constant $C=C(\varepsilon)>0$ (independent of $u$ ). Let us choose $\varepsilon=$ $1 / 2|\sigma|$, so that

$$
a(u, u) \geq \frac{1}{2}\|u\|_{H^{1}}^{2}-C|\sigma|\|u\|_{L^{2}}^{2} .
$$

Hence the new sesquilinear form $\widetilde{\mathfrak{a}}(u, v)=\mathrm{a}(u, v)+\mathrm{C}|\sigma|\langle u, v\rangle_{\mathrm{L}^{2}}$ is elliptic. We apply the discrete spectral theorem to this new form, and then obtain the eigenvalues of a by subtracting $\mathrm{C}|\sigma|$ (with the same ONB of eigenfunctions).

## Eigenfunction expansions in the $\mathrm{L}^{2}$ and $\mathrm{H}^{1}$ norms

The $L^{2}$-ONB of eigenfunctions $\left\{u_{j}\right\}$ of the Laplacian satisfies

$$
\begin{equation*}
f=\sum_{j}\left\langle f, u_{j}\right\rangle_{L^{2}} u_{j} \tag{5.3}
\end{equation*}
$$

with convergence in $L^{2}(\Omega)$, for all $f$ in the following spaces:

$$
f \in \begin{cases}\mathrm{H}_{0}^{1}(\Omega) & \text { for Dirichlet } \\ \mathrm{H}^{1}(\Omega) & \text { for Neumann } \\ \mathrm{H}^{1}(\Omega) & \text { for Robin. }\end{cases}
$$

Importantly, this expansion (5.3) converges not only in the $\mathrm{L}^{2}$-norm, but also in the $\mathrm{H}^{1}$-norm.
Proof. In the Dirichlet case we have

$$
\left\langle f, u_{j}\right\rangle_{L^{2}}=\frac{1}{\lambda_{j}+1} a\left(f, u_{j}\right)
$$

by choosing $v=\mathrm{f}$ in the eigenfunction equation. Hence

$$
\left\langle f, u_{j}\right\rangle_{L^{2}} u_{j}=a\left(f, u_{j} / \sqrt{\lambda_{j}+1}\right) u_{j} / \sqrt{\lambda_{j}+1} .
$$

We know $\left\{u_{j} / \sqrt{\lambda_{j}+1}\right\}$ is an ONB for $\mathcal{K}=H_{0}^{1}(\Omega)$ under the $\mathfrak{a}$-inner product, by Theorem 4.1, and so expansion (5.3) converges in the $\mathrm{H}^{1}$-norm. The argument is exactly the same for the Neumann and Robin applications, except using $\mathrm{H}^{1}(\Omega)$.

## Invariance of eigenvalues under translation, rotation and reflection, and scaling under dilation

Eigenvalues of the Laplacian remain invariant when the domain $\Omega$ is translated, rotated or reflected, as one sees by a straightforward change of variable in either the classical or weak formulation of the eigenvalue problem. Physically, this invariance simply means that a vibrating membrane is unaware of any coordinate system we impose upon it.

Dilations do change the eigenvalues, of course, by a simple rescaling relation: $\lambda_{j}(t \Omega)=t^{-2} \lambda_{j}(\Omega)$ for each $j$ and all $t>0$, and similarly for the Robin and Neumann eigenvalues. We understand the scale factor $t^{-2}$ physically by recalling that large drums vibrate at low tones.

## BiLaplacian - vibrating plates

The fourth order wave equation $\phi_{\mathrm{tt}}=-\Delta \Delta \phi$ describes the transverse vibrations of a rigid plate. (In one dimension, this equation simplifies to the beam equation: $\left.\phi_{\mathrm{tt}}=-\phi^{\prime \prime \prime \prime}\right)$. After separating out the time variable, one arrives at the eigenvalue problem for the biLaplacian:

$$
\Delta \Delta \mathfrak{u}=\Lambda \mathfrak{u} \quad \text { in } \Omega
$$

We will prove existence of an orthonormal basis of eigenfunctions. For simplicity, we treat only the Dirichlet case, which has boundary conditions

$$
\mathfrak{u}=|\nabla \mathfrak{u}|=0 \quad \text { on } \partial \Omega .
$$

(The Neumann "natural" boundary conditions are rather complicated, for the biLaplacian.)
$\Omega=$ bounded domain in $\mathbb{R}^{\text {d }}$
$\mathcal{H}=\mathrm{L}^{2}(\Omega)$
$\mathcal{K}=\mathrm{H}_{0}^{2}(\Omega)=$ completion of $\mathrm{C}_{0}^{\infty}(\Omega)$ under the inner product

$$
\langle u, v\rangle_{H^{2}}=\int_{\Omega}\left[\sum_{m, n=1}^{d} u_{x_{m} x_{n}} v_{x_{m} x_{n}}+\sum_{m=1}^{d} u_{x_{m}} v_{x_{m}}+u v\right] d x .
$$

Density: $\mathrm{H}_{0}^{2}$ contains $\mathrm{C}_{0}^{\infty}$, which is dense in $\mathrm{L}^{2}$.
Compact imbedding: $H_{0}^{2} \hookrightarrow H_{0}^{1} \hookrightarrow L^{2}$ and the second imbedding is compact.

Sesquilinear form: define

$$
a(u, v)=\int_{\Omega}\left[\sum_{m, n=1}^{d} u_{x_{m} x_{n}} v_{x_{m} x_{n}}+u v\right] d x, \quad u, v \in H_{0}^{2}(\Omega)
$$

Clearly $a$ is symmetric and continuous on $H_{0}^{2}(\Omega)$.
Ellipticity: $\|\mathfrak{u}\|_{\mathrm{H}^{2}}^{2} \leq(\mathrm{d}+1) \mathfrak{a}(u, u)$, because integration by parts gives

$$
\begin{aligned}
\int_{\Omega} \sum_{m=1}^{d} u_{x_{m}}^{2} d x & =-\sum_{m=1}^{d} \int_{\Omega} u_{x_{m} x_{m}} u d x \\
& \leq \sum_{m=1}^{d} \int_{\Omega}\left[u_{x_{m} x_{m}}^{2}+u^{2}\right] d x \\
& \leq a(u, u) d
\end{aligned}
$$

The Spectral Theorem 4.1 gives an $\operatorname{ONB}\left\{\mathfrak{u}_{\mathfrak{j}}\right\}$ for $\mathrm{L}^{2}(\Omega)$ and corresponding eigenvalues which we denote $\gamma_{j}=\Lambda_{j}+1>0$ satisfying

$$
\mathrm{a}\left(u_{\mathrm{j}}, v\right)=\left(\Lambda_{\mathrm{j}}+1\right)\left\langle u_{j}, v\right\rangle_{\mathrm{L}^{2}} \quad \forall v \in \mathrm{H}_{0}^{2}(\Omega)
$$

Equivalently,

$$
\int_{\Omega} \sum_{m, n=1}^{\mathrm{d}}\left(u_{j}\right)_{x_{\mathrm{m}} x_{n}} v_{x_{\mathrm{m}} x_{n}} \mathrm{~d} x=\Lambda_{\mathrm{j}} \int_{\Omega} u_{j} v \mathrm{~d} x \quad \forall v \in \mathrm{H}_{0}^{2}(\Omega)
$$

That is,

$$
\sum_{m, n=1}^{d}\left(u_{j}\right)_{x_{m} x_{m} x_{n} x_{n}}=\Lambda_{j} u_{j}
$$

weakly, which says

$$
\Delta \Delta u_{j}=\Lambda_{j} u_{j}
$$

weakly. Hence $\mathfrak{u}_{\mathrm{j}}$ is a weak eigenfunction of the biLaplacian with eigenvalue $\Lambda_{j}$. Elliptic regularity gives that $u_{j}$ is $C^{\infty}$-smooth, and hence satisfies the eigenfunction equation classically. The boundary condition $\mathfrak{u}_{j}=\left|\nabla \mathfrak{u}_{j}\right|=0$ is satisfied in the sense of Sobolev spaces (since $\boldsymbol{u}_{j}$ and each partial derivative $\left(u_{j}\right)_{\chi_{m}}$ belong to $\left.H_{0}^{1}\right)$, and the boundary condition is satisfied classically on any smooth portion of $\partial \Omega$, again by elliptic regularity.

The eigenvalues are nonnegative, with

$$
\Lambda_{j}=\frac{\int_{\Omega}\left|D^{2} u_{j}\right|^{2} d x}{\int_{\Omega} u_{j}^{2} d x} \geq 0
$$

as we see by choosing $v=u_{j}$ in the weak formulation and writing $D^{2} u=$ [ $\left.\mathrm{u}_{\mathrm{x}_{\mathfrak{m}} \mathrm{x}_{\mathrm{n}}}\right]_{\mathfrak{m}, \mathfrak{n}=1}^{\mathrm{d}}$ for the Hessian matrix.

Further, $\Lambda_{j}>0$ because: if $\Lambda_{j}=0$ then $\left(u_{j}\right)_{x_{m} x_{n}} \equiv 0$ by the last formula, so that $\left(u_{j}\right)_{x_{m}} \equiv 0$ by the Sobolev inequality for $H_{0}^{1}$ applied to $\left(u_{j}\right)_{x_{m}}$, and hence $\mathfrak{u}_{j} \equiv 0$ by the same Sobolev inequality, which gives a contradiction. Hence

$$
0<\Lambda_{1} \leq \Lambda_{2} \leq \Lambda_{3} \leq \cdots \rightarrow \infty
$$

## Compact resolvents

The essence of the proof of the Spectral Theorem 4.1 is to show that the inverse operator B is compact, which means for our differential operators that the inverse is a compact integral operator. For example, in the Neumann Laplacian application we see that $(-\Delta+1)^{-1}$ is compact from $L^{2}(\Omega)$ to $\mathrm{H}^{1}(\Omega)$. So is $(-\Delta+\alpha)^{-1}$ for any positive $\alpha$, but $\alpha=0$ does not give an invertible operator because the Neumann Laplacian has nontrivial kernel, with $-\Delta(c)=0$ for every constant $c$.

Thus for the Neumann Laplacian, the resolvent operator

$$
R_{\lambda}=(-\Delta-\lambda)^{-1}
$$

is compact whenever $\lambda$ is negative.

## Chapter 6

## Natural boundary conditions

## Goal

To understand how the Neumann and Robin boundary conditions arise "naturally" from the weak eigenfunction equation.

## Dirichlet boundary conditions

are imposed directly by our choice of function space $H_{0}^{1}(\Omega)$, since each function in that space is a limit of functions with compact support in $\Omega$.

## Neumann boundary conditions

The weak form of the Neumann eigenequation for the Laplacian, from Chapter 5, is:

$$
\begin{equation*}
\int_{\Omega} \nabla \mathfrak{u} \cdot \nabla v \mathrm{~d} x=\mu \int_{\Omega} u v \mathrm{~d} x \quad \forall v \in \mathrm{H}^{1}(\Omega) . \tag{6.1}
\end{equation*}
$$

From this formula we showed that $-\Delta \mathfrak{u}=\mu u$ weakly and hence classically, by using only functions $v$ that vanish on the boundary, meaning $v \in \mathrm{H}_{0}^{1}(\Omega)$.

To deduce the Neumann boundary condition $\partial u / \partial n=0$, we will take $v$ not to vanish on the boundary. Assume for simplicity that the boundary is smooth, so that $u$ extends smoothly to $\bar{\Omega}$. Green's formula (integration by parts) applied to (6.1) implies that

$$
\int_{\Omega}(-\Delta u) v \mathrm{~d} x+\int_{\partial \Omega} \frac{\partial u}{\partial n} v \mathrm{dS}=\int_{\Omega}(\mu u) v \mathrm{~d} x \quad \forall v \in \mathrm{C}^{\infty}(\bar{\Omega})
$$

Since $-\Delta \mathfrak{u}=\mu u$, we deduce

$$
\int_{\partial \Omega} \frac{\partial u}{\partial \mathrm{n}} v \mathrm{dS}=0 \quad \forall v \in \mathrm{C}^{\infty}(\bar{\Omega})
$$

One may choose $v \in \mathrm{C}^{\infty}(\bar{\Omega})$ to equal the normal derivative of $u$ on the boundary (meaning $\left.v\right|_{\partial \Omega}=\frac{\partial u}{\partial n}$ ), or alternatively one may use density of $\left.\mathrm{C}^{\infty}(\bar{\Omega})\right|_{\partial \Omega}$ in $L^{2}(\partial \Omega)$; either way one concludes that

$$
\frac{\partial u}{\partial n}=0 \quad \text { on } \partial \Omega
$$

which is the Neumann boundary condition.
Note. If the boundary is only piecewise smooth, then one merely applies the above reasoning on the smooth portions of the boundary, to show the Neumann condition holds there.

## Robin boundary conditions

Integrating by parts in the Robin eigenfunction equation

$$
\int_{\Omega} \nabla \mathfrak{u} \cdot \nabla v \mathrm{dx}+\sigma \int_{\partial \Omega} \mathrm{uvdS}=\rho \int_{\Omega} \mathrm{uvdx} \quad \forall v \in \mathrm{H}^{1}(\Omega)
$$

(that is, applying Green's formula to this equation) and then using that $-\Delta u=\rho u$ gives that

$$
\int_{\partial \Omega}\left(\frac{\partial u}{\partial n}+\sigma u\right) v \mathrm{dS}=0 \quad \forall v \in \mathrm{C}^{\infty}(\bar{\Omega})
$$

Like above, we obtain the Robin boundary condition

$$
\frac{\partial u}{\partial n}+\sigma u=0 \quad \text { on } \partial \Omega
$$

at least on smooth portions of the boundary.

## BiLaplacian — natural boundary conditions

Natural boundary conditions for the biLaplacian can be derived similarly [Chasman, §5]. They are much more complicated than for the Laplacian.

## Chapter 7

## Application: ONB of eigenfunctions for the Laplacian with magnetic field

## Goals

To apply the spectral theorem from Chapter 4 to the magnetic Laplacian (the Schrödinger operator for a particle in the presence of a classical magnetic field).

## Magnetic Laplacian

Take a bounded domain $\Omega$ in $\mathbb{R}^{\mathrm{d}}$, with $\mathrm{d}=2$ or $\mathrm{d}=3$. We seek an ONB of eigenfunctions and eigenvalues for the magnetic Laplacian

$$
\begin{aligned}
(i \nabla+\vec{A})^{2} u & =\beta u \quad \text { in } \Omega, \\
u & =0 \quad \text { on } \partial \Omega,
\end{aligned}
$$

where $\mathfrak{u}(x)$ is complex-valued and

$$
\vec{A}: \mathbb{R}^{\mathrm{d}} \rightarrow \mathbb{R}^{\mathrm{d}}
$$

is a given bounded vector field.
Physically, $\vec{A}$ represents the vector potential, whose curl equals the magnetic field: $\nabla \times \vec{A}=\vec{B}$. Note that in 2 dimensions, one extends $\vec{A}=$
$\left(A_{1}, A_{2}\right)$ to a 3-vector $\left(A_{1}, A_{2}, 0\right)$ before taking the curl, so that the field $\vec{B}=\left(0,0, \frac{\partial A_{2}}{\partial x_{1}}-\frac{\partial A_{1}}{\partial x_{2}}\right)$ cuts vertically through the plane of the domain. For a brief explanation of how the magnetic Laplacian arises from the correspondence between classical energy functions and quantum mechanical Hamiltonians, see [ReedSimon2, p. 173].

Now we choose the Hilbert spaces and sesquilinear form. Consider only the Dirichlet boundary condition, for simplicity:
$\mathcal{H}=\mathrm{L}^{2}(\Omega ; \mathbb{C})$ (complex valued functions), with inner product

$$
\langle u, v\rangle_{\mathrm{L}^{2}}=\int_{\Omega} u \bar{v} \mathrm{dx}
$$

$\mathcal{K}=\mathrm{H}_{0}^{1}(\Omega ; \mathbb{C})$ with inner product

$$
\langle u, v\rangle_{\mathrm{H}^{1}}=\int_{\Omega}[\nabla u \cdot \overline{\nabla v}+u \bar{v}] \mathrm{d} x
$$

Density: $\mathcal{K}$ contains $\mathrm{C}_{0}^{\infty}$, which is dense in $\mathrm{L}^{2}$.
Continuous imbedding $\mathrm{H}_{0}^{1} \hookrightarrow \mathrm{~L}^{2}$ is trivial, since $\|\mathfrak{u}\|_{\mathrm{L}^{2}} \leq\|\mathfrak{u}\|_{\mathrm{H}^{1}}$, and the imbedding is compact by the Rellich-Kondrachov Theorem [GilbargTrudinger, Theorem 7.22].

Sesquilinear form: define

$$
\mathrm{a}(\mathrm{u}, v)=\int_{\Omega}(i \nabla+\vec{A}) u \cdot \overline{(i \nabla+\vec{A}) v} \mathrm{~d} x+\mathrm{C} \int_{\Omega} u \bar{v} \mathrm{~d} x, \quad u, v \in \mathrm{H}_{0}^{1}(\Omega ; \mathbb{C})
$$

with constant $C=\|\vec{A}\|_{L^{\infty}}^{2}+\frac{1}{2}$. Clearly a is symmetric and continuous on $H_{0}^{1}$. Ellipticity:

$$
\begin{aligned}
\mathrm{a}(\mathrm{u}, \mathfrak{u}) & =\int_{\Omega}\left[|\nabla \mathfrak{u}|^{2}+2 \operatorname{Re}(\mathfrak{i} \nabla \mathfrak{u} \cdot \vec{A} \bar{u})+|\vec{A}|^{2}|\mathfrak{u}|^{2}+C|u|^{2}\right] \mathrm{d} x \\
& \geq \int_{\Omega}\left[|\nabla \mathfrak{u}|^{2}-2|\nabla \mathfrak{u}||\vec{A} \| \mathfrak{u}|+2|\vec{A}|^{2}|\mathfrak{u}|^{2}+\frac{1}{2}|\mathfrak{u}|^{2}\right] \mathrm{d} x \\
& \geq \int_{\Omega}\left[\frac{1}{2}|\nabla u|^{2}+\frac{1}{2}|\mathfrak{u}|^{2}\right] \mathrm{d} x \\
& =\frac{1}{2}\|\mathfrak{u}\|_{\mathrm{H}^{1}}^{2}
\end{aligned}
$$

The Spectral Theorem 4.1 gives an ONB $\left\{\mathbf{u}_{j}\right\}$ for $L^{2}(\Omega ; \mathbb{C})$ and corresponding eigenvalues which we denote $\gamma_{j}=\beta_{j}+C>0$ satisfying

$$
\int_{\Omega}(i \nabla+\vec{A}) u_{j} \cdot \overline{(i \nabla+\vec{A}) v)} d x=\beta_{j} \int_{\Omega} u_{j} \bar{v} d x \quad \forall v \in H_{0}^{1}(\Omega ; \mathbb{C})
$$

In particular,

$$
(i \nabla+\vec{A})^{2} u_{j}=\beta_{j} u_{j}
$$

weakly (and hence classically, assuming smoothness of the vector potential $\vec{A}$ ), so that $u_{j}$ is an eigenfunction of the magnetic Laplacian $(i \nabla+\vec{A})^{2}$ with eigenvalue $\beta_{j}$. We have

$$
\beta_{1} \leq \beta_{2} \leq \beta_{3} \leq \cdots \rightarrow \infty
$$

The eigenvalues satisfy

$$
\beta_{j}=\frac{\int_{\Omega}\left|(i \nabla+\vec{A}) u_{j}\right|^{2} d x}{\int_{\Omega}\left|u_{j}\right|^{2} d x},
$$

as we see by choosing $v=\mathfrak{u}_{\mathrm{j}}$ in the weak formulation. Hence the eigenvalues are all nonnegative.

In fact $\beta_{1}>0$ if the magnetic field vanishes nowhere, as we will show by proving the contrapositive. If $\beta_{1}=0$ then $(i \nabla+\vec{A}) u_{1} \equiv 0$, which implies $\vec{A}=-\mathfrak{i} \nabla \log u_{1}$ wherever $u_{1}$ is nonzero. Then $\nabla \times \vec{A}=0$ wherever $u_{1}$ is nonzero, since the curl of a gradient vanishes identically. (Here we assume $\mathfrak{u}_{1}$ is twice continuously differentiable.) Thus the magnetic field vanishes somewhere, as we wanted to show.

Aside. The preceding argument works regardless of the boundary condition. In the case of Dirichlet boundary conditions, one need not assume the magnetic field is nonvanishing, because the above argument and the reality of $\stackrel{\rightharpoonup}{A}$ together imply that if $\beta_{1}=0$ then $\left|\mathfrak{u}_{1}\right|$ is constant, which is impossible since $u_{1}=0$ on the boundary.

## Gauge invariance

Many different vector potentials can generate the same magnetic field. For example, in 2 dimensions the potentials

$$
\vec{A}=\left(0, x_{1}\right), \quad \vec{A}=\left(-x_{2}, 0\right), \quad \vec{A}=\frac{1}{2}\left(-x_{2}, x_{1}\right)
$$

all generate the same (constant) magnetic field: $\nabla \times \vec{A}=(0,0,1)$. Indeed, adding any gradient vector $\nabla \mathrm{f}$ to the potential leaves the magnetic field unchanged, since the curl of a gradient equals zero. This phenomenon goes by the name of gauge invariance.

How is the spectral theory of the magnetic Laplacian affected by gauge invariance? The sesquilinear form definitely changes when we replace $\vec{A}$ with $\vec{A}+\nabla f$. Fortunately, the new eigenfunctions are related to the old by a unitary transformation, as follows. Suppose $f$ is $C^{1}$-smooth on the closure of the domain. For any trial function $u \in H_{0}^{1}(\Omega ; \mathbb{C})$ we note that the modulated function $e^{\text {if }} \mathfrak{u}$ also belongs to $\mathrm{H}_{0}^{1}(\Omega ; \mathbb{C})$, and that

$$
(i \nabla+\vec{A}) u=(i \nabla+\vec{A}+\nabla f)\left(e^{i f} u\right)
$$

Thus if we write a for the original sesquilinear form and $\widetilde{\mathfrak{a}}$ for the analogous form coming from the vector potential $\vec{A}+\nabla \mathrm{f}$, we deduce

$$
\mathrm{a}(u, v)=\widetilde{\mathfrak{a}}\left(e^{\text {if }} u, e^{\text {if }} v\right)
$$

for all trial functions $u, v$. Since also $\langle u, v\rangle_{\mathrm{L}^{2}}=\left\langle e^{i f} u, e^{i f} v\right\rangle_{\mathrm{L}^{2}}$, we find that the ONB of eigenfunctions $u_{j}$ associated with a transforms to an ONB of eigenfunctions $e^{i f} u_{j}$ associated with $\widetilde{\mathbf{a}}$. The eigenvalues (energy levels) $\beta_{j}$ are unchanged by this transformation.

For geometric invariance of the spectrum with respect to rotations, reflections and translations, and for a discussion of the Neumann and Robin situations, see [LaugesenLiangRoy, Appendix A].

## Higher dimensions

In dimensions $d \geq 4$ we identify the vector potential $\vec{A}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ with a 1-form

$$
A=A_{1} d x_{1}+\cdots+A_{d} d x_{d}
$$

and obtain the magnetic field from the exterior derivative:

$$
\mathrm{B}=\mathrm{d} \mathrm{~A} .
$$

Apart from that, the spectral theory proceeds as in dimensions 2 and 3.

## Chapter 8

## Application: ONB of eigenfunctions for Schrödinger in a confining well

## Goal

To apply the spectral theorem from Chapter 4 to the harmonic oscillator and more general confining potentials in higher dimensions.

Schrödinger operator with potential growing to infinity
We treat a locally bounded, real-valued potential $\mathrm{V}(\mathrm{x})$ on $\mathbb{R}^{\mathrm{d}}$ that grows at infinity:

$$
-\mathrm{C} \leq \mathrm{V}(\mathrm{x}) \rightarrow \infty \quad \text { as }|\mathrm{x}| \rightarrow \infty,
$$

for some constant $C>0$. For example, $\mathrm{V}(\mathrm{x})=|\mathrm{x}|^{2}$ gives the harmonic oscillator.

We aim to prove existence of an ONB of eigenfunctions and eigenvalues for

$$
\begin{aligned}
(-\Delta+\mathrm{V}) \mathbf{u} & =\mathrm{Eu} & & \text { in } \mathbb{R}^{\mathrm{d}} \\
\mathfrak{u} & \rightarrow 0 & & \text { as }|x| \rightarrow \infty
\end{aligned}
$$

$$
\begin{aligned}
& \Omega=\mathbb{R}^{\mathrm{d}} \\
& \mathcal{H}=\mathrm{L}^{2}\left(\mathbb{R}^{\mathrm{d}}\right), \text { inner product }\langle u, v\rangle_{\mathrm{L}^{2}}=\int_{\mathbb{R}^{\mathrm{d}}} u v \mathrm{~d} x .
\end{aligned}
$$

$\mathcal{K}=\mathrm{H}^{1}\left(\mathbb{R}^{\mathrm{d}}\right) \cap \mathrm{L}^{2}(|\mathrm{~V}| \mathrm{dx})$ under the inner product

$$
\langle u, v\rangle_{\mathcal{K}}=\int_{\mathbb{R}^{\mathrm{d}}}[\nabla \mathrm{u} \cdot \nabla v+(1+|\mathrm{V}|) \mathfrak{u v}] \mathrm{d} x
$$

Density: $\mathcal{K}$ contains $C_{0}^{\infty}$, which is dense in $L^{2}$.
Continuous imbedding $\mathcal{K} \hookrightarrow \mathrm{L}^{2}$ is trivial, since $\|\mathfrak{u}\|_{\mathrm{L}^{2}} \leq\|\mathfrak{u}\|_{\mathcal{K}}$. To prove the imbedding is compact:

Proof that imbedding is compact. Suppose $\left\{\mathrm{f}_{\mathrm{k}}\right\}$ is a bounded sequence in $\mathcal{K}$, say with $\left\|f_{k}\right\|_{\mathcal{K}} \leq M$ for all $k$. We must prove the existence of a subsequence converging in $L^{2}\left(\mathbb{R}^{d}\right)$.

The sequence is bounded in $H^{1}(B(R))$ for each ball $B(R) \subset \mathbb{R}^{d}$ that is centered at the origin. Take $R=1$. The Rellich-Kondrachov theorem provides a subsequence that converges in $L^{2}(B(1))$. Repeating with $R=2$ provides a sub-subsequence converging in $L^{2}(B(2))$. Continue in this fashion and then consider the diagonal subsequence, to obtain a subsequence that converges in $L^{2}(B(R))$ for each $R>0$.

We will show this subsequence converges in $L^{2}\left(\mathbb{R}^{d}\right)$. Denote it by $\left\{f_{\mathrm{k}_{\ell}}\right\}$. Let $\varepsilon>0$. Since $V(x)$ grows to infinity, we may choose $R$ so large that $V(x) \geq 1 / \varepsilon$ when $|x| \geq R$. Then

$$
\begin{aligned}
\int_{\mathbb{R}^{d} \backslash B(R)} f_{k_{\ell}}^{2} d x & \leq \varepsilon \int_{\mathbb{R}^{d} \backslash B(R)} f_{k_{\ell}}^{2} V d x \\
& \leq \varepsilon\left\|f_{k_{\ell}}\right\|_{\mathcal{K}}^{2} \\
& \leq \varepsilon M^{2}
\end{aligned}
$$

for all $\ell$. Since also $\left\{f_{k_{\ell}}\right\}$ converges on $B(R)$, we have

$$
\limsup _{\ell, m \rightarrow \infty}\left\|f_{\mathrm{k}_{\ell}}-f_{\mathrm{k}_{\mathrm{m}}}\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{d}\right)}=\limsup _{\ell, \mathrm{m} \rightarrow \infty}\left\|\mathrm{f}_{\mathrm{k}_{\ell}}-\mathrm{f}_{\mathrm{k}_{\mathrm{m}}}\right\|_{\mathrm{L}^{2}\left(\mathbb{R}^{d} \backslash \mathrm{~B}(\mathrm{R})\right)} \leq 2 \sqrt{\varepsilon} M .
$$

Therefore $\left\{f_{k_{\ell}}\right\}$ is Cauchy in $L^{2}\left(\mathbb{R}^{d}\right)$, and hence converges.
Sesquilinear form: define

$$
\mathrm{a}(\mathrm{u}, v)=\int_{\mathbb{R}^{\mathrm{d}}}[\nabla \mathrm{u} \cdot \nabla v+\mathrm{vuv}] \mathrm{d} x+(2 \mathrm{C}+1) \int_{\mathbb{R}^{\mathrm{d}}} u v \mathrm{dx}, \quad \mathrm{u}, v \in \mathcal{K} .
$$

Clearly $a$ is symmetric and continuous on $\mathcal{K}$.
Ellipticity: $a(u, u) \geq\|u\|_{\mathcal{K}}^{2}$, since $V+2 C+1 \geq 1+|V|$.

The Spectral Theorem 4.1 gives an ONB $\left\{\mathbf{u}_{j}\right\}$ for $L^{2}\left(\mathbb{R}^{d}\right)$ and corresponding eigenvalues which we denote $\gamma_{j}=E_{j}+2 C+1>0$ satisfying

$$
\int_{\mathbb{R}^{\mathrm{d}}}\left[\nabla \mathrm{u}_{\mathrm{j}} \cdot \nabla v+\mathrm{vu}_{\mathrm{j}} v\right] \mathrm{d} x=\mathrm{E}_{\mathrm{j}} \int_{\mathbb{R}^{\mathrm{d}}} \mathrm{u}_{\mathrm{j}} v \mathrm{dx} \quad \forall v \in \mathcal{K} .
$$

In particular,

$$
-\Delta u_{j}+V u_{j}=E_{j} u_{j}
$$

weakly (and hence classically, assuming smoothness of V ), so that $\mathfrak{u}_{j}$ is an eigenfunction of the Schrödinger operator $-\Delta+V$, with eigenvalue $E_{j}$. We have

$$
\mathrm{E}_{1} \leq \mathrm{E}_{2} \leq \mathrm{E}_{3} \leq \cdots \rightarrow \infty .
$$

The boundary condition $u_{j} \rightarrow 0$ at infinity is interpreted to mean, more precisely, that $u_{j}$ belongs to the space $H^{1}\left(\mathbb{R}^{d}\right) \cap L^{2}(|V| d x)$. This condition suffices to rule out the existence of any other eigenvalues for the harmonic oscillator, for example, as one can show by direct estimation [Strauss].

The eigenvalues satisfy

$$
\mathrm{E}_{\mathrm{j}}=\frac{\int_{\mathbb{R}^{\mathrm{d}}}\left(\left|\nabla \mathrm{u}_{\mathrm{j}}\right|^{2}+\mathrm{V} \mathrm{u}_{\mathrm{j}}^{2}\right) \mathrm{d} x}{\int_{\mathbb{R}^{\mathrm{d}}} u_{\mathrm{j}}^{2} \mathrm{~d} x}
$$

as we see by choosing $v=u_{j}$ in the weak formulation. Hence if $V \geq 0$ then the eigenvalues are all positive.

## Chapter 9

## Variational characterizations of eigenvalues

## Goal

To obtain minimax and maximin characterizations of the eigenvalues of the sesquilinear form in Chapter 4.

References [Bandle] Section III.1.2

Motivation and hypotheses. How can one estimate the eigenvalues if the spectrum cannot be computed explicitly? We will develop two complementary variational characterizations of eigenvalues. The intuition for these characterizations comes from the special case of eigenvalues of a Hermitian (or real symmetric) matrix $A$, for which the sesquilinear form is $a(u, v)=A u \cdot \bar{v}$ and the first eigenvalue is

$$
\lambda_{1}=\min _{v \neq 0} \frac{A v \cdot \bar{v}}{v \cdot \bar{v}} .
$$

We will work under the assumptions of the discrete spectral theorem in Chapter 4, for the sesquilinear form $\mathbf{a}$. Recall the ordering

$$
\lambda_{1} \leq \lambda_{2} \leq \lambda_{3} \leq \cdots \rightarrow \infty .
$$

Poincaré's minimax characterization of the eigenvalues

Define the Rayleigh quotient of $\mathfrak{u}$ to be

$$
\frac{\mathfrak{a}(u, u)}{\langle u, u\rangle_{\mathcal{H}}} .
$$

We claim $\lambda_{1}$ equals the minimum value of the Rayleigh quotient:

$$
\begin{equation*}
\lambda_{1}=\min _{f \in \mathcal{K} \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}} \tag{9.1}
\end{equation*}
$$

This characterization of the first eigenvalue is the Rayleigh principle.
More generally, each eigenvalue is given by a minimax formula known as the Poincaré principle:

$$
\begin{equation*}
\lambda_{j}=\min _{S} \max _{f \in S \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}} \tag{9.2}
\end{equation*}
$$

where $S$ ranges over all $\mathfrak{j}$-dimensional subspaces of $\mathcal{K}$.
Remark. The Rayleigh and Poincaré principles provide upper bounds on eigenvalues, since they expresses $\lambda_{j}$ as a minimum. More precisely, we obtain an upper bound on $\lambda_{j}$ by choosing $S$ to be any $\mathfrak{j}$-dimensional subspace and evaluating the maximum of the Rayleigh quotient over $f \in S$.

Proof of Poincaré principle. First we prove the Rayleigh principle for the first eigenvalue. Let $\mathrm{f} \in \mathcal{K}$. Then f can be expanded in terms of the ONB of eigenvectors as

$$
f=\sum_{j} c_{j} u_{j}
$$

where $\mathbf{c}_{\mathfrak{j}}=\left\langle\mathbf{f}, \boldsymbol{u}_{\mathfrak{j}}\right\rangle_{\mathcal{H}}$. This series converges in both $\mathcal{H}$ and $\mathcal{K}$ (as we proved in Chapter 4). Hence we may substitute it into the Rayleigh quotient to obtain

$$
\begin{align*}
\frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}} & =\frac{\sum_{j, k} c_{j} \overline{c_{k}} a\left(u_{j}, u_{k}\right)}{\sum_{j, k} c_{j} \overline{c_{k}}\left\langle u_{j}, u_{k}\right\rangle_{\mathcal{H}}} \\
& =\frac{\sum_{j}\left|c_{j}\right|^{2} \lambda_{j}}{\sum_{j}\left|c_{j}\right|^{2}} \tag{9.3}
\end{align*}
$$

since the eigenvectors $\left\{u_{j}\right\}$ are orthonormal in $\mathcal{H}$ and the collection $\left\{u_{j} / \sqrt{\lambda_{j}}\right\}$ is $\mathfrak{a}$-orthonormal in $\mathcal{K}$ (that is, $\left.\mathfrak{a}\left(\mathfrak{u}_{\mathfrak{j}}, \mathfrak{u}_{\mathrm{k}}\right)=\boldsymbol{\lambda}_{\boldsymbol{j}} \boldsymbol{\delta}_{\mathfrak{j} k}\right)$. The expression (9.3) is
obviously greater than or equal to $\lambda_{1}$, with equality when $f=u_{1}$, and so we have proved the Rayleigh principle (9.1).

Next we prove the minimax formula (9.2) for $\mathfrak{j}=2$. (We leave the case of higher $\mathfrak{j}$-values as an exercise.) Choose $S=\left\{c_{1} u_{1}+c_{2} u_{2}: c_{1}, c_{2}\right.$ scalars $\}$ to be the span of the first two eigenvectors. Then

$$
\max _{f \in S \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}}=\max _{\left(c_{1}, c_{2}\right) \neq(0,0)} \frac{\sum_{j=1}^{2}\left|c_{j}\right|^{2} \lambda_{j}}{\sum_{j=1}^{2}\left|c_{j}\right|^{2}}=\lambda_{2} .
$$

Hence the minimum on the right side of (9.2) is $\leq \lambda_{2}$.
To prove the opposite inequality, consider an arbitrary 2-dimensional subspace $S \subset \mathcal{K}$. This subspace contains a nonzero vector $g$ that is orthogonal to $u_{1}$ (since given a basis $\left\{v_{1}, v_{2}\right\}$ for the subspace, we can find scalars $d_{1}, d_{2}$ not both zero such that $g=d_{1} v_{1}+d_{2} v_{2}$ satisfies $0=d_{1}\left\langle v_{1}, u_{1}\right\rangle_{\mathcal{H}}+d_{2}\left\langle v_{2}, u_{1}\right\rangle_{\mathcal{H}}=$ $\left\langle\boldsymbol{g}, \mathfrak{u}_{1}\right\rangle_{\mathcal{H}}$ ). Then $\mathrm{c}_{1}=0$ in the expansion for g , and so by (9.3),

$$
\frac{a(g, g)}{\langle g, g\rangle_{\mathcal{H}}}=\frac{\sum_{j=2}^{\infty}\left|c_{j}\right|^{2} \lambda_{j}}{\sum_{j=2}^{\infty}\left|c_{j}\right|^{2}} \geq \lambda_{2}
$$

Hence

$$
\max _{f \in S \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}} \geq \frac{a(g, g)}{\langle g, g\rangle_{\mathcal{H}}} \geq \lambda_{2}
$$

which implies that the minimum on the right side of (9.2) is $\geq \lambda_{2}$.
Variational characterization of eigenvalue sums. The sum of the first $n$ eigenvalues has a simple "minimum" characterization, similar to the Rayleigh principle for the first eigenvalue, but now involving pairwise orthogonal trial functions:

$$
\begin{align*}
& \lambda_{1}+\cdots+\lambda_{n}  \tag{9.4}\\
& =\min \left\{\frac{\mathfrak{a}\left(f_{1}, f_{1}\right)}{\left\langle f_{1}, f_{1}\right\rangle_{\mathcal{H}}}+\cdots+\frac{a\left(f_{n}, f_{n}\right)}{\left\langle f_{n}, f_{n}\right\rangle_{\mathcal{H}}}: f_{j} \in \mathcal{K} \backslash\{0\},\left\langle f_{j}, f_{k}\right\rangle_{\mathcal{H}}=0 \text { when } \mathfrak{j} \neq k\right\} .
\end{align*}
$$

See Bandle's book for the proof and related results [Bandle, Section III.1.2].

## Courant's maximin characterization

The eigenvalues are given also by a maximin formula known as the Courant principle:

$$
\begin{equation*}
\lambda_{\mathrm{j}}=\max _{\mathrm{S}} \min _{\mathrm{f} \in \mathrm{~S}^{ \pm} \backslash\{0\}} \frac{\mathrm{a}(\mathrm{f}, \mathrm{f})}{\langle\mathrm{f}, \mathrm{f}\rangle_{\mathcal{H}}} \tag{9.5}
\end{equation*}
$$

where this time $S$ ranges over all $(\mathfrak{j}-1)$-dimensional subspaces of $\mathcal{K}$.
Remark. The Courant principle provide lower bounds on eigenvalues, since it expresses $\lambda_{j}$ as a maximum. The lower bounds are difficult to compute, however, because $S^{\perp}$ is an infinite dimensional space.

Sketch of proof of Courant principle. The Courant principle reduces to Rayleigh's principle when $\mathfrak{j}=1$, since in that case $S$ is the zero subspace and $S^{\perp}=\mathcal{K}$.

Now take $\mathfrak{j}=2$ (we leave the higher values of $\mathfrak{j}$ as an exercise). For the " $\leq$ " direction of the proof, we choose $S$ to be the 1 -dimensional space spanned by the first eigenvector $u_{1}$. Then every $f \in S^{\perp}$ has $c_{1}=\left\langle f, u_{1}\right\rangle_{\mathcal{H}}=0$ and so

$$
\lambda_{2} \leq \min _{f \in S^{\perp} \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}}
$$

by expanding $f=\sum_{j=2}^{\infty} \mathfrak{c}_{j} \mathfrak{u}_{j}$ and computing as in our proof of the Poincaré principle.

For the " $\geq$ " direction of the proof, consider an arbitrary 1-dimensional subspace $S$ of $\mathcal{K}$. Then $S^{\perp}$ contains some vector of the form $f=c_{1} u_{1}+c_{2} u_{2}$ with at least one of $c_{1}$ or $c_{2}$ nonzero. Hence

$$
\min _{f \in S^{ \pm} \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}} \leq \frac{\sum_{j=1}^{2}\left|c_{j}\right|^{2} \lambda_{j}}{\sum_{j=1}^{2}\left|c_{j}\right|^{2}} \leq \lambda_{2}
$$

as desired.

## Eigenvalues as critical values of the Rayleigh quotient

Even if we did not know the existence of an ONB of eigenvectors we could still prove the Rayleigh principle, by the following alternative approach. Define $\lambda^{*}$ to equal the infimum of the Rayleigh quotient:

$$
\lambda^{*}=\inf _{f \in \mathcal{K} \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}} .
$$

We will prove $\lambda^{*}$ is an eigenvalue.
First, choose an infimizing sequence $\left\{\mathrm{f}_{\mathrm{k}}\right\}$ normalized with $\left\|\mathrm{f}_{\mathrm{k}}\right\|_{\mathcal{H}}=1$, so that

$$
\mathrm{a}\left(\mathrm{f}_{\mathrm{k}}, \mathrm{f}_{\mathrm{k}}\right) \rightarrow \lambda^{*}
$$

By weak compactness of closed balls in the Hilbert space $\mathcal{K}$, we may suppose $\mathrm{f}_{\mathrm{k}}$ converges weakly in $\mathcal{K}$ to some $u \in \mathcal{K}$. Hence $\mathrm{f}_{\mathrm{k}}$ also converges weakly in $\mathcal{H}$ to $u$ (because if $F(\cdot)$ is any bounded linear functional on $\mathcal{H}$ then it is also a bounded linear functional on $\mathcal{K}$ ). We may further suppose $f_{k}$ converges in $\mathcal{H}$ to some $v \in \mathcal{H}$ (by compactness of the imbedding $\mathcal{K} \hookrightarrow \mathcal{H}$ ) and then $\mathrm{f}_{\mathrm{k}}$ converges weakly in $\mathcal{H}$ to $v$, which forces $v=u$. To summarize: $f_{k} \rightharpoonup u$ weakly in $\mathcal{K}$ and $f_{k} \rightarrow \mathfrak{u}$ in $\mathcal{H}$. In particular, $\|\mathfrak{u}\|_{\mathcal{H}}=1$. Therefore we have

$$
\begin{aligned}
0 & \leq a\left(f_{k}-u, f_{k}-u\right) \\
& =a\left(f_{k}, f_{k}\right)-2 \operatorname{Re} a\left(f_{k}, u\right)+a(u, u) \\
& \rightarrow \lambda^{*}-2 \operatorname{Re} a(u, u)+a(u, u) \quad \text { using weak convergence } f_{k} \rightharpoonup u \\
& =\lambda^{*}-a(u, u) \\
& \leq 0
\end{aligned}
$$

by definition of $\lambda^{*}$ as an infimum.
We have shown that the infimum defining $\lambda^{*}$ is actually a minimum,

$$
\lambda^{*}=\min _{f \in \mathcal{K} \backslash\{0\}} \frac{a(f, f)}{\langle f, f\rangle_{\mathcal{H}}},
$$

and that the minimum is attained when $f=u$.
Our second task is to show $\mathfrak{u}$ is an eigenvector with eigenvalue $\lambda^{*}$. Let $\nu \in \mathcal{K}$ be arbitrary and use $\mathrm{f}=u+\varepsilon v$ as a trial function in the Rayleigh quotient; since $u$ gives the minimizer, the derivative at $\varepsilon=0$ must equal zero by the first derivative test from calculus:

$$
0=\left.\frac{d}{d \varepsilon} \frac{a(u+\varepsilon v, u+\varepsilon v)}{\langle u+\varepsilon v, u+\varepsilon v\rangle_{\mathcal{H}}}\right|_{\varepsilon=0}=2 \operatorname{Re} a(u, v)-\lambda^{*} 2 \operatorname{Re}\langle u, v\rangle_{\mathcal{H}} .
$$

The same equation holds with Im instead of Re, as we see by replacing $v$ with $i v$. (This last step is unnecessary when working with real Hilbert spaces, of course.) Hence

$$
a(u, v)=\lambda^{*}\langle u, v\rangle_{\mathcal{H}} \quad \forall v \in \mathcal{K},
$$

which means $\mathfrak{u}$ is an eigenvector for the sesquilinear form $\mathfrak{a}$, with eigenvalue $\lambda^{*}$.

Aside. The higher eigenvalues $\left(\lambda_{j}\right.$ for $\left.\mathfrak{j}>1\right)$ can be obtained by a similar process, minimizing the Rayleigh quotient on the orthogonal complement of the span of the preceding eigenfunctions $\left(u_{1}, \ldots, \mathfrak{u}_{j-1}\right)$.

## Chapter 10

## Monotonicity properties of eigenvalues

## Goal

To apply Poincaré's minimax principle to the Laplacian and related operators, and hence to establish monotonicity results for Dirichlet and Neumann eigenvalues of the Laplacian, and a diamagnetic comparison for the magnetic Laplacian.

References [Bandle]

## Laplacian, biLaplacian, and Schrödinger operators

Applying the Rayleigh principle (9.1) to the examples in Chapters 5-8 gives:

$$
\begin{array}{ll}
\lambda_{1}=\min _{f \in H_{0}^{1}(\Omega)} \frac{\int_{\Omega}|\nabla f|^{2} d x}{\int_{\Omega} f^{2} d x} & \text { Dirichlet Laplacian on } \Omega \\
\rho_{1}=\min _{f \in H^{1}(\Omega)} \frac{\int_{\Omega}|\nabla f|^{2} d x+\sigma \int_{\partial \Omega} f^{2} d S}{\int_{\Omega} f^{2} d x} & \text { Robin Laplacian on } \Omega \\
\mu_{1}=\min _{f \in H^{1}(\Omega)} \frac{\int_{\Omega}|\nabla f|^{2} d x}{\int_{\Omega} f^{2} d x} & \text { Neumann Laplacian on } \Omega
\end{array}
$$

$$
\begin{array}{rlrl}
\Lambda_{1} & =\min _{f \in H_{0}^{2}(\Omega)} \frac{\int_{\Omega} \sum_{m, n=1}^{d} f_{x_{m} x_{n}}^{2} d x}{\int_{\Omega} f^{2} d x} & & \text { Dirichlet biLaplacian on } \Omega \\
& =\min _{f \in H_{0}^{2}(\Omega)} \frac{\int_{\Omega}(\Delta f)^{2} d x}{\int_{\Omega} f^{2} d x}, & & \\
\beta_{1} & =\min _{f \in H_{0}^{1}(\Omega)} \frac{\int_{\mathbb{R}^{d}}|i \nabla f+\vec{A} f|^{2} d x}{\int_{\mathbb{R}^{d}}|f|^{2} d x} & \text { magnetic Laplacian } \\
E_{1} & =\min _{f \in H^{1}\left(\mathbb{R}^{\mathrm{d}}\right) \cap L^{2}(|V| d x)} \frac{\int_{\mathbb{R}^{d}}\left(|\nabla f|^{2}+V f^{2}\right) d x}{\int_{\mathbb{R}^{d}} f^{2} d x} & \text { Schrödinger with potential } \\
& V(x) \text { growing to infinity. }
\end{array}
$$

The Poincaré principle applies too, giving formulas for the higher eigenvalues and hence implying certain monotonicity relations, as follows.

## Neumann $\leq$ Robin $\leq$ Dirichlet

Free membranes give lower tones than partially free and fixed membranes:
Theorem 10.1 (Neumann-Robin-Dirichlet comparison). Let $\Omega$ be a bounded domain in $\mathbb{R}^{\mathrm{d}}$ with Lipschitz boundary, and fix $\sigma>0$.

Then the Neumann eigenvalues of the Laplacian lie below their Robin counterparts, which in turn lie below the Dirichlet eigenvalues:

$$
\mu_{j} \leq \rho_{j} \leq \lambda_{j} \quad \forall j \geq 1
$$

Proof. Poincaré's minimax principle gives the formulas

$$
\begin{aligned}
& \mu_{\mathrm{j}}=\min _{\mathrm{U}} \max _{\mathrm{f} \in \mathrm{U} \backslash\{0\}} \frac{\int_{\Omega}|\nabla \mathrm{f}|^{2} \mathrm{dx}}{\int_{\Omega} \mathrm{f}^{2} \mathrm{dx}} \\
& \rho_{\mathrm{j}}=\min _{\mathrm{T}} \max _{\mathrm{f} \in \mathrm{~T} \backslash\{0\}} \frac{\int_{\Omega}|\nabla \mathrm{f}|^{2} \mathrm{dx}+\sigma \int_{\partial \Omega} f^{2} \mathrm{dS}}{\int_{\Omega} \mathrm{f}^{2} \mathrm{dx}} \\
& \lambda_{\mathrm{j}}=\min _{\mathrm{S}} \max _{\mathrm{f} \in S \backslash\{0\}} \frac{\int_{\Omega}|\nabla \mathrm{f}|^{2} \mathrm{dx}}{\int_{\Omega} \mathrm{f}^{2} d x}
\end{aligned}
$$

where $S$ ranges over all $j$-dimensional subspaces of $H_{0}^{1}(\Omega)$, and $T$ and $U$ range over all $\mathfrak{j}$-dimensional subspaces of $\mathrm{H}^{1}(\Omega)$.

Clearly $\mu_{j} \leq \rho_{j}$. Further, every subspace $S$ is also a valid $T$, since $H_{0}^{1} \subset H^{1}$. Thus the minimum for $\rho_{j}$ is taken over a larger class of subspaces. Since for $f \in H_{0}^{1}$ the boundary term vanishes in the Rayleigh quotient for $\rho_{j}$, we conclude that $\rho_{j} \leq \lambda_{j}$.

## Domain monotonicity for Dirichlet spectrum

Making a drum smaller increases its frequencies of vibration:
Theorem 10.2. Let $\Omega$ and $\widetilde{\Omega}$ be bounded domains in $\mathbb{R}^{\mathrm{d}}$, and denote the eigenvalues of the Dirichlet Laplacian on these domains by $\lambda_{j}$ and $\widetilde{\lambda}_{j}$, respectively.

If $\Omega \supset \widetilde{\Omega}$ then

$$
\lambda_{j} \leq \widetilde{\lambda}_{j} \quad \forall j \geq 1
$$

Proof. Poincaré's minimax principle gives that

$$
\begin{aligned}
& \lambda_{j}=\min _{S} \max _{f \in S \backslash\{0\}} \frac{\int_{\Omega}|\nabla f|^{2} d x}{\int_{\Omega} f^{2} d x} \\
& \widetilde{\lambda}_{j}=\min _{\widetilde{\widetilde{S}}} \max _{f \in \widetilde{S} \backslash\{0\}} \frac{\int_{\widetilde{\Omega}} \mid \nabla f f^{2} d x}{\int_{\widetilde{\Omega}} f^{2} d x}
\end{aligned}
$$

where $S$ ranges over all j-dimensional subspaces of $H_{0}^{1}(\Omega)$ and $\widetilde{S}$ ranges over all j-dimensional subspaces of $\mathrm{H}_{0}^{1}(\widetilde{\Omega})$.

Every subspace $\widetilde{S}$ is also a valid $S$, since $H_{0}^{1}(\widetilde{\Omega}) \subset H_{0}^{1}(\Omega)$ (noting that any approximating function in $\mathrm{C}_{0}^{\infty}(\widetilde{\Omega})$ belongs also to $\mathrm{C}_{0}^{\infty}(\Omega)$ by extension by 0.) Therefore $\lambda_{j} \leq \widetilde{\lambda}_{j}$.

## Restricted reverse monotonicity for Neumann spectrum

The monotonicity proof breaks down in the Neumann case because $H^{1}(\widetilde{\Omega})$ is not a subspace of $H^{1}(\Omega)$. More precisely, while one can extend a function in $\mathrm{H}^{1}(\widetilde{\Omega})$ to belong to $\mathrm{H}^{1}(\Omega)$, the extended function must generally be nonzero outside $\widetilde{\Omega}$, and so its $L^{2}$ norm and Dirichlet integral will differ from those of the original function.

Furthermore, counterexamples to domain monotonicity are easy to construct for Neumann eigenvalues, as the figure below shows with a rectangle contained in a square. In that example, the square has side length 1 and hence $\mu_{2}=\pi^{2}$, while the rectangle has side length $\sqrt{2}(0.9)$ and so $\widetilde{\mu}_{2}=\pi^{2} /(1.62)$, which is smaller than $\mu_{2}$.

Nonetheless, monotonicity does holds in a certain restricted situation, although the inequality is reversed - the smaller drum has lower tones:


Theorem 10.3. Let $\Omega$ and $\widetilde{\Omega}$ be bounded Lipschitz domains in $\mathbb{R}^{d}$, and denote the eigenvalues of the Neumann Laplacian on these domains by $\mu_{\mathrm{j}}$ and $\widetilde{\mu}_{j}$, respectively.

If $\widetilde{\Omega} \subset \Omega$ and $\Omega \backslash \widetilde{\Omega}$ has measure zero, then

$$
\widetilde{\mu}_{\mathrm{j}} \leq \mu_{\mathrm{j}} \quad \forall \mathrm{j} \geq 1
$$

One might imagine the smaller domain $\widetilde{\Omega}$ as being constructed by removing a hypersurface of measure zero from $\Omega$, thus introducing an additional boundary surface. Reverse monotonicity then makes perfect sense, because the additional boundary, on which values are not specified for the eigenfunctions, enables the eigenfunctions to "relax" and hence lowers the eigenvalues.


Introducing additional boundary surfaces to a Dirichlet problem would have the opposite effect: the eigenfunctions would be further constrained, and the eigenvalues raised.

Proof. Poincaré's minimax principle gives that

$$
\begin{aligned}
& \mu_{j}=\min _{S} \max _{f \in S \backslash\{0\}} \frac{\int_{\Omega}|\nabla f|^{2} d x}{\int_{\Omega} f^{2} d x} \\
& \widetilde{\mu}_{j}=\min _{\widetilde{S}} \max _{f \in \widetilde{S} \backslash\{0\}} \frac{\int_{\Omega}|\nabla f|^{2} d x}{\int_{\Omega} f^{2} d x}
\end{aligned}
$$

where $S$ ranges over all j-dimensional subspaces of $H^{1}(\Omega)$ and $\widetilde{S}$ ranges over all $\mathfrak{j}$-dimensional subspaces of $\mathrm{H}^{1}(\widetilde{\Omega})$.

Every subspace $S$ is also a valid $\widetilde{S}$, since each $f \in H^{1}(\Omega)$ restricts to a function in $\mathrm{H}^{1}(\widetilde{\Omega})$ that has the same $\mathrm{H}^{1}$-norm (using here that $\Omega \backslash \widetilde{\Omega}$ has measure zero). Therefore $\widetilde{\mu}_{j} \leq \mu_{j}$.

## Diamagnetic comparison for the magnetic Laplacian

Imposing a magnetic field always raises the ground state energy.
Theorem 10.4 (Diamagnetic comparison).

$$
\beta_{1} \geq \lambda_{1}
$$

First we prove a pointwise comparison.
Lemma 10.5 (Diamagnetic inequality).

$$
|(i \nabla+\vec{A}) f| \geq|\nabla| f| |
$$

Proof of Lemma 10.5. Write $f$ in polar form as $f=R e^{i \Theta}$. Then

$$
\begin{aligned}
|i \nabla f+\vec{A} f|^{2} & =\left|i e^{i \Theta} \nabla R-R e^{i \Theta} \nabla \Theta+\vec{A} R e^{i \Theta}\right|^{2} \\
& =|i \nabla R-R \nabla \Theta+\vec{A} R|^{2} \\
& =|\nabla R|^{2}+R^{2}|\nabla \Theta-\vec{A}|^{2} \\
& \geq|\nabla R|^{2}=|\nabla| f| |^{2} .
\end{aligned}
$$

Proof of Theorem 10.4. The proof is immediate from the diamagnetic inequality in Lemma 10.5 and the Rayleigh principles for $\beta_{1}$ and $\lambda_{1}$ at the beginning of this chapter. Note we can assume $f \geq 0$ in the Rayleigh principle for $\lambda_{1}$, since the first Dirichlet eigenfunction can be taken nonnegative [GilbargTrudinger, Theorem 8.38].

## Chapter 11

## Weyl's asymptotic for high eigenvalues

## Goal

To determine the rate of growth of eigenvalues of the Laplacian.

References [Arendt]; [CourantHilbert] Section VI. 4

## Notation

The asymptotic notation $\alpha_{j} \sim \beta_{j}$ means

$$
\lim _{j \rightarrow \infty} \frac{\alpha_{j}}{\beta_{j}}=1
$$

Write $\mathrm{V}_{\mathrm{d}}$ for the volume of the unit ball in d-dimensions.

## Growth of eigenvalues

The eigenvalues of the Laplacian grow at a rate $\boldsymbol{c j}^{2 / d}$ where the constant depends only on the volume of the domain, independent of the boundary conditions.

Theorem 11.1 (Weyl's law). Let $\Omega$ be a bounded domain in $\mathbb{R}^{\mathrm{d}}$ with piecewise smooth boundary. As $\mathfrak{j} \rightarrow \infty$ the eigenvalues grow according to:

$$
\lambda_{j} \sim \rho_{j} \sim \mu_{j} \sim \begin{cases}(\pi j /|\Omega|)^{2} & (d=1) \\ 4 \pi j /|\Omega| & (d=2) \\ \left(6 \pi^{2} j /|\Omega|\right)^{2 / 3} & (d=3)\end{cases}
$$

and more generally,

$$
\lambda_{j} \sim \rho_{j} \sim \mu_{j} \sim 4 \pi^{2}\left(\frac{j}{V_{d}|\Omega|}\right)^{2 / d} \quad(d \geq 1)
$$

Here $|\Omega|$ denotes the d -dimensional volume of the domain, in other words its length when $\mathrm{d}=1$ and area when $\mathrm{d}=2$.

In 1 dimension the theorem is proved by the explicit formulas for the eigenvalues in Chapter 2. We will prove the theorem in 2 dimensions, by a technique known as "Dirichlet-Neumann bracketing". The higher dimensional proof is similar.

An alternative proof using small-time heat kernel asymptotics can be found (for example) in the survey paper by Arendt et al. [Arendt, §1.6].

Proof of Weyl aymptotic - Step 1: rectangular domains. In view of the Neu-mann-Robin-Dirichlet comparison (Theorem 10.1), we need only prove Weyl's law for the Neumann and Dirichlet eigenvalues. We provided a proof in Proposition 2.1, for rectangles.

Proof of Weyl aymptotic - Step 2: finite union of rectangles. Next we suppose $R_{1}, \ldots, R_{n}$ are disjoint rectangular domains and put

$$
\begin{aligned}
& \widetilde{\Omega}=\cup_{m=1}^{n} R_{m}, \\
& \Omega=\operatorname{Int}\left(\cup_{m=1}^{n} \overline{R_{m}}\right) .
\end{aligned}
$$



For example, if $R_{1}$ and $R_{2}$ are adjacent squares of side length 1 , then $\widetilde{\Omega}$ is the disjoint union of those squares whereas $\Omega$ is the $2 \times 1$ rectangular domain formed from the interior of their union.

Admittedly $\widetilde{\Omega}$ is not connected, but the spectral theory of the Laplacian remains valid on a finite union of disjoint domains: the eigenfunctions are simply the eigenfunctions of each of the component domains extended to be zero on the other components, and the spectrum equals the union of the spectra of the individual components. (On an infinite union of disjoint domains, on the other hand, one would lose compactness of the imbedding $H^{1} \hookrightarrow L^{2}$, and the zero eigenvalue of the Neumann Laplacian would have infinite multiplicity.)

Write $\widetilde{\lambda}_{j}$ and $\widetilde{\mu}_{j}$ for the Dirichlet and Neumann eigenvalues of $\widetilde{\Omega}$.
Then by the restricted reverse Neumann monotonicity (Theorem 10.3), Neumann-Robin-Dirichlet comparison (Theorem 10.1) and Dirichlet monotonicity (Theorem 10.2), we deduce that

$$
\widetilde{\mu}_{j} \leq \mu_{j} \leq \rho_{j} \leq \lambda_{j} \leq \widetilde{\lambda}_{j} \quad \forall j \geq 1
$$

Hence if we can prove Weyl's law

$$
\begin{equation*}
\widetilde{\mu}_{j} \sim \tilde{\lambda}_{j} \sim \frac{4 \pi j}{|\Omega|} \tag{11.1}
\end{equation*}
$$

for the union-of-rectangles domain $\widetilde{\Omega}$, then Weyl's law will follow for the original domain $\Omega$.

Define the eigenvalue counting functions of the rectangle $R_{m}$ to be

$$
\begin{aligned}
\mathrm{N}_{\text {Neu }}\left(\alpha ; R_{m}\right) & =\#\left\{j \geq 1: \mu_{j}\left(R_{m}\right) \leq \alpha\right\}, \\
\mathrm{N}_{\text {Dir }}\left(\alpha ; R_{m}\right) & =\#\left\{j \geq 1: \lambda_{j}\left(R_{m}\right) \leq \alpha\right\} .
\end{aligned}
$$

We know from Weyl's law for rectangles (Step 1 of the proof above) that

$$
\begin{equation*}
N_{\text {Neu }}\left(\alpha ; R_{m}\right) \sim N_{\text {Dir }}\left(\alpha ; R_{m}\right) \sim \frac{\left|R_{m}\right|}{4 \pi} \alpha \tag{11.2}
\end{equation*}
$$

as $\alpha \rightarrow \infty$.
The spectrum of $\widetilde{\Omega}$ is the union of the spectra of the $R_{m}$, and so (here comes the key step in the proof!) the eigenvalue counting functions of $\widetilde{\Omega}$
equal the sums of the corresponding counting functions of the rectangles:

$$
\begin{aligned}
\mathrm{N}_{\mathrm{Neu}}(\alpha ; \widetilde{\Omega}) & =\sum_{m=1}^{n} \mathrm{~N}_{\mathrm{Neu}}\left(\alpha ; \mathrm{R}_{\mathrm{m}}\right), \\
\mathrm{N}_{\mathrm{Dir}}(\alpha ; \widetilde{\Omega}) & =\sum_{\mathrm{m}=1}^{n} \mathrm{~N}_{\mathrm{Dir}}\left(\alpha ; \mathrm{R}_{\mathrm{m}}\right) .
\end{aligned}
$$

Combining these sums with the asymptotic (11.2) shows that

$$
\mathrm{N}_{\mathrm{Neu}}(\alpha ; \widetilde{\Omega}) \sim\left(\sum_{m=1}^{n} \frac{\left|\mathrm{R}_{\mathrm{m}}\right|}{4 \pi}\right) \alpha=\frac{|\Omega|}{4 \pi} \alpha
$$

and similarly

$$
\mathrm{N}_{\mathrm{Dir}}(\alpha ; \widetilde{\Omega}) \sim \frac{|\Omega|}{4 \pi} \alpha
$$

as $\alpha \rightarrow \infty$. We can invert these last two asymptotic formulas with the help of Lemma 2.2, thus obtaining Weyl's law (11.1) for $\widetilde{\Omega}$.

Proof of Weyl aymptotic - Step 3: approximation of arbitrary domains. Lastly we suppose $\Omega$ is an arbitrary domain with piecewise smooth boundary. The idea is to approximate $\Omega$ with a union-of-rectangles domain such as in Step 2 , such that the volume of the approximating domain is within $\varepsilon$ of the volume of $\Omega$. We refer to the text of Courant and Hilbert for the detailed proof [CourantHilbert, §VI.4.4].

## Chapter 12

## Pólya's conjecture and the Berezin-Li-Yau Theorem

## Goal

To describe Polya's conjecture about Weyl's law, and to state the "tiling domain" and "summed" versions that are known to hold.

References [Kellner, Laptev, Pólya]
Pólya's conjecture
Weyl's law (Theorem 11.1) says that

$$
\lambda_{j} \sim \frac{4 \pi j}{|\Omega|} \sim \mu_{j} \quad \text { as } j \rightarrow \infty
$$

for a bounded plane domain $\Omega$ with piecewise smooth boundary. (We restrict to plane domains, in this chapter, for simplicity.)

Pólya conjectured that these asymptotic formulas hold as inequalities.
Conjecture 12.1 ([Pólya], 1960).

$$
\lambda_{\mathrm{j}} \geq \frac{4 \pi \mathrm{j}}{|\Omega|} \geq \mu_{\mathrm{j}} \quad \forall \mathrm{j} \geq 1
$$

The conjecture remains open even for a disk.
Pólya proved the Dirichlet part of the inequality for tiling domains [Pólya], and Kellner did the same for the Neumann part [Kellner]. Recall that a
"tiling domain" covers the plane with congruent copies of itself (translations, rotations and reflections). For example, parallelograms and triangles are tiling domains, as are many variants of these domains (a fact that M. C. Escher exploited in his artistic creations).

Pólya and Kellner's proofs are remarkably simple, using a rescaling argument together with Weyl's law.

For arbitrary domains, Pólya's conjecture has been proved only for $\lambda_{1}, \lambda_{2}$ (see [Henrot, Th. 3.2.1 and (4.3)]) and for $\mu_{1}, \mu_{2}, \mu_{3}$ (see [Girouard]). The conjecture remains open for $\mathfrak{j} \geq 3$ (Dirichlet) and $\mathfrak{j} \geq 4$ (Neumann).

## Berezin-Li-Yau results

The major progress for arbitrary domains has been on a "summed" version of the conjecture. (Quite often in analysis, summing or integrating an expression produces a significantly more tractable quantity.) Li and Yau [LiYau] proved that

$$
\sum_{k=1}^{j} \lambda_{k} \geq \frac{2 \pi j^{2}}{|\Omega|}
$$

which is only slightly smaller than the quantity $(2 \pi /|\Omega|) \mathfrak{j}(\mathfrak{j}+1)$ that one gets by summing the left side of the Pólya conjecture. An immediate consequence is a Weyl-type inequality for Dirichlet eigenvalues:

$$
\lambda_{j} \geq \frac{2 \pi j}{|\Omega|}
$$

by combining the very rough estimate $j \lambda_{j} \geq \sum_{k=1}^{j} \lambda_{k}$ with the Li-Yau inequality. The last formula has $2 \pi$ whereas Pólya's conjecture demands $4 \pi$, and so we see the conjecture is true up to a factor of 2 , at worst.

Similar results hold for Neumann eigenvalues.
A somewhat more general approach had been obtained earlier by Berezin. For more information, consult the work of Laptev [Laptev] and a list of open problems from recent conferences [AIM].

## Chapter 13

## Case study: stability of steady states for reaction-diffusion PDEs

## Goal

To linearize a nonlinear reaction-diffusion PDE around a steady state, and study the spectral theory of the linearized operator by time-map methods.

References [Schaaf] Section 4.1

## Reaction-diffusion PDEs

Assume throughout this section that $f(y)$ is a smooth function on $\mathbb{R}$. Let $X>0$. We study the reaction-diffusion PDE

$$
\begin{equation*}
u_{t}=u_{x x}+f(u) \tag{13.1}
\end{equation*}
$$

on the interval $(0, X)$ with Dirichlet boundary conditions $u(0)=u(X)=0$. Physical interpretations include: (i) $u=$ temperature and $f=$ rate of heat generation, (ii) $u=$ chemical concentration and $f=$ reaction rate of chemical creation.

Intuitively, the 2 nd order diffusion term in the PDE is stabilizing (since $u_{t}=u_{x x}$ is the usual diffusion equation), whereas the Oth order reaction term can be destabilizing (since solutions to $\mathfrak{u}_{t}=\mathbf{f}(\mathfrak{u})$ will grow, when $\mathbf{f}$ is positive). Thus the reaction-diffusion PDE features a competition between
stabilizing and destabilizing effects. This competition leads to nonconstant steady states, and interesting stability behavior.

Steady states. If $U(x)$ is a steady state, then

$$
\begin{equation*}
\mathrm{U}^{\prime \prime}+\mathrm{f}(\mathrm{U})=0, \quad 0<x<\mathrm{X} \tag{13.2}
\end{equation*}
$$

More than one steady state can exist. For example if $f(0)=0$ then $U \equiv 0$ is a steady state, but nonconstant steady states might exist too, such as $\mathrm{U}(\mathrm{x})=\sin \mathrm{x}$ when $X=\pi$ and $\mathrm{f}(\mathrm{y})=\mathrm{y}$.

## Linearized PDE

We perturb a steady state by considering

$$
u=u+\varepsilon \phi
$$

where the perturbation $\phi(x, t)$ is assumed to satisfy the Dirichlet BC $\phi=0$ at $x=0, L$, for each $t$. Substituting $u$ into the equation (13.1) gives

$$
\begin{aligned}
0+\varepsilon \phi_{\mathrm{t}} & =\left(\mathrm{U}_{x x}+\varepsilon \phi_{x x}\right)+\mathrm{f}(\mathrm{U}+\varepsilon \phi) \\
& =\mathrm{U}_{x x}+\varepsilon \phi_{x x}+\mathrm{f}(\mathrm{U})+\mathrm{f}^{\prime}(\mathrm{U}) \varepsilon \phi+\mathrm{O}\left(\varepsilon^{2}\right)
\end{aligned}
$$

The leading terms, of order $\varepsilon^{0}$, equal zero by the steady state equation for U. We discard terms of order $\varepsilon^{2}$ and higher. The remaining terms, of order $\varepsilon^{1}$, give the linearized equation:

$$
\begin{equation*}
\phi_{\mathrm{t}}=\phi_{x x}+\mathrm{f}^{\prime}(\mathrm{U}) \phi . \tag{13.3}
\end{equation*}
$$

That is,

$$
\phi_{\mathrm{t}}=-\mathrm{L} \phi
$$

where L is the symmetric linear operator

$$
\mathrm{L} w=-w_{x x}-\mathrm{f}^{\prime}(\mathrm{U}) w
$$

Separation of variables gives (formally) solutions of the form

$$
\phi=\sum_{j} c_{j} e^{-\tau_{j} \mathrm{t}} w_{j}(x),
$$

where the eigenvalues $\tau_{j}$ and Dirichlet eigenfunctions $\boldsymbol{w}_{j}$ satisfy

$$
\mathrm{L} w_{j}=\tau_{j} w_{j}
$$

with $w_{j}(0)=w_{j}(X)=0$.
Thus the steady state U of the reaction-diffusion PDE is
linearly unstable if $\tau_{1}<0$
because the perturbation $\phi$ grows to infinity, whereas the steady state is

## linearly stable if $\tau_{1} \geq 0$

because $\phi$ remains bounded in that case.
To make these claims rigorous, we study the spectrum of L.

## Spectrum of L

We take:

$$
\begin{aligned}
& \Omega=(0, X) \\
& \mathcal{H}=\mathrm{L}^{2}(0, X), \text { inner product }\langle u, v\rangle_{\mathrm{L}^{2}}=\int_{0}^{\mathrm{X}} u \mathfrak{u} \mathrm{~d} x \\
& \mathcal{K}=\mathrm{H}_{0}^{1}(0, X) \text {, inner product } \\
& \qquad\langle u, v\rangle_{\mathrm{H}^{1}}=\int_{0}^{\mathrm{X}}\left(u^{\prime} v^{\prime}+u v\right) \mathrm{d} \chi
\end{aligned}
$$

Compact imbedding $\mathrm{H}_{0}^{1} \hookrightarrow \mathrm{~L}^{2}$ by Rellich-Kondrachov Symmetric sesquilinear form

$$
a(u, v)=\int_{0}^{x}\left(u^{\prime} v^{\prime}-f^{\prime}(u) u v+C u v\right) d x
$$

where $C>0$ is chosen larger than $\left\|f^{\prime}\right\|_{L^{\infty}}+1$. Proof of ellipticity:

$$
a(u, u) \geq \int_{0}^{x}\left(\left(u^{\prime}\right)^{2}+u^{2}\right) d x=\|u\|_{H^{1}}^{2}
$$

by choice of C .
The Spectral Theorem 4.1 now yields an ONB of eigenfunctions $\left\{w_{j}\right\}$ with eigenvalues $\gamma_{j}$ such that

$$
\mathrm{a}\left(w_{\mathrm{j}}, v\right)=\gamma_{\mathrm{j}}\left\langle w_{\mathrm{j}}, v\right\rangle_{\mathrm{L}^{2}} \quad \forall v \in \mathrm{H}_{0}^{1}(0, \mathrm{X}) .
$$

Writing $\gamma_{j}=\tau_{j}+C$ we get

$$
\int_{0}^{x}\left(w_{j}^{\prime} v^{\prime}-\mathrm{f}^{\prime}(\mathrm{U}) w_{\mathrm{j}} v\right) \mathrm{d} x=\tau_{\mathrm{j}} \int_{0}^{X} w_{\mathrm{j}} v \mathrm{~d} x \quad \forall v \in \mathrm{H}_{0}^{1}(0, \mathrm{X}) .
$$

These eigenfunctions satisfy $\mathrm{L} w_{j}=\tau_{j} w_{j}$ weakly, and hence also classically.

## Stability of the zero steady state.

Assume $\mathrm{f}(0)=0$, so that $\mathrm{U} \equiv 0$ is a steady state. Its stability is easily determined, as follows.

The linearized operator is $L w=-w^{\prime \prime}-f^{\prime}(0) w$, which on the interval $(0, X)$ has Dirichlet eigenvalues

$$
\tau_{j}=\left(\frac{j \pi}{X}\right)^{2}-f^{\prime}(0)
$$

Thus the zero steady state is linearly unstable if and only if

$$
\left(\frac{\pi}{X}\right)^{2}<f^{\prime}(0) .
$$

Thus we may call the reaction-diffusion PDE "long-wave unstable" when $f^{\prime}(0)>0$, because then the zero steady state is unstable with respect to perturbations of sufficiently long wavelength X. On short intervals, the Dirichlet BCs are strong enough to stabilize the steady state.

## Sufficient conditions for linearized instability of nonconstant steady states

Our first instability criterion is structural, meaning it depends on properties of the reaction function $f$ rather than on properties of the particular steady state U.

Theorem 13.1. Assume the steady state U is nonconstant, and that $\mathrm{f}(0)=$ $0, f^{\prime \prime}(0)=0$ and $f^{\prime \prime \prime}>0$. Then $\tau_{1}<0$.

For example, the theorem shows that nonconstant steady states are unstable when $f(y)=y^{3}$.

Proof. First we collect facts about boundary values, to be used later in the proof when we integrate by parts:

$$
\begin{array}{rlrl}
\mathrm{U} & =0 \text { at } \mathrm{x} & =0, \mathrm{X} & \\
\mathrm{f}(\mathrm{U}) & =0 \text { at } \mathrm{x}=0, \mathrm{X} & & \text { since } \mathrm{f}(0)=0, \\
\mathrm{U}^{\prime \prime} & =0 \text { at } \mathrm{x} & =0, \mathrm{X} & \\
\text { because } \mathrm{U}^{\prime \prime}=-\mathrm{f}(\mathrm{U}) \\
\mathrm{f}^{\prime \prime}(\mathrm{U}) & =0 \text { at } \mathrm{x} & =0, \mathrm{X} & \\
\text { since } \mathrm{f}^{\prime \prime}(0)=0 .
\end{array}
$$

The Rayleigh principle for L says that

$$
\tau_{1}=\min \left\{\frac{\int_{0}^{x}\left(\left(w^{\prime}\right)^{2}-f^{\prime}(\mathrm{U}) w^{2}\right) d x}{\int_{0}^{x} w^{2} d x}: w \in \mathrm{H}_{0}^{1}(0, X)\right\} .
$$

We choose a trial function

$$
w=\mathrm{U}^{\prime \prime}
$$

which is not the zero function, since U is nonconstant. Then the numerator of the Rayleigh quotient for $w$ is

$$
\begin{aligned}
& \int_{0}^{x}\left(\left(U^{\prime \prime \prime}\right)^{2}-f^{\prime}(U)\left(U^{\prime \prime}\right)^{2}\right) d x \\
& =\int_{0}^{x}\left(-U^{\prime \prime \prime \prime}-f^{\prime}(U) U^{\prime \prime}\right) U^{\prime \prime} d x \\
& \text { by parts } \\
& =\int_{0}^{x} f^{\prime \prime}(U)\left(U^{\prime}\right)^{2} U^{\prime \prime} d x
\end{aligned} \quad \text { by the steady state equation (13.2) } \quad \begin{array}{ll}
=\frac{1}{3} \int_{0}^{x} f^{\prime \prime}(U)\left[\left(U^{\prime}\right)^{3}\right]^{\prime} d x & \text { by parts } \\
=-\frac{1}{3} \int_{0}^{x} f^{\prime \prime \prime}(U)\left(U^{\prime}\right)^{4} d x &
\end{array}
$$

since $f^{\prime \prime \prime}>0$ and $U$ is nonconstant. Hence $\tau_{1}<0$, by the Rayleigh principle.

Motivation for the choice of trial function. Our trial function $w=\mathrm{U}^{\prime \prime}$ corresponds to a perturbation $u=U+\varepsilon \phi=U+\varepsilon U^{\prime \prime}$, which tends (when $\varepsilon>0$ ) to push the steady state towards the constant function. The opposite perturbation $(\varepsilon<0)$ would tend to make the solution grow even further away from the constant steady state.

The next instability criterion, rather than being structural, depends on particular properties of the steady state.

Theorem 13.2 ([Schaaf, Proposition 4.1.2]). Assume the nonconstant steady state U changes sign on $(0, X)$. Then $\tau_{1}<0$.

For example, suppose $f(y)=y$ so that the steady state equation is $U^{\prime \prime}+$ $\mathrm{U}=0$. If $\mathrm{X}=2 \pi$ then the steady state $\mathrm{U}=\sin x$ is linearly unstable, by the theorem. Of course, for that example we can compute the spectrum of L exactly: the lowest eigenfunction is $w=\sin (x / 2)$ with eigenvalue $\tau_{1}=$ $\left(\frac{1}{2}\right)^{2}-1<0$.

Proof. If U changes sign then it has a positive local maximum and a negative local minimum in $(0, X)$, recalling that $\mathrm{U}=0$ at the endpoints. Obviously $\mathrm{U}^{\prime}$ must be nonzero at some point between these local extrema, and so there exist points $0<x_{1}<x_{2}<X$ such that

$$
\mathrm{u}^{\prime}\left(\mathrm{x}_{1}\right)=\mathrm{U}^{\prime}\left(\mathrm{x}_{2}\right)=0
$$

and $\mathrm{U}^{\prime} \neq 0$ on $\left(\mathrm{x}_{1}, \mathrm{x}_{2}\right)$. Define a trial function

$$
w= \begin{cases}\mathrm{u}^{\prime} & \text { on }\left(x_{1}, x_{2}\right) \\ 0 & \text { elsewhere }\end{cases}
$$

(We motivate this choice of trial function at the end of the proof.) Then $w$ is piecewise smooth, and is continuous since $w=\mathrm{U}^{\prime}=0$ at $x_{1}$ and $x_{2}$. Therefore $w \in H_{0}^{1}(0, X)$, and $w \not \equiv 0$ since $U^{\prime} \neq 0$ on ( $x_{1}, x_{2}$ ).

The numerator of the Rayleigh quotient for $w$ is

$$
\begin{aligned}
\int_{0}^{x}\left(\left(w^{\prime}\right)^{2}-f^{\prime}(\mathrm{U}) w^{2}\right) d x & =\int_{x_{1}}^{x_{2}}\left(-w^{\prime \prime}-\mathrm{f}^{\prime}(\mathrm{U}) w\right) w d x \quad \text { by parts } \\
& =0
\end{aligned}
$$

since

$$
\begin{equation*}
-w^{\prime \prime}=-\mathrm{U}^{\prime \prime \prime}=(\mathrm{f}(\mathrm{U}))^{\prime}=\mathrm{f}^{\prime}(\mathrm{U}) \mathrm{U}^{\prime}=\mathrm{f}^{\prime}(\mathrm{U}) w \tag{13.4}
\end{equation*}
$$

Hence $\tau_{1} \leq 0$, by the Rayleigh principle for the first eigenvalue.
Suppose $\tau_{1}=0$, so that the Rayleigh quotient of $w$ equals $\tau_{1}$. Then $w$ must be an eigenfunction with eigenvalue $\tau_{1}$ (because substituting $w=$
$\sum_{j} c_{j} w_{j}$ into the Rayleigh quotient would give a value larger than $\tau_{1}$, if $c_{j}$ were nonzero for any term with eigenvalue larger than $\tau_{1}$ ).

Since the eigenfunction $w$ must be smooth, the slopes of $w$ from the left and the right at $x_{2}$ must agree, which means $w^{\prime}\left(x_{2}\right)=w^{\prime}\left(x_{2}+\right)=0$. Thus $w\left(x_{2}\right)=w^{\prime}\left(x_{2}\right)=0$ and $w$ satisfies the second order linear ODE (13.4) on $\left(x_{1}, x_{2}\right)$. Therefore $w \equiv 0$ on ( $x_{1}, x_{2}$ ), by uniqueness, which contradicts our construction of $w$. We conclude $\tau_{1}<0$.

Motivation for the choice of trial function. The steady state equation reads $\mathrm{U}^{\prime \prime}+\mathrm{f}(\mathrm{U})=0$, and differentiating shows that $\mathrm{U}^{\prime}$ lies in the nullspace of the linearized operator L :

$$
\mathrm{LU}^{\prime}=-\left(\mathrm{U}^{\prime}\right)^{\prime \prime}-\mathrm{f}^{\prime}(\mathrm{U}) \mathrm{U}^{\prime}=0 .
$$

In other words, $\mathrm{U}^{\prime}$ is an eigenfunction with eigenvalue O , which almost proves instability (since instability would correspond to a negative eigenvalue). Of course, the eigenfunction $\mathrm{U}^{\prime}$ does not satisfy the Dirichlet boundary conditions at the endpoints, and hence we must restrict to the subinterval ( $x_{1}, x_{2}$ ), in the proof above, in order to obtain a valid trial function.

## Time maps and linearized stability

Next we derive instability criteria that are almost necessary and sufficient. These conditions depend on the time map for a family of steady states.

Parameterize the steady states by their slope at the left endpoint: given $s \neq 0$, write $\mathrm{U}_{\mathrm{s}}(x)$ for the steady state on $\mathbb{R}$ (if it exists) satisfying

$$
\mathrm{U}_{s}(0)=0, \quad \mathrm{U}_{s}^{\prime}(0)=s, \quad \mathrm{U}_{s}(x)=0 \text { for some } x>0 .
$$

Define the time map to give the first point or "time" $x$ at which the steady state hits the axis:

$$
\mathrm{T}(\mathrm{~s})=\min \left\{x>0: \mathrm{U}_{s}(x)=0\right\} .
$$

If $\mathrm{U}_{\mathrm{s}}$ exists for some $s \neq 0$ then it exists for all nonzero $s$-values in a neighborhood, and the time map is smooth on that neighborhood [Schaaf, Proposition 4.1.1]. The time map can be determined numerically by plotting solutions with different initial slopes, as the figures below show. In the first figure the time map is decreasing, whereas in the second it increases.

Monotonicity of the time maps determines stability of the steady state:


Theorem 13.3 ([Schaaf, Proposition 4.1.3]). The steady state $\mathrm{U}_{\mathrm{s}}$ is linearly unstable on the interval $(0, \mathrm{~T}(\mathrm{~s}))$ if $\mathrm{sT}^{\prime}(\mathrm{s})<0$, and is linearly stable if $\mathrm{sT}^{\prime}(\mathrm{s})>0$.
Proof. We begin by differentiating the family of steady states with respect to the parameter $s$, and obtaining some properties of that function. Then we treat the "instability" and "stability" parts of the theorem separately.

Write $s_{0} \neq 0$ for a specific value of $s$, in order to reduce notational confusion. Let $X=T\left(s_{0}\right)$. Define a function

$$
v=\left.\frac{\partial \mathrm{U}_{\mathrm{s}}}{\partial \mathrm{~s}}\right|_{\mathrm{s}=s_{0}}
$$

on $(0, X)$, where we use that $U_{s}(x)$ is jointly smooth in $(x, s)$. Then

$$
\begin{equation*}
v^{\prime \prime}+\mathrm{f}^{\prime}(\mathrm{U}) v=0 \tag{13.5}
\end{equation*}
$$

as one sees by differentiating the steady state equation (13.2) with respect to $s$, and writing U for $\mathrm{U}_{\mathrm{s}_{0}}$.

At the left endpoint we have

$$
v(0)=0, \quad v^{\prime}(0)=1
$$

because $\mathrm{U}_{\mathrm{s}}(0)=0, \mathrm{U}_{\mathrm{s}}^{\prime}(0)=\mathrm{s}$ for all s .
We do not expect $v$ to vanish at the right endpoint, but we can calculate its value there to be

$$
v(X)=s_{0} \mathrm{~T}^{\prime}\left(\mathrm{s}_{0}\right)
$$

as follows. First, differentiating the equation $0=U_{s}(T(s))$ gives that

$$
\begin{aligned}
0 & =\frac{\partial}{\partial \mathrm{s}} \mathrm{U}_{s}(\mathrm{~T}(\mathrm{~s})) \\
& =\frac{\partial \mathrm{U}_{s}}{\partial \mathrm{~s}}(\mathrm{~T}(\mathrm{~s}))+\mathrm{U}_{\mathrm{s}}^{\prime}(\mathrm{T}(\mathrm{~s})) \mathrm{T}^{\prime}(\mathrm{s}) \\
& =v(\mathrm{~T}(\mathrm{~s}))+\mathrm{U}_{\mathrm{s}}^{\prime}(\mathrm{T}(\mathrm{~s})) \mathrm{T}^{\prime}(\mathrm{s})
\end{aligned}
$$

Note the steady state $\mathrm{U}_{\mathrm{s}}$ is symmetric about the midpoint of the interval $(0, \mathrm{~T}(\mathrm{~s}))$ (exercise; use that $\mathrm{U}_{\mathrm{s}}=0$ at both endpoints and that the steady state equation is invariant under $\chi \mapsto-x$, so that steady states must be symmetric about any local maximum point). Thus $\mathrm{U}_{\mathrm{s}}^{\prime}(\mathrm{T}(\mathrm{s}))=-\mathrm{U}_{\mathrm{s}}^{\prime}(0)=$ $-s$, and evaluating the last displayed formula at $s=s_{0}$ then gives that $0=v(\mathrm{X})-\mathrm{s}_{0} \mathrm{~T}^{\prime}\left(\mathrm{s}_{0}\right)$, as we wanted.

Proof of instability. Assume $s_{0} T^{\prime}\left(s_{0}\right)<0$. Then $v(X)<0$. Since $v^{\prime}(0)=1$ we know $v(x)$ is positive for small values of $x$, and so some $x_{2} \in(0, X)$ exists at which $v\left(x_{2}\right)=0$. Define a trial function

$$
w= \begin{cases}v & \text { on }\left[0, x_{2}\right) \\ 0 & \text { elsewhere }\end{cases}
$$

Then $w$ is piecewise smooth, and is continuous since $v=0$ at $x_{2}$. Note $w(0)=0$. Therefore $w \in \mathrm{H}_{0}^{1}(0, X)$, and $w \not \equiv 0$.

Hence $\tau_{1}<0$ by arguing as in the proof of Theorem 13.2, except with $\mathrm{x}_{1}=0$.
[Motivation for the choice of trial function. Differentiating the steady state equation $\mathrm{U}^{\prime \prime}+\mathrm{f}(\mathrm{U})=0$ with respect to s shows that $\partial \mathrm{U} / \partial \mathrm{s}$ is an eigenfunction with eigenvalue zero:

$$
\mathrm{L}\left(\frac{\partial \mathrm{U}}{\partial \mathrm{~s}}\right)=-\left(\frac{\partial \mathrm{U}}{\partial s}\right)^{\prime \prime}-\mathrm{f}^{\prime}(\mathrm{U}) \frac{\partial \mathrm{U}}{\partial \mathrm{~s}}=0 .
$$

In other words, $\partial \mathrm{U} / \partial \mathrm{s}$ lies in the nullspace of the linearized operator. It does not satisfy the Dirichlet boundary condition at the right endpoint, but we handled that issue in the proof above by restricting to the subinterval ( $0, x_{2}$ ), in order to obtain a valid trial function.]

Proof of stability. Assume $\mathrm{s}_{0} \mathrm{~T}^{\prime}\left(\mathrm{s}_{0}\right)>0$, so that $v(\mathrm{X})>0$. Define $\sigma=$ $-v^{\prime}(\mathrm{X}) / v(\mathrm{X})$. Then

$$
v(0)=0, \quad v^{\prime}(X)+\sigma v(X)=0
$$

which is a mixed Dirichlet-Robin boundary condition. We will show later that $v$ is a first eigenfunction for $L$, under this mixed condition, with eigenvalue is $\rho_{1}=0$ (since $L v=0$ by (13.5)).

By adapting our Dirichlet-to-Robin monotonicity result (Theorem 10.1) one deduces that

$$
\tau_{1} \geq \rho_{1}=0
$$

which gives linearized stability of the steady state U .
To show $v$ is a first eigenfunction for L , as used above, we first show $v$ is positive on ( $0, X$ ). Apply the steady state equation (13.2) to $\mathrm{U}_{s}$, and multiply by $\mathrm{U}_{\mathrm{s}}^{\prime}$ and integrate to obtain the energy equation

$$
\begin{equation*}
\frac{1}{2}\left(\mathrm{U}_{\mathrm{s}}^{\prime}\right)^{2}+\mathrm{F}\left(\mathrm{U}_{\mathrm{s}}\right)=\frac{1}{2} \mathrm{~s}^{2} \tag{13.6}
\end{equation*}
$$

where $F$ is an antiderivative of $f$ chosen with $F(0)=0$. Differentiating with respect to $s$ at $s=s_{0}$ gives that

$$
\mathrm{U}^{\prime} v^{\prime}+\mathrm{f}(\mathrm{U}) v=\mathrm{s}_{0} .
$$

Hence if $v$ vanishes at some $x_{0} \in(0, X)$ then $\mathrm{U}^{\prime} v^{\prime}=s_{0} \neq 0$ and so $v^{\prime}\left(x_{0}\right) \neq$ 0 . Thus at any two successive zeros of $v$, we know $v^{\prime}$ has opposite signs. Therefore $\mathrm{U}^{\prime}$ has opposite signs too, because $\mathrm{U}^{\prime} v^{\prime}=s_{0}$ at the zeros. It is straightforward to show from (13.6) that U increases on $[0, \mathrm{X} / 2$ ] and decreases on $[X / 2, X]$, and so after the zero of $v$ at $x=0$ the next zero (if it exists) can only be $>X / 2$, and the one after that must be $>X$. Since we know $v(x)$ is positive for small $X$ and that $v(X)>0$, we conclude $v$ has no zeros in $(0, X)$ and hence is positive there.

The first eigenfunction of L with mixed Dirichlet-Robin boundary condition is positive, and it is the unique positive eigenfunction (adapt the argument in [GilbargTrudinger, Theorem 8.38]). Since the eigenfunction $v$ is positive, we conclude that it is the first Dirichlet-Robin eigenfunction, as desired.

## Chapter 14

## Case study: stability of steady states for thin fluid film PDEs

## Goal

To linearize a particular nonlinear PDE around a steady state, and develop the spectral theory of the linearized operator.

References [LaugesenPugh1, LaugesenPugh2]

## Thin fluid film PDE

The evolution of a thin layer of fluid (such as paint) on a flat substrate (such as the ceiling) can be modeled using the thin fluid film PDE:

$$
h_{t}=-\left(f(h) h_{x x x}\right)_{x}-\left(g(h) h_{x}\right)_{x}
$$

where $h(x, t)>0$ measures the thickness of the fluid, and the smooth, positive coefficient functions $f$ and $g$ represent surface tension and gravitational effects (or substrate-fluid interactions). For simplicity we assume $f \equiv 1$, so that the equation becomes

$$
\begin{equation*}
h_{t}=-h_{x x x x}-\left(g(h) h_{x}\right)_{x} . \tag{14.1}
\end{equation*}
$$

We will treat the case of general g , but readers are welcome to focus on the special case $g(y)=y^{p}$ for some $p \in \mathbb{R}$.

Solutions are known to exist for small time, given positive smooth initial data. But films can "rupture" in finite time, meaning $h(x, t) \searrow 0$ as $t \nearrow T$, for some coefficient functions $g$ (for example, for $g \equiv 1$ ).

Intuitively, the 4th order "surface tension" term in the PDE is stabilizing (since $h_{t}=-h_{x x x x}$ is the usual 4th order diffusion equation) whereas the 2nd order "gravity" term is destabilizing (since $h_{t}=-h_{x x}$ is the backwards heat equation). Thus the thin film PDE features a competition between stabilizing and destabilizing effects. This competition leads to nonconstant steady states, and interesting stability behavior.

Periodic BCs and conservation of fluid. Fix $X>0$ and assume $h$ is $X$-periodic with respect to $x$. Then the total volume of fluid is conserved, since

$$
\begin{aligned}
\frac{d}{d t} \int_{0}^{x} h(x, t) d x & =-\int_{0}^{x}\left(h_{x x x x}+\left(g(h) h_{x}\right)_{x}\right) d x \\
& =-\left.\left(h_{x x x}+g(h) h_{x}\right)\right|_{x=0} ^{x=x} \\
& =0
\end{aligned}
$$

by periodicity.

Nonconstant steady states. Every constant function is a steady state of (14.1). We discuss the stability of these steady states at the end of the chapter.

To find nonconstant steady states, substitute $h=H(x)$ and solve:

$$
\begin{align*}
-H_{x x x x}-\left(g(H) H_{x}\right)_{x} & =0  \tag{14.2}\\
H_{x x x}+g(H) H_{x} & =\alpha \\
H_{x x}+G(H) & =\beta+\alpha x \\
H_{x x}+G(H) & =\beta
\end{align*}
$$

where G is an antiderivative of g ; here $\alpha=0$ is forced because the left side of the equation $\left(\mathrm{H}_{x x}+\mathrm{G}(\mathrm{H})\right)$ is periodic. This last equation describes a nonlinear oscillator, and it is well known how to construct solutions (one multiplies by $H_{x}$ and integrates). For example, when $g \equiv(2 \pi / X)^{2}$ we have steady states $H(x)=($ const. $)+\cos (2 \pi x / X)$. For the general case see [LaugesenPugh1].

Assume from now on that $\mathrm{H}(\mathrm{x})$ is a nonconstant steady state with period X .

## Linearized PDE

We perturb a steady state by considering

$$
h=H+\varepsilon \phi
$$

where the perturbation $\phi(x, t)$ is assumed to have mean value zero $\int_{0}^{x} \phi(x, t) d x=$ 0 ), so that fluid is conserved. Substituting $h$ into the equation (14.1) gives

$$
\begin{aligned}
0+\varepsilon \phi_{\mathrm{t}} & =-\left(\mathrm{H}_{x x x x}+\varepsilon \phi_{x x x x}\right)-\left(\mathrm{g}(\mathrm{H}+\varepsilon \phi)\left(\mathrm{H}_{x}+\varepsilon \phi_{x}\right)\right)_{x} \\
& =-\mathrm{H}_{x x x x}-\left(\mathrm{g}(\mathrm{H}) \mathrm{H}_{x}\right)_{x}-\varepsilon\left[\phi_{x x x}+\mathrm{g}(\mathrm{H}) \phi_{x}+\mathrm{g}^{\prime}(\mathrm{H}) \mathrm{H}_{x} \phi\right]_{x}+\mathrm{O}\left(\varepsilon^{2}\right) .
\end{aligned}
$$

The leading terms, of order $\varepsilon^{0}$, equal zero by the steady state equation for H. We discard terms of order $\varepsilon^{2}$ and higher. The remaining terms, of order $\varepsilon^{1}$, give the linearized equation:

$$
\begin{equation*}
\phi_{\mathrm{t}}=-\left[\phi_{x x}+\mathrm{g}(\mathrm{H}) \phi\right]_{x x} . \tag{14.3}
\end{equation*}
$$

Unfortunately, the operator on the right side is not symmetric (meaning it does not equal its formal adjoint). To make it symmetric, we "integrate up" the equation, as follows. Write

$$
\phi=\psi_{x}
$$

where $\psi$ is $X$-periodic (since $\phi$ has mean value zero). We may suppose $\psi$ has mean value zero at each time, by adding to $\psi$ a suitable function of $t$.

Substituting $\phi=\psi_{x}$ into (14.3) gives that

$$
\begin{aligned}
\psi_{\mathrm{t} x} & =-\left[\psi_{x x x}+\mathrm{g}(\mathrm{H}) \psi_{x}\right]_{x x} \\
\psi_{\mathrm{t}} & =-\left[\psi_{x x x}+\mathrm{g}(\mathrm{H}) \psi_{x}\right]_{x}
\end{aligned}
$$

(noting the constant of integration must equal 0 , by integrating both sides and using periodicity). Thus

$$
\psi_{t}=-\mathrm{L} \psi
$$

where $L$ is the symmetric operator

$$
\mathrm{L} w=w_{x x x x}+\left(\mathrm{g}(\mathrm{H}) w_{x}\right)_{x} .
$$

Separation of variables gives (formally) solutions of the form

$$
\psi=\sum_{j} c_{j} e^{-\tau_{j} t} w_{j}(x), \quad \phi=\sum_{j} c_{j} e^{-\tau_{j} t} w_{j}^{\prime}(x),
$$

where the eigenvalues $\tau_{j}$ and periodic eigenfunctions $w_{j}$ satisfy

$$
\mathrm{L} w_{j}=\tau_{j} w_{j}
$$

We conclude that the steady state H of the thin fluid film PDE is

## linearly unstable if $\tau_{1}<0$

because the perturbation $\phi$ grows to infinity, whereas the steady state is
linearly stable if $\tau_{1} \geq 0$
because $\phi$ remains bounded in that case. Remember these stability claims relate only to mean zero (volume preserving) perturbations.

To make these claims more rigorous, we need to understand the eigenvalue problem for L.

## Spectrum of L

We take:
$\Omega=\mathbb{T}=\mathbb{R} /(X \mathbb{Z})=$ torus of length $X$, so that functions on $\Omega$ are Xperiodic
$\mathcal{H}=\mathrm{L}^{2}(\mathbb{T})$, inner product $\langle u, v\rangle_{\mathrm{L}^{2}}=\int_{0}^{x} u v \mathrm{~d} x$
$\mathcal{K}=\mathrm{H}^{2}(\mathbb{T}) \cap\left\{u: \int_{0}^{x} u d x=0\right\}$, with inner product

$$
\langle u, v\rangle_{\mathrm{H}^{2}}=\int_{0}^{x}\left(u^{\prime \prime} v^{\prime \prime}+u^{\prime} v^{\prime}+u v\right) d x
$$

Compact imbedding $\mathcal{K} \hookrightarrow \mathrm{L}^{2}$ by Rellich-Kondrachov
Symmetric sesquilinear form

$$
a(u, v)=\int_{0}^{x}\left(u^{\prime \prime} v^{\prime \prime}-g(H) u^{\prime} v^{\prime}+C u v\right) d x
$$

where $\mathrm{C}>0$ is a sufficiently large constant to be chosen below.

Proof of ellipticity: The quantity $a(u, u)$ has a term of the form $-\left(u^{\prime}\right)^{2}$, whereas for $\|\mathfrak{u}\|_{\boldsymbol{H}^{2}}^{2}$ we need $+\left(\mathfrak{u}^{\prime}\right)^{2}$. To get around this obstacle we "hide" the $-\left(u^{\prime}\right)^{2}$ term inside the terms of the form $\left(u^{\prime \prime}\right)^{2}$ and $u^{2}$. Specifically,

$$
\begin{align*}
\int_{0}^{x}\left(u^{\prime}\right)^{2} d x & =-\int_{0}^{x} u^{\prime \prime} u d x \\
& \leq \int_{0}^{x}\left(\delta\left(u^{\prime \prime}\right)^{2}+(4 \delta)^{-1} u^{2}\right) d x \tag{14.4}
\end{align*}
$$

for any $\delta>0$. Here we used "Cauchy-with- $\delta$, which is the observation that for any $\alpha, \beta \in \mathbb{R}$,

$$
0 \leq\left(\sqrt{\delta} \alpha \pm(4 \delta)^{-1 / 2} \beta\right)^{2} \quad \Longrightarrow \quad|\alpha \beta| \leq \delta \alpha^{2}+(4 \delta)^{-1} \beta^{2}
$$

Next,

$$
\begin{aligned}
& a(u, u) \\
& \geq \int_{0}^{x}\left(\left(u^{\prime \prime}\right)^{2}-\left(\|g(H)\|_{L^{\infty}}+\frac{1}{2}\right)\left(u^{\prime}\right)^{2}+\frac{1}{2}\left(u^{\prime}\right)^{2}+C u^{2}\right) d x \\
& \geq \int_{0}^{x}\left(\left[1-\left(\|g(H)\|_{L^{\infty}}+\frac{1}{2}\right) \delta\right]\left(u^{\prime \prime}\right)^{2}+\frac{1}{2}\left(u^{\prime}\right)^{2}+\left[C-\left(\|g(H)\|_{L^{\infty}}+\frac{1}{2}\right)(4 \delta)^{-1}\right] u^{2}\right) d x
\end{aligned}
$$

$$
\begin{equation*}
\geq \frac{1}{2}\|u\|_{\mathrm{H}^{2}}^{2} \tag{14.4}
\end{equation*}
$$

provided we choose $\delta$ sufficiently small (depending on H ) and then choose C sufficiently large. Thus ellipticity holds.

The Spectral Theorem 4.1 now yields an ONB of eigenfunctions $\left\{w_{j}\right\}$ with eigenvalues $\gamma_{j}$ such that

$$
\mathrm{a}\left(w_{\mathrm{j}}, v\right)=\gamma_{\mathrm{j}}\left\langle w_{\mathrm{j}}, v\right\rangle_{\mathrm{L}^{2}} \quad \forall v \in \mathcal{K} .
$$

Writing $\gamma_{j}=\tau_{j}+C$ we get

$$
\int_{0}^{x}\left(w_{j}^{\prime \prime} v^{\prime \prime}-\mathrm{g}(\mathrm{H}) w_{\mathrm{j}}^{\prime} v^{\prime}\right) \mathrm{d} x=\tau_{\mathrm{j}} \int_{0}^{x} w_{\mathrm{j}} v \mathrm{~d} x \quad \forall v \in \mathcal{K} .
$$

These eigenfunctions satisfy $\mathrm{L} \boldsymbol{w}_{j}=\tau_{j} \boldsymbol{w}_{\mathrm{j}}$ weakly, and hence also classically (by elliptic regularity, since H and g are smooth).

Zero eigenvalue due to translational symmetry. We will show that $\tau=0$ is always an eigenvalue, with eigenfunction $u=H-\bar{H}$, where the constant $\overline{\mathrm{H}}$ equals the mean value of the steady state H . Indeed,

$$
\mathrm{Lu}=\mathrm{L}(\mathrm{H}-\overline{\mathrm{H}})=\mathrm{H}_{x x x x}+\left(\mathrm{g}(\mathrm{H}) \mathrm{H}_{x}\right)_{x}=0
$$

by the steady state equation (14.2).
This zero eigenvalue arises from a translational perturbation of the steady state, because choosing

$$
h=H(x+\varepsilon)=H(x)+\varepsilon H^{\prime}(x)+O\left(\varepsilon^{2}\right)
$$

gives rise to $\phi=\mathrm{H}^{\prime}$ and hence $\psi=\mathrm{H}-\overline{\mathrm{H}}$.

## Sufficient condition for linearized instability of nonconstant steady state H

Theorem 14.1 ([LaugesenPugh2, Th. 3]). If g is strictly convex then $\tau_{1}<0$.
For example, the theorem shows that nonconstant steady states are unstable with respect to volume-preserving perturbations if $g(y)=y^{p}$ with either $p>1$ or $p<0$.

Incidentally, the theorem is essentially the same as Theorem 13.1 for the reaction-diffusion PDE, simply writing $g$ instead of $f^{\prime}$ and noting that our periodic boundary conditions take care of the boundary terms in the integrations by parts.

Proof. The Rayleigh principle for L says that

$$
\tau_{1}=\min \left\{\frac{\int_{0}^{x}\left(\left(w^{\prime \prime}\right)^{2}-g(H)\left(w^{\prime}\right)^{2}\right) d x}{\int_{0}^{x} w^{2} d x}: w \in \mathrm{H}^{2}(\mathbb{T}) \backslash\{0\}, \int_{0}^{x} w \mathrm{~d} x=0\right\}
$$

We choose

$$
w=\mathrm{H}^{\prime},
$$

which is not the zero function since H is nonconstant, and note $w$ has mean value zero (as required), by periodicity of H . Then the numerator of the

Rayleigh quotient for $w$ is

$$
\begin{aligned}
& \int_{0}^{x}\left(\left(H^{\prime \prime \prime}\right)^{2}-g(H)\left(H^{\prime \prime}\right)^{2}\right) d x \\
& =\int_{0}^{x}\left(-H^{\prime \prime \prime \prime}-g(H) H^{\prime \prime}\right) H^{\prime \prime} d x \\
& \text { by parts } \\
& =\int_{0}^{x} g^{\prime}(H)\left(H^{\prime}\right)^{2} H^{\prime \prime} d x
\end{aligned} \quad \text { by the steady state equation (14.2) } \quad \begin{array}{ll}
=\frac{1}{3} \int_{0}^{x} g^{\prime}(H)\left[\left(H^{\prime}\right)^{3}\right]^{\prime} d x & \text { by parts } \\
=-\frac{1}{3} \int_{0}^{x} g^{\prime \prime}(H)\left(H^{\prime}\right)^{4} d x & \\
<0 &
\end{array}
$$

by convexity of g and since $\mathrm{H}^{\prime} \not \equiv 0$. Hence $\tau_{1}<0$, by the Rayleigh principle.

Motivation for the choice of trial function. Our trial function $w=\mathrm{H}^{\prime}$ corresponds to a perturbation $\phi=\mathrm{H}^{\prime \prime}$. This perturbation $h=\mathrm{H}+\varepsilon \phi=\mathrm{H}+\varepsilon \mathrm{H}^{\prime \prime}$ tends to push the steady state towards the constant function. The opposite perturbation would tend to push the steady state towards a "droplet" solution that equals 0 at some point. Thus our instability proof in Theorem 14.1 suggests (in the language of dynamical systems) that a heteroclinic connection might exist between the nonconstant steady state and the constant steady state, and similarly between the nonconstant steady state and a droplet steady state.

Linear stability of nonconstant steady states. It is more difficult to prove stability results, because lower bounds on the first eigenvalue are more difficult to prove (generally) than upper bounds.

See [LaugesenPugh2, §3.2] for some results when $g(y)=y^{p}, 0<p \leq 3 / 4$, based on time-map monotonicity ideas from the theory of reaction diffusion equations (see Chapter 13).

## Stability of constant steady states.

Let $\overline{\mathrm{H}}>0$ be constant. Then $\mathrm{H} \equiv \overline{\mathrm{H}}$ is a constant steady state. Its stability is easily determined, as follows.

Linearizing gives

$$
\phi_{t}=-\phi_{x x x x}-g(\overline{\mathrm{H}}) \phi_{x x}
$$

by (14.3), where the right side is linear and symmetric with constant coefficients.

We substitute the periodic Fourier mode $\phi=e^{-\tau t} \exp (2 \pi i k x / X)$, where $k \in \mathbb{Z}$ (and $k \neq 0$ since our perturbations have mean value zero), obtaining the eigenvalue

$$
\tau=\left(\frac{2 \pi \mathrm{k}}{\mathrm{X}}\right)^{2}\left(\left(\frac{2 \pi \mathrm{k}}{\mathrm{X}}\right)^{2}-\mathrm{g}(\overline{\mathrm{H}})\right)
$$

If $\mathrm{g} \leq 0$ (which means the second order term in the thin film PDE behaves like a forwards heat equation), then $\tau \geq 0$ for each $k$, and so all constant steady states are linearly stable.

If $g>0$ and $\left(\frac{2 \pi}{x}\right)^{2} \geq g(\overline{\mathrm{H}})$ then $\tau \geq 0$ for each $k$, and so the constant steady states $\overline{\mathrm{H}}$ is linearly stable.

If $\mathrm{g}>0$ and $\left(\frac{2 \pi}{\mathrm{x}}\right)^{2}<\mathrm{g}(\overline{\mathrm{H}})$ then the constant steady states $\overline{\mathrm{H}}$ is linearly unstable with respect to the $k=1$ mode (and possibly other modes too). In particular, this occurs if $X$ is large enough. Hence we call the thin film PDE "long-wave unstable" if $\mathrm{g}>0$, since constant steady states are then unstable with respect to perturbations of sufficiently long wavelength.

## Part II

## Continuous Spectrum

## Looking ahead to continuous spectrum

The discrete spectral theory in Part I of the course generated, in each application,

- eigenfunctions $\left\{\mathfrak{u}_{j}\right\}$ with "discrete" spectrum $\lambda_{1}, \lambda_{2}, \lambda_{3}, \ldots$ satisfying $L u_{j}=\lambda_{j} u_{j}$ where $L$ is a symmetric differential operator, together with
- a spectral decomposition (or "resolution") of each $f \in L^{2}$ into a sum of eigenfunctions: $f=\sum_{j}\left\langle f, u_{j}\right\rangle u_{j}$.

These constructions depended heavily on symmetry of the differential operator L (which ensured symmetry of the sesquilinear form a) and on compactness of the imbedding of the Hilbert space $\mathcal{K}$ into $\mathcal{H}$.

For the remainder of the course we retain the symmetry assumption on the operator, but drop the compact imbedding assumption. The resulting "continuous" spectrum leads to a decomposition of $f \in L^{2}$ into an integral of "almost eigenfunctions".

We begin with examples, and later put the examples in context by developing some general spectral theory for unbounded, selfadjoint differential operators.

## Chapter 15

## Computable example: Laplacian (free Schrödinger) on all of space

## Goal

To determine for the Laplacian on Euclidean space its continuous spectrum $[0, \infty)$, and the associated spectral decomposition of $L^{2}$.

## Spectral decomposition

The Laplacian $\mathrm{L}=-\Delta$ on a bounded domain has discrete spectrum, as we saw in Chapters 2 and 5 . When the domain expands to be all of space, though, the Laplacian has no eigenvalues at all. For example in 1 dimension, solutions of $-\mathfrak{u}^{\prime \prime}=\lambda u$ are linear combinations of $e^{ \pm i \sqrt{\lambda x}}$, which oscillates if $\lambda>0$, or is constant if $\lambda=0$, or grows in one direction or the other if $\lambda \in \mathbb{C} \backslash[0, \infty)$. In none of these situations does $u$ belong to $L^{2}$. (In all dimensions we can argue as follows: if $-\Delta \boldsymbol{u}=\lambda \boldsymbol{u}$ and $\boldsymbol{u} \in \mathrm{L}^{2}$ then by taking Fourier transforms, $4 \pi^{2}|\xi|^{2} \widehat{\mathfrak{u}}(\xi)=\widehat{\mathfrak{u}}(\xi)$ a.e., and so $\widehat{\mathfrak{u}}=0$ a.e. Thus no $\mathrm{L}^{2}$-eigenfunctions exist.)

A fundamental difference between the whole space case and the case of bounded domains is that the imbedding $\mathrm{H}^{1}\left(\mathbb{R}^{\mathrm{d}}\right) \hookrightarrow \mathrm{L}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)$ is not compact. For example, given any nonzero $f \in H^{1}(\mathbb{R})$, the functions $f(\cdot-k)$ are bounded in $L^{2}(\mathbb{R})$, but have no $L^{2}$-convergent subsequence as $k \rightarrow \infty$. Hence the discrete spectral theorem (Theorem 4.1) is inapplicable.

Nevertheless, the Laplacian $-\Delta$ on $\mathbb{R}^{\mathrm{d}}$ has:

1. generalized eigenfunctions

$$
v_{\omega}(x)=e^{2 \pi i \omega \cdot x}, \quad \omega \in \mathbb{R}^{\mathrm{d}}
$$

(note that $v_{\omega}$ is bounded, but it is not an eigenfunction since $v_{\omega} \notin \mathrm{L}^{2}$ ) which satisfy the eigenfunction equation $-\Delta v_{\omega}=\lambda \nu_{\omega}$ with generalized eigenvalue

$$
\lambda=\lambda(\omega)=4 \pi^{2}|\omega|^{2}
$$

## Generalized eigenfn: $\operatorname{Re} \boldsymbol{v}_{\mathbf{2}}(\mathbf{x})$


2. and a spectral decomposition

$$
\mathrm{f}=\int_{\mathbb{R}^{\mathrm{d}}}\left\langle\mathrm{f}, v_{\omega}\right\rangle v_{\omega} \mathrm{d} \omega, \quad \forall \mathrm{f} \in \mathrm{~L}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)
$$

Proof of spectral decomposition. Since

$$
\left\langle f, v_{\omega}\right\rangle=\int_{\mathbb{R}^{\mathrm{d}}} f(x) e^{-2 \pi i \omega \cdot x} d x=\widehat{f}(\omega),
$$

the spectral decomposition simply says

$$
f(x)=\int_{\mathbb{R}^{\mathrm{d}}} \widehat{f}(\omega) e^{2 \pi i \omega \cdot x} d \omega,
$$

which is the Fourier inversion formula.

## Application of spectral decomposition

One may solve evolution equations by separating variables: for example, the heat equation $u_{t}=\Delta u$ with initial condition $h(x)$ has solution

$$
u(x, t)=\int_{\mathbb{R}^{d}} \widehat{h}(\omega) e^{-\lambda(\omega) t} v_{\omega}(x) d \omega
$$

Note the analogy to the series solution by separation of variables, in the case of discrete spectrum.

Aside. Typically, one evaluates the last integral (an inverse Fourier transform) and thus obtains a convolution of the initial data $h$ and the fundamental solution of the heat equation (which is the inverse transform of $e^{-\lambda(\omega) t}$ ).

Continuous spectrum $=[0, \infty)$
The generalized eigenvalue $\lambda \geq 0$ is "almost" an eigenvalue, in two senses:

- the eigenfunction equation $(-\Delta-\lambda) \mathfrak{u}=0$ does not have a solution in $\mathrm{L}^{2}$, but it does have a solution $v_{\omega} \in \mathrm{L}^{\infty}$,
- a Weyl sequence exists for $-\Delta$ and $\lambda$, meaning there exist functions $w_{n}$ such that
(W1) $\left\|(-\Delta-\lambda) w_{\mathrm{n}}\right\|_{\mathrm{L}^{2}} \rightarrow 0$ as $\mathrm{n} \rightarrow \infty$,
(W2) $\left\|w_{n}\right\|_{\mathrm{L}^{2}}=1$,
(W3) $w_{n} \rightharpoonup 0$ weakly in $L^{2}$ as $n \rightarrow \infty$.
We prove existence of a Weyl sequence in the Proposition below. Later we will define the continuous spectrum to consist of those $\lambda$-values for which a Weyl sequence exists. Thus the continuous spectrum of $-\Delta$ is precisely the nonnegative real axis. Recall it is those values of $\lambda$ that entered into our spectral decomposition earlier in the chapter.

Remark. Existence of a Weyl sequence ensures that $(-\Delta-\lambda)$ does not have a bounded inverse from $L^{2} \rightarrow L^{2}$, for if we write $f_{n}=(-\Delta-\lambda) w_{n}$ then

$$
\frac{\left\|(-\Delta-\lambda)^{-1} f_{n}\right\|_{L^{2}}}{\left\|f_{n}\right\|_{L^{2}}}=\frac{\left\|w_{n}\right\|_{L^{2}}}{\left\|(-\Delta-\lambda) w_{n}\right\|_{L^{2}}} \rightarrow \infty
$$

as $n \rightarrow \infty$, by (W1) and (W2). In this way, existence of a Weyl sequence is similar to existence of an eigenfunction, which also prevents invertibility of $(-\Delta-\lambda)$.

| $\mathbb{C}$ |  |
| :---: | :---: |
|  | $\operatorname{spec}(-\Delta)$ |

Proposition 15.1 (Weyl sequences for negative Laplacian). A Weyl sequence exists for $-\Delta$ and $\lambda \in \mathbb{C}$ if and only if $\lambda \in[0, \infty)$.
Proof. " $\Longleftarrow " ~ F i x ~ \lambda \in[0, \infty)$ and choose $\omega \in \mathbb{R}^{\text {d }}$ with $4 \pi|\omega|^{2}=\lambda$. Take a cut-off function $\kappa \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\kappa \equiv 1$ on the unit ball $B(1)$ and $\kappa \equiv 0$ on $\mathbb{R}^{d} \backslash B(2)$. Define a cut-off version of the generalized eigenfunction, by

$$
w_{n}=c_{n} \kappa\left(\frac{x}{n}\right) v_{\omega}(x)
$$

where the normalizing constant is

$$
c_{n}=\frac{1}{n^{d / 2}\|\kappa\|_{L^{2}}} .
$$



First we prove (W1). We have

$$
\begin{aligned}
& (\lambda+\Delta) w_{n} \\
& =c_{n}\left(\lambda v_{\omega}+\Delta v_{\omega}\right) \kappa\left(\frac{x}{n}\right)+2 \frac{c_{n}}{n} \nabla v_{\omega}(x) \cdot(\nabla \kappa)\left(\frac{x}{n}\right)+\frac{c_{n}}{n^{2}} v_{\omega}(x)(\Delta \kappa)\left(\frac{x}{n}\right) .
\end{aligned}
$$

The first term vanishes because $\Delta v_{\omega}=-4 \pi|\omega|^{2} v_{\omega}$ pointwise. In the third term, note that $\nu_{\omega}$ is a bounded function, and that a change of variable shows

$$
\frac{\mathrm{c}_{\mathrm{n}}}{\mathrm{n}^{2}}\left\|(\Delta \kappa)\left(\frac{\dot{1}}{\mathrm{n}}\right)\right\|_{\mathrm{L}^{2}}=\frac{1}{\mathrm{n}^{2}} \frac{\|\Delta K\|_{\mathrm{L}^{2}}}{\|\kappa\|_{\mathrm{L}^{2}}} \rightarrow 0
$$

The second term similarly vanishes in the limit, as $\mathfrak{n} \rightarrow \infty$. Hence $(\lambda+$ $\Delta) w_{n} \rightarrow 0$ in $\mathrm{L}^{2}$, which is ( W 1 ).

For (W2) we simply observe that $\left|v_{\omega}(x)\right|=1$ pointwise, so that $\left\|w_{n}\right\|_{L^{2}}=$ 1 by a change of variable, using the definition of $\mathrm{c}_{\mathrm{n}}$.

To prove (W3), take $f \in L^{2}$ and let $R>0$. We decompose $f$ into "near" and "far" components, as $f=g+h$ where $g=f 1_{B(R)}$ and $h=f 1_{\mathbb{R}^{d} \backslash B(R)}$. Then

$$
\left\langle\mathrm{f}, w_{\mathrm{n}}\right\rangle_{\mathrm{L}^{2}}=\left\langle\mathrm{g}, w_{\mathrm{n}}\right\rangle_{\mathrm{L}^{2}}+\left\langle\mathrm{h}, w_{\mathrm{n}}\right\rangle_{\mathrm{L}^{2}} .
$$

We have

$$
\left|\left\langle\boldsymbol{g}, w_{n}\right\rangle_{\mathrm{L}^{2}}\right| \leq \mathrm{c}_{\mathrm{n}}\|\kappa\|_{\mathrm{L}^{\infty}}\|\mathrm{g}\|_{\mathrm{L}^{1}} \rightarrow 0
$$

as $n \rightarrow \infty$, since $c_{n} \rightarrow 0$. Also, by Cauchy-Schwarz and (W2) we see

$$
\limsup _{n \rightarrow \infty}\left|\left\langle h, w_{n}\right\rangle_{\mathrm{L}^{2}}\right| \leq\|h\|_{\mathrm{L}^{2}} .
$$

This last quantity can be made arbitrarily small by letting $R \rightarrow \infty$, and so $\lim _{n \rightarrow \infty}\left\langle\mathrm{f}, w_{n}\right\rangle_{\mathrm{L}^{2}}=0$. That is, $w_{n} \rightharpoonup 0$ weakly.
$" \Longrightarrow "$ Assume $\lambda \in \mathbb{C} \backslash[0, \infty)$, and let

$$
\delta=\operatorname{dist}(\lambda,[0, \infty))
$$

so that $\delta>0$.
Suppose (W1) holds, and write $g_{n}=(-\Delta-\lambda) w_{n}$. Then

$$
\begin{aligned}
\widehat{g_{n}}(\xi) & =\left(4 \pi^{2}|\xi|^{2}-\lambda\right) \widehat{w_{n}}(\xi) \\
\widehat{w_{n}}(\xi) & =\frac{1}{\left(4 \pi^{2}|\xi|^{2}-\lambda\right)} \widehat{g_{n}}(\xi) \\
\left|\widehat{w_{n}}(\xi)\right| & \leq \delta^{-1}\left|\widehat{g_{n}}(\xi)\right|
\end{aligned}
$$

and hence

$$
\begin{aligned}
\left\|w_{\mathrm{n}}\right\|_{\mathrm{L}^{2}}=\left\|\widehat{w_{\mathrm{n}}}\right\|_{\mathrm{L}^{2}} & \leq \delta^{-1}\left\|\widehat{\mathrm{~g}_{\mathrm{n}}}\right\|_{\mathrm{L}^{2}} \\
& =\delta^{-1}\left\|g_{\mathrm{n}}\right\|_{\mathrm{L}^{2}} \\
& \rightarrow 0
\end{aligned}
$$

by (W1). Thus (W2) does not hold.
(Aside. The calculations above show, in fact, that $(-\Delta-\lambda)^{-1}$ is bounded from $\mathrm{L}^{2} \rightarrow \mathrm{~L}^{2}$ with norm bound $\delta^{-1}$, when $\lambda \notin[0, \infty)$.)

## Chapter 16

## Computable example: Schrödinger with a bounded potential well

Goal
To show that the Schrödinger operator

$$
\mathrm{L}=-\frac{\mathrm{d}^{2}}{\mathrm{~d} x^{2}}-2 \operatorname{sech}^{2} x
$$

in 1 dimension has a single negative eigenvalue (discrete spectrum) as well as nonnegative continuous spectrum $[0, \infty)$. The spectral decomposition will show the potential is reflectionless.


Reference [Keener] Section 7.5

## Discrete spectrum $=\{-1\}$

We claim -1 is an eigenvalue of $L$ with eigenfunction sech $x$. This fact can be checked directly, but we will proceed more systematically by factoring the Schrödinger operator with the help of the first order operators

$$
\begin{aligned}
& \mathrm{L}^{+}=-\frac{\mathrm{d}}{\mathrm{dx}}+\tanh x \\
& \mathrm{~L}^{-}=\frac{\mathrm{d}}{\mathrm{dx}}+\tanh x
\end{aligned}
$$

We compute

$$
\begin{aligned}
\mathrm{L}^{+} \mathrm{L}^{-}-1 & =\left(-\frac{\mathrm{d}}{\mathrm{dx}}+\tanh x\right)\left(\frac{\mathrm{d}}{\mathrm{dx}}+\tanh x\right)-1 \\
& =-\frac{\mathrm{d}^{2}}{\mathrm{dx}}-(\tanh x)^{\prime}+\tanh ^{2} x-1 \\
& =-\frac{\mathrm{d}^{2}}{\mathrm{dx} x^{2}}-2 \operatorname{sech}^{2} x \\
& =\mathrm{L}
\end{aligned}
$$

since $(\tanh )^{\prime}=\operatorname{sech}^{2}$ and $1-\tanh ^{2}=\operatorname{sech}^{2}$. Thus

$$
\begin{equation*}
\mathrm{L}=\mathrm{L}^{+} \mathrm{L}^{-}-1 \tag{16.1}
\end{equation*}
$$

It follows that functions in the kernel of $\mathrm{L}^{-}$are eigenfunctions of L with eigenvalue $\lambda=-1$. To find the kernel we solve:

$$
\begin{aligned}
\mathrm{L}^{-} v & =0 \\
v^{\prime}+(\tanh x) v & =0 \\
(\cosh x) v^{\prime}+(\sinh x) v & =0 \\
(\cosh x) v & =\text { const. } \\
v & =\mathrm{c} \operatorname{sech} x
\end{aligned}
$$

Clearly sech $x \in L^{2}(\mathbb{R})$, since sech decays exponentially. Thus -1 lies in the discrete spectrum of L, with eigenfunction sech $x$.

Are there any other eigenvalues? No! Argue as follows. By composing

## Bound state, with energy -1


$\mathrm{L}^{+}$and $\mathrm{L}^{-}$in the reverse order we find

$$
\begin{align*}
\mathrm{L}^{-} \mathrm{L}^{+}-1 & =\left(\frac{\mathrm{d}}{\mathrm{dx}}+\tanh x\right)\left(-\frac{\mathrm{d}}{\mathrm{dx}}+\tanh x\right)-1 \\
& =-\frac{\mathrm{d}^{2}}{\mathrm{dx} x^{2}}+(\tanh x)^{\prime}+\tanh ^{2} x-1 \\
& =-\frac{\mathrm{d}^{2}}{\mathrm{~d} x^{2}} . \tag{16.2}
\end{align*}
$$

From (16.1) and (16.2) we deduce

$$
-\frac{\mathrm{d}^{2}}{\mathrm{dx}} \mathrm{x}^{-}=\mathrm{L}^{-} \mathrm{L}
$$

Thus if $\mathrm{L} v=\lambda v$ then $-\frac{\mathrm{d}^{2}}{\mathrm{dx}}\left(\mathrm{L}^{-} v\right)=\mathrm{L}^{-} \mathrm{L} v=\lambda\left(\mathrm{L}^{-} v\right)$. By solving for $\mathrm{L}^{-} v$ in terms of $e^{ \pm i \sqrt{\lambda} x}$, and then integrating to obtain $\nu$, we conclude after some thought (omitted) that the only way for $v$ to belong to $\mathrm{L}^{2}(\mathbb{R})$ is to have $\mathrm{L}^{-} v=0$ and hence $v=\mathrm{c} \operatorname{sech} x$, so that $\lambda=-1$.

## Continuous spectrum $\supset[0, \infty)$

Let $\lambda \in[0, \infty)$. Generalized eigenfunctions with $\mathrm{L} v=\lambda \nu$ certainly exist: choose $\omega \in \mathbb{R}$ with $\lambda=4 \pi^{2} \omega^{2}$ and define

$$
v(x)=\mathrm{L}^{+}\left(e^{2 \pi i \omega x}\right)=(\tanh x-2 \pi i \omega) e^{2 \pi i \omega x},
$$

which is bounded but not square integrable. We compute

$$
\begin{aligned}
\mathrm{L} v & =\left(\mathrm{L}^{+} \mathrm{L}^{-}-1\right) \mathrm{L}^{+}\left(e^{2 \pi i \omega x}\right) \\
& =\mathrm{L}^{+}\left(\mathrm{L}^{-} \mathrm{L}^{+}-1\right)\left(e^{2 \pi i \omega x}\right) \\
& =-\mathrm{L}^{+} \frac{\mathrm{d}^{2}}{\mathrm{dx} x^{2}}\left(e^{2 \pi i \omega x}\right) \\
& =\mathrm{L}^{+}\left(4 \pi^{2} \omega^{2} e^{2 \pi i \omega x}\right) \\
& =\lambda \nu,
\end{aligned}
$$

which verifies that $v(x)$ is a generalized eigenfunction.
We can further prove existence of a Weyl sequence for $L$ and $\lambda$ by adapting Lemma 15.1 " $\Longleftarrow$ ", using the same Weyl functions $w_{\mathfrak{n}}(x)$ as for the free Schrödinger operator $-\Delta$. The only new step in the proof, for proving $\|(\mathrm{L}-$ $\lambda) w_{n} \|_{L^{2}} \rightarrow 0$ in (W1), is to observe that

$$
\begin{aligned}
\left|2 \operatorname{sech}^{2} x w_{n}(x)\right| & =2{c_{n}}\left|\kappa\left(\frac{x}{n}\right) e^{2 \pi i \omega x}\right| \operatorname{sech}^{2} x \\
& \leq 2 c_{n}\|\kappa\|_{L^{\infty}} \operatorname{sech}^{2} x \\
& \rightarrow 0
\end{aligned}
$$

in $L^{2}(\mathbb{R})$ as $n \rightarrow \infty$, because $c_{n} \rightarrow 0$. (Note. This part of the proof works not only for the sech ${ }^{2}$ potential, but for any potential belonging to $L^{2}$.)

We have shown that the continuous spectrum contains $[0, \infty)$. We will prove the reverse inclusion at the end of the chapter.

Generalized eigenfunctions as traveling waves. The eigenfunction ("bound state") $v(x)=\operatorname{sech} x$ with eigenvalue ("energy") -1 produces a standing wavefunction

$$
u=e^{i t} \operatorname{sech} x
$$

satisfying the time-dependent Schrödinger equation

$$
\mathfrak{i} u_{t}=\mathrm{Lu}
$$

The generalized eigenfunction

$$
\begin{equation*}
v(x)=(\tanh x-2 \pi i \omega) e^{2 \pi i \omega x} \tag{16.3}
\end{equation*}
$$

with generalized eigenvalue $\lambda=4 \pi^{2} \omega^{2}$ similarly produces a standing wave

$$
u=e^{-i 4 \pi^{2} \omega^{2} t}(\tanh x-2 \pi i \omega) e^{2 \pi i \omega x}
$$



More usefully, we rewrite this formula as a traveling plane wave multiplied by an $x$-dependent amplitude:

$$
\begin{equation*}
u=(\tanh x-2 \pi i \omega) e^{2 \pi i \omega(x-2 \pi \omega t)} \tag{16.4}
\end{equation*}
$$

The amplitude factor serves to quantify the effect of the potential on the traveling wave: in the absence of a potential, the amplitude would be identically 1 , since the plane wave $e^{2 \pi i \omega(x-2 \pi \omega t)}$ solves the free Schrödinger equation $\mathfrak{i u} u_{t}=-\Delta u$.

Reflectionless nature of the potential, and a nod to scattering theory. One calls the potential $-2 \operatorname{sech}^{2} x$ "reflectionless" because the rightmoving wave in (16.4) passes through the potential with none of its energy reflected into a left-moving wave. In other words, the generalized eigenfunction (16.3) has the form $c e^{2 \pi i \omega x}$ both as $x \rightarrow-\infty$ and as $x \rightarrow \infty$ (with different constants, it turns out, although the constants are equal in magnitude).

This reflectionless property is unusual. A typical Schrödinger potential would produce generalized eigenfunctions equalling approximately

$$
c_{I} e^{2 \pi i \omega x}+c_{R} e^{-2 \pi i \omega x} \quad \text { as } x \rightarrow-\infty
$$

and

$$
\mathrm{c}_{\mathrm{T}} \mathrm{e}^{2 \pi i \omega x} \quad \text { as } x \rightarrow \infty
$$

(or similarly with the roles of $\pm \infty$ interchanged). Here $\left|\mathbf{c}_{\boldsymbol{I}}\right|$ is the amplitude of the incident right-moving wave, $\left|\mathbf{c}_{\boldsymbol{R}}\right|$ is the amplitude of the left-moving wave reflected by the potential, and $\left|\mathbf{c}_{T}\right|$ is the amplitude of the right-moving wave transmitted through the potential. Conservation of L2-energy demands that

$$
\left|c_{I}\right|^{2}=\left|{c_{R}}\right|^{2}+\left|c_{T}\right|^{2} .
$$

For a gentle introduction to this "scattering theory" see [Keener, Section 7.5]. Then one can proceed to the book-length treatment in [ReedSimon3].

## Spectral decomposition of $L^{2}$

Analogous to an orthonormal expansion in terms of eigenfunctions, we have:

## Theorem 16.1.

$$
\mathrm{f}=\frac{1}{2}\langle\mathrm{f}, \operatorname{sech}\rangle \operatorname{sech}+\int_{\mathbb{R}}\left\langle\mathrm{f}, \mathrm{~L}^{+} v_{\omega}\right\rangle \mathrm{L}^{+} v_{\omega} \frac{\mathrm{d} \omega}{1+4 \pi^{2} \omega^{2}}, \quad \forall \mathrm{f} \in \mathrm{~L}^{2}(\mathbb{R})
$$

where $\mathrm{L}^{+} v_{\omega}(x)=(\tanh \mathrm{x}-2 \pi i \omega) \mathrm{e}^{2 \pi i \omega x}$ is the generalized eigenfunction at frequency $\omega$.

The discrete part of the decomposition has the same form as the continuous part, in fact, because sech $=-\mathrm{L}^{+}(\sinh )$.

Proof. We will sketch the main idea of the proof, and leave it to the reader to make the argument rigorous.

By analogy with an orthonormal expansion in the discrete case, we assume that $f \in L^{2}(\mathbb{R})$ has a decomposition in terms of the eigenfunction sech $x$ and the generalized eigenfunctions $\mathrm{L}^{+} \nu_{\omega}$ in the form

$$
f=c\langle f, \operatorname{sech}\rangle \operatorname{sech}+\int_{\mathbb{R}} m_{f}(\omega)\left\langle f, L^{+} v_{\omega}\right\rangle L^{+} v_{\omega} d \omega
$$

where the coefficient $c$ and multiplier $m_{f}(\boldsymbol{\omega})$ are to be determined.
Taking the inner product with sech $x$ implies that $\mathrm{c}=\frac{1}{2}$, since $\|$ sech $\|_{\mathrm{L}^{2}(\mathbb{R})}^{2}=$ 2 and $\left\langle\mathrm{L}^{+} v_{\omega}\right.$, sech $\rangle=\left\langle v_{\omega}, \mathrm{L}^{-}\right.$sech $\rangle=0$.

Next we annihilate the sech term by applying $\mathrm{L}^{-}$to both sides:

$$
\mathrm{L}^{-} \mathrm{f}=\mathrm{L}^{-}\left(\int_{\mathbb{R}} m_{\mathrm{f}}(\omega)\left\langle\mathrm{f}, \mathrm{~L}^{+} v_{\omega}\right\rangle \mathrm{L}^{+} v_{\omega} \mathrm{d} \omega\right)
$$

Note that by integration by parts,

$$
\left.\left\langle f, L^{+} v_{\omega}\right\rangle=\left\langle L^{-} f, v_{\omega}\right\rangle=\widehat{\left(L^{-} f\right.}\right)(\omega)
$$

Hence

$$
\begin{aligned}
\mathrm{L}^{-} \mathrm{f} & =\mathrm{L}^{-}\left(\int_{\mathbb{R}} m_{\mathrm{f}}(\boldsymbol{\omega}) \widehat{\left(\mathrm{L}^{-f}\right)}(\boldsymbol{\omega}) \mathrm{L}^{+} v_{\omega} \mathrm{d} \omega\right) \\
& =\int_{\mathbb{R}} m_{\mathrm{f}}(\omega) \widehat{\left(\mathrm{L}^{-f}\right)}(\omega) \mathrm{L}^{-} \mathrm{L}^{+} v_{\omega} d \omega \\
& =\int_{\mathbb{R}} m_{\mathrm{f}}(\omega) \widehat{\left(\mathrm{L}^{-f}\right)}(\omega)\left(1+4 \pi^{2} \omega^{2}\right) v_{\omega} d \omega
\end{aligned}
$$

by (16.2). Thus the multiplier should be $\mathfrak{m}_{\mathrm{f}}(\omega)=1 /\left(1+4 \pi^{2} \omega^{2}\right)$, in order for Fourier inversion to hold. This argument shows the necessity of the formula in the theorem, and one can show sufficiency by suitably reversing the steps.

The theorem implies a Plancherel type identity.

## Corollary 16.2.

$$
\|f\|_{L^{2}}^{2}=\frac{1}{2} \left\lvert\,\left.\langle f, \text { sech }\rangle\right|^{2}+\int_{\mathbb{R}}\left|\left\langle f, L^{+} v_{\omega}\right\rangle\right|^{2} \frac{d \omega}{1+4 \pi^{2} \omega^{2}}\right., \quad \forall f \in L^{2}(\mathbb{R}) .
$$

Proof. Take the inner product of f with the formula in Theorem 16.1.

## Continuous spectrum $=[0, \infty)$

Earlier we showed that the continuous spectrum contains $[0, \infty)$. For the reverse containment, suppose $\lambda \notin[0, \infty)$ and $\lambda \neq-1$. Then $L-\lambda$ is invertible on $\mathrm{L}^{2}$, with

$$
(L-\lambda)^{-1} f=-\frac{1}{\lambda+1} \frac{1}{2}\langle f, \text { sech }\rangle \operatorname{sech}+\int_{\mathbb{R}} \frac{\left\langle f, L^{+} v_{\omega}\right\rangle}{4 \pi^{2} \omega^{2}-\lambda} L^{+} v_{\omega} \frac{d \omega}{1+4 \pi^{2} \omega^{2}}
$$

as one sees by applying $L-\lambda$ to both sides and recalling Theorem 16.1. To check the boundedness of this inverse, note that

$$
\begin{aligned}
\left\|(\mathrm{L}-\lambda)^{-1} \mathrm{f}\right\|_{\mathrm{L}^{2}}^{2} & =\frac{1}{|\lambda+1|^{2}} \frac{1}{2} \left\lvert\,\left.\langle f, \text { sech }\rangle\right|^{2}+\int_{\mathbb{R}} \frac{\left|\left\langle f, \mathrm{~L}^{+} v_{\omega}\right\rangle\right|^{2}}{\left|4 \pi^{2} \omega^{2}-\lambda\right|^{2}} \frac{\mathrm{~d} \omega}{1+4 \pi^{2} \omega^{2}}\right. \\
& \leq \frac{1}{|\lambda+1|^{2}} \frac{1}{2} \left\lvert\,\left.\langle f, \text { sech }\rangle\right|^{2}+\frac{1}{\operatorname{dist}(\lambda,[0, \infty))^{2}} \int_{\mathbb{R}}\left|\left\langle f, \mathrm{~L}^{+} v_{\omega}\right\rangle\right|^{2} \frac{\mathrm{~d} \omega}{1+4 \pi^{2} \omega^{2}}\right. \\
& \leq \text { (const.) }\|f\|_{\mathrm{L}^{2}}^{2},
\end{aligned}
$$

where we used Corollary 16.2.
The boundedness of $(\mathrm{L}-\lambda)^{-1}$ implies that the Weyl conditions (W1) and (W2) cannot both hold. Thus no Weyl sequence can exist for $\lambda$, so that $\lambda$ does not belong to the continuous spectrum.

Next suppose $\lambda=-1$. If a Weyl sequence $w_{n}$ exists, then

$$
\left\langle w_{n}, \text { sech }\right\rangle_{\mathrm{L}^{2}} \rightarrow 0 \quad \text { as } n \rightarrow \infty
$$

by the weak convergence in (W3). Hence if we project away from the $\lambda=-1$ eigenspace by defining

$$
y_{n}=w_{n}-\frac{1}{2}\left\langle w_{n}, \operatorname{sech}\right\rangle_{\mathrm{L}^{2}} \text { sech } \quad \text { and } \quad z_{n}=y_{n} /\left\|y_{n}\right\|_{L^{2}}
$$

then we find $\left\|y_{n}\right\|_{\mathrm{L}^{2}} \rightarrow 1$ and $\left\|z_{n}\right\|_{\mathrm{L}^{2}}=1$, with $\left\langle z_{\mathrm{n}} \text {, sech }\right\rangle_{\mathrm{L}^{2}}=0$. Also

$$
(\mathrm{L}+1) z_{\mathrm{n}}=(\mathrm{L}+1) \mathrm{y}_{\mathrm{n}} /\left\|\mathrm{y}_{\mathrm{n}}\right\|_{\mathrm{L}^{2}}=(\mathrm{L}+1) w_{\mathrm{n}} /\left\|y_{n}\right\|_{\mathrm{L}^{2}} \rightarrow 0
$$

in $L^{2}$. Thus $z_{n}$ satisfies (W1) and (W2) and lies in the orthogonal complement of the eigenspace spanned by sech. A contradiction now follows from the boundedness of $(L+1)^{-1}$ on that orthogonal complement (with the boundedness being proved by the same argument as above for $\lambda \neq-1$ ). This contradiction shows that no such Weyl sequence $\boldsymbol{w}_{n}$ can exist, and so -1 does not belong to the continuous spectrum.

Note. The parallels with our derivation of the continuous spectrum for the Laplacian in Chapter 15 are instructive.

## Chapter 17

## Selfadjoint, unbounded linear operators

## Goal

To develop the theory of unbounded linear operators on a Hilbert space, and to define selfadjointness for such operators.

References [GustafsonSigal] Sections 1.5, 2.4
[HislopSigal] Chapters 4, 5

## Motivation

Now we should develop some general theory, to provide context for the examples computed in Chapters 15 and 16.

We begin with a basic principle of calculus:
integration makes functions better, while differentiation makes them worse.

More precisely, integral operators are bounded (generally speaking), while differential operators are unbounded. For example, $e^{2 \pi i n x}$ has norm 1 in $L^{2}[0,1]$ while its derivative $\frac{d}{d x} e^{2 \pi i n x}=2 \pi i n e^{2 \pi i n x}$ has norm that grows with n . The unboundedness of such operators prevents us from applying the spectral theory of bounded operators on a Hilbert space.

Further, differential operators are usually defined only on a (dense) subspace of our natural function spaces. In particular, we saw in our study
of discrete spectra that the Laplacian is most naturally studied using the Sobolev space $\mathrm{H}^{1}$, even though the Laplacian involves two derivatives and $\mathrm{H}^{1}$-functions are guaranteed only to possess a single derivative.

To meet these challenges, we will develop the theory of densely defined, unbounded linear operators, along with the notion of adjoints and selfadjointness for such operators.

## Domains and inverses of (unbounded) operators

Take a complex Hilbert space $\mathcal{H}$ with inner product $\langle\cdot, \cdot\rangle$. Suppose $\mathcal{A}$ is a linear operator (not necessarily bounded) from a subspace $\mathrm{D}(\mathcal{A}) \subset \mathcal{H}$ into $\mathcal{H}$ :

$$
A: D(A) \rightarrow \mathcal{H}
$$

Call $D(A)$ the domain of $A$.
An operator $B$ with domain $D(B)$ is called the inverse of $A$ if

- $D(B)=\operatorname{Ran}(A), D(A)=\operatorname{Ran}(B)$, and
- $B A=\operatorname{id}_{\operatorname{Ran}(B)}, A B=\operatorname{id}_{\operatorname{Ran}(A)}$.

Write $A^{-1}$ for this inverse, if it exists. Obviously $A^{-1}$ is unique, if it exists, because in that case $\mathcal{A}$ is bijective.

Further say $A$ is invertible if $A^{-1}$ exists and is bounded on $\mathcal{H}$ (meaning that $A^{-1}$ exists, $\operatorname{Ran}(A)=\mathcal{H}$, and $A^{-1}: \mathcal{H} \rightarrow \mathcal{H}$ is a bounded linear operator).
Example. Consider the operator $A=-\Delta+1$ with domain $H^{2}\left(\mathbb{R}^{d}\right) \subset L^{2}\left(\mathbb{R}^{d}\right)$. Invertibility is proved using the Fourier transform: let $D(B)=L^{2}\left(\mathbb{R}^{d}\right)$, and define a bounded operator $B: L^{2} \rightarrow L^{2}$ by

$$
\widehat{\mathrm{Bf}}(\xi)=\left(1+4 \pi^{2}|\xi|^{2}\right)^{-1} \widehat{f}(\xi) .
$$

One can check that $\operatorname{Ran}(B)=H^{2}\left(\mathbb{R}^{d}\right)=D(A)$. Notice $B A=i d_{H^{2}}, A B=$ $\operatorname{id}_{L^{2}}$. The second identity implies that $\operatorname{Ran}(A)=L^{2}=D(B)$.

## Adjoint of an (unbounded) operator

Call $A$ symmetric if

$$
\begin{equation*}
\langle A f, g\rangle=\langle f, A g\rangle, \quad \forall f, g \in D(A) \tag{17.1}
\end{equation*}
$$

Symmetry is a simpler concept than selfadjointness, which requires the operator and its adjoint to have the same domain, as we now explain.

First we define a subspace
$D\left(A^{*}\right)=\{f \in \mathcal{H}$ : the linear functional $g \mapsto\langle f, A g\rangle$ is bounded on $D(A)\}$.
Assume from now on that $A$ is densely defined, meaning $D(A)$ is dense in $\mathcal{H}$. Then for each $f \in D\left(A^{*}\right)$, the bounded linear functional $g \mapsto\langle f, A g\rangle$ is defined on a dense subspace of $\mathcal{H}$ and hence extends uniquely to a bounded linear functional on all of $\mathcal{H}$. By the Riesz Representation Theorem, that linear functional can be represented as the inner product of $g$ against a unique element of $\mathcal{H}$, which we call $\boldsymbol{A}^{*}$ f. Hence

$$
\begin{equation*}
\langle\mathrm{f}, \mathrm{Ag}\rangle=\left\langle A^{*} \mathrm{f}, \mathrm{~g}\right\rangle, \quad \forall \mathrm{f} \in \mathrm{D}\left(A^{*}\right), \quad \mathrm{g} \in \mathrm{D}(\mathrm{~A}) \tag{17.2}
\end{equation*}
$$

Clearly this operator $A^{*}: D\left(A^{*}\right) \rightarrow \mathcal{H}$ is linear. We call it the adjoint of $A$.
Lemma 17.1. If $A$ is a densely defined linear operator and $\lambda \in \mathbb{C}$, then $(A-\lambda)^{*}=A^{*}-\bar{\lambda}$.

We leave the (easy) proof to the reader. Implicit in the proof is that domains are unchanged by subtracting a constant: $D(A-\lambda)=D(A)$ and $\mathrm{D}\left(\left((A-\lambda)^{*}\right)=\mathrm{D}\left(A^{*}\right)\right.$.

The kernel of the adjoint complements the range of the original operator, as follows.

Proposition 17.2. If $A$ is a densely defined linear operator then $\overline{\operatorname{Ran}(\mathcal{A})} \oplus$ $\operatorname{ker}\left(\boldsymbol{A}^{*}\right)=\mathcal{H}$.

Proof. Clearly $\operatorname{ker}\left(\boldsymbol{A}^{*}\right) \subset \operatorname{Ran}(A)^{\perp}$, because if $\mathrm{f} \in \operatorname{ker}\left(\boldsymbol{A}^{*}\right)$ then $\boldsymbol{A}^{*} \mathrm{f}=0$ and so for all $g \in D(A)$ we have

$$
\langle\mathrm{f}, \mathrm{Ag}\rangle=\left\langle A^{*} \mathrm{f}, \mathrm{~g}\right\rangle=0
$$

To prove the reverse inclusion, $\operatorname{Ran}(A)^{\perp} \subset \operatorname{ker}\left(A^{*}\right)$, suppose $h \in \operatorname{Ran}(A)^{\perp}$. For all $g \in D(A)$ we have $\langle h, A g\rangle=0$. In particular, $h \in D\left(A^{*}\right)$. Hence

$$
\left\langle A^{*} \mathrm{~h}, \mathrm{~g}\right\rangle=\langle\mathrm{h}, \mathrm{Ag}\rangle=0 \quad \forall \mathrm{~g} \in \mathrm{D}(\mathrm{~A})
$$

and so from density of $D(A)$ we conclude $A^{*} h=0$. That is, $h \in \operatorname{ker}\left(A^{*}\right)$.
We have shown $\operatorname{Ran}(A)^{\perp}=\operatorname{ker}\left(A^{*}\right)$, and so (since the orthogonal complement is unaffected by taking the closure) $\overline{\operatorname{Ran}(A)}{ }^{\perp}=\operatorname{ker}\left(\mathcal{A}^{*}\right)$. The proposition follows immediately.

We will need later that the graph of the adjoint, $\left\{\left(f, \mathcal{A}^{*} f\right): f \in D\left(\mathcal{A}^{*}\right)\right\}$, is closed in $\mathcal{H} \times \mathcal{H}$.

Theorem 17.3. If $A$ is a densely defined linear operator then $A^{*}$ is a closed operator.

Proof. Suppose $f_{n} \in D\left(A^{*}\right)$ with $f_{n} \rightarrow f, A^{*} f_{n} \rightarrow g$, for some $f, g \in \mathcal{H}$. To prove the graph of $A^{*}$ is closed, we must show $f \in D\left(A^{*}\right)$ with $A^{*} f=g$.

For each $h \in D(A)$ we have

$$
\langle f, A h\rangle=\lim _{n}\left\langle f_{n}, A h\right\rangle=\lim _{n}\left\langle A^{*} f_{n}, h\right\rangle=\langle g, h\rangle .
$$

Thus the map $h \mapsto\langle f, A h\rangle$ is bounded for $h \in D(A)$. Hence $f \in D\left(A^{*}\right)$, and using the last calculation we see

$$
\left\langle A^{*} f, h\right\rangle=\langle f, A h\rangle=\langle g, h\rangle
$$

for all $h \in D(A)$. Density of the domain implies $A^{*} f=g$, as we wanted.

## Selfadjointness

Call $A$ selfadjoint if $A^{*}=A$, meaning $D\left(A^{*}\right)=D(A)$ and $A^{*}=A$ on their common domain.

Selfadjoint operators have closed graphs, due to closedness of the adjoint in Theorem 17.3. Thus:

Proposition 17.4. If a densely defined linear operator $\mathcal{A}$ is selfadjoint then it is closed.

The relation between selftadjointness and symmetry is clear:
Proposition 17.5. The densely defined linear operator $\mathcal{A}$ is selfadjoint if and only if it is symmetric and $\mathrm{D}(\mathrm{A})=\mathrm{D}\left(\mathrm{A}^{*}\right)$.

Proof. " $\Longrightarrow$ " If $A^{*}=A$ then the adjoint relation (17.2) reduces immediately to the symmetry relation (17.1).
" " The symmetry relation (17.1) together with the adjoint relation (17.2) implies that $\langle A f, g\rangle=\left\langle A^{*} f, g\right\rangle$ for all $f, g \in D(A)=D\left(A^{*}\right)$. Since $D(\mathcal{A})$ is dense in $\mathcal{H}$, we conclude $\mathcal{A f}=A^{*} f$.

For bounded operators, selfadjointness and symmetry are equivalent.

Lemma 17.6. If a linear operator $\mathcal{A}$ is bounded on $\mathcal{H}$, then it is selfadjoint if and only if it is symmetric.

Proof. Boundedness of $\mathcal{A}$ ensures that $\mathrm{D}\left(\mathcal{A}^{*}\right)=\mathcal{H}=\mathrm{D}(\mathcal{A})$, and so the adjoint relation (17.2) holds for all $\mathrm{f}, \mathrm{g} \in \mathcal{H}$. Thus $\mathcal{A}^{*}=\mathcal{A}$ is equivalent to symmetry.

## Example: selfadjointness for Schrödinger operators

Let $\mathrm{L}=-\Delta+\mathrm{V}$ be a Schrödinger operator with potential $\mathrm{V}(\mathrm{x})$ that is bounded and real-valued. Choose the domain to be $\mathrm{D}(\mathrm{L})=\mathrm{H}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)$ in the Hilbert space $L^{2}\left(\mathbb{R}^{d}\right)$. This Schrödinger operator is selfadjoint.

Proof. Density of $\mathrm{D}(\mathrm{L})$ follows from density in $\mathrm{L}^{2}$ of the smooth functions with compact support.

Our main task is to determine the domain of $L^{*}$. Fix $f, g \in H^{2}\left(\mathbb{R}^{d}\right)$. From the integration by parts formula $\langle f, \Delta g\rangle_{\mathrm{L}^{2}}=\langle\Delta f, g\rangle_{\mathrm{L}^{2}}$ (which one may alternatively prove with the help of the Fourier transform), one deduces that

$$
\left|\langle f, \Delta g\rangle_{\mathrm{L}^{2}}\right|=\left|\langle\Delta \mathrm{f}, \mathrm{~g}\rangle_{\mathrm{L}^{2}}\right| \leq\|f\|_{\mathrm{H}^{2}}\|\mathrm{~g}\|_{\mathrm{L}^{2}} .
$$

Also $\left|\langle f, V g\rangle_{\mathrm{L}^{2}}\right| \leq\|f\|_{\mathrm{L}^{2}}\|V\|_{\mathrm{L}^{\infty}}\|\mathrm{g}\|_{\mathrm{L}^{2}}$. Hence the linear functional $\mathrm{g} \mapsto\langle\mathrm{f}, \mathrm{Lg}\rangle_{\mathrm{L}^{2}}$ is bounded on $g \in D(L)$. Therefore $f \in D\left(L^{*}\right)$, which tells us $H^{2}\left(\mathbb{R}^{d}\right) \subset$ $\mathrm{D}\left(\mathrm{L}^{*}\right)$.

To prove the reverse inclusion, fix $f \in D\left(L^{*}\right)$. Then

$$
\left|\langle\mathrm{f}, \mathrm{Lg}\rangle_{\mathrm{L}^{2}}\right| \leq(\text { const. })\|\mathrm{g}\|_{\mathrm{L}^{2}}, \quad \forall \mathrm{~g} \in \mathrm{D}(\mathrm{~L})=\mathrm{H}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)
$$

Since the potential V is bounded, the last formula still holds if we replace V with 1, so that

$$
\left.\left|\langle\mathrm{f},(-\Delta+1) \mathrm{g}\rangle_{\mathrm{L}^{2}}\right| \leq \text { (const. }\right)\|\mathrm{g}\|_{\mathrm{L}^{2}}, \quad \forall \mathrm{~g} \in \mathrm{H}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)
$$

Taking Fourier transforms gives

$$
\left|\left\langle\widehat{\mathfrak{f}},\left(1+4 \pi^{2}|\xi|^{2}\right) \widehat{\boldsymbol{g}}\right\rangle_{\mathrm{L}^{2}}\right| \leq(\text { const. })\|\widehat{\boldsymbol{g}}\|_{\mathrm{L}^{2}}, \quad \forall \mathrm{~g} \in \mathrm{H}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)
$$

In particular, we may suppose $\widehat{g}=h \in C_{0}^{\infty}\left(\mathbb{R}^{d}\right)$, since every such $\widehat{g}$ gives $g \in H^{2}\left(\mathbb{R}^{d}\right)$. Hence

$$
\left.\left|\left\langle\left(1+4 \pi^{2}|\xi|^{2}\right) \widehat{f}, h\right\rangle_{\mathrm{L}^{2}}\right| \leq \text { (const. }\right)\|h\|_{\mathrm{L}^{2}}, \quad \forall \mathrm{~h} \in \mathrm{C}_{0}^{\infty}\left(\mathbb{R}^{\mathrm{d}}\right)
$$

Taking the supremum of the left side over all $h$ with $L^{2}$-norm equal to 1 shows that

$$
\left\|\left(1+4 \pi^{2}|\xi|^{2}\right) \widehat{f}\right\|_{L^{2}} \leq \text { (const.) }
$$

Hence $(1+|\xi|)^{2} \widehat{f} \in L^{2}\left(\mathbb{R}^{d}\right)$, which means $f \in H^{2}\left(\mathbb{R}^{d}\right)$. Thus $D\left(L^{*}\right) \subset H^{2}\left(\mathbb{R}^{d}\right)$.
Now that we know the domains of L and $L^{*}$ agree, we have only to check symmetry, and that is straightforward. When $f, g \in H^{2}\left(\mathbb{R}^{d}\right)$ we have

$$
\begin{aligned}
\langle\mathrm{Lf}, \mathrm{~g}\rangle & =-\langle\Delta \mathrm{f}, \mathrm{~g}\rangle_{\mathrm{L}^{2}}+\langle\mathrm{Vf}, \mathrm{~g}\rangle_{\mathrm{L}^{2}} \\
& =-\langle\mathrm{f}, \Delta \mathrm{~g}\rangle_{\mathrm{L}^{2}}+\langle\mathrm{f}, \mathrm{Vg}\rangle_{\mathrm{L}^{2}} \\
& =\langle\mathrm{f}, \mathrm{Lg}\rangle_{\mathrm{L}^{2}}
\end{aligned}
$$

where we integrated by parts and used that $\mathrm{V}(\mathrm{x})$ is real-valued.

## Chapter 18

## Spectra: discrete and continuous

## Goal

To develop the spectral theory of selfadjoint unbounded linear operators.
References [GustafsonSigal] Sections 2.4, 5.1
[HislopSigal] Chapters 1, 5, 7
[Rudin] Chapter 13

## Resolvent set, and spectrum

Let $\mathcal{A}$ be a densely defined linear operator on a complex Hilbert space $\mathcal{H}$, as in the preceding chapter. The operator $A-\lambda$ has domain $D(A)$, for each constant $\lambda \in \mathbb{C}$. Define the resolvent set
$\operatorname{res}(A)=\{\lambda \in \mathbb{C}: A-\lambda$ is invertible (has a bounded inverse defined on $\mathcal{H})\}$.
For $\lambda$ in the resolvent set, we call the inverse $(A-\lambda)^{-1}$ the resolvent operator.

The spectrum is defined as the complement of the resolvent set:

$$
\operatorname{spec}(A)=\mathbb{C} \backslash \operatorname{res}(A)
$$

For example, if $\lambda$ is an eigenvalue of $A$ then $\lambda \in \operatorname{spec}(\mathcal{A})$, because if $A f=\lambda f$ for some $\mathrm{f} \neq 0$, then $(A-\lambda) f=0$ and so $A-\lambda$ is not injective, and hence is not invertible.

Proposition 18.1 ([HislopSigal, Theorem 1.2]). The resolvent set is open, and hence the spectrum is closed.

We omit the proof.
The next result generalizes the fact that Hermitian matrices have only real eigenvalues.

Theorem 18.2. If $A$ is selfadjoint then its spectrum is real: $\operatorname{spec}(\mathcal{A}) \subset \mathbb{R}$.
Proof. We prove the contrapositive. Suppose $\lambda \in \mathbb{C}$ has nonzero imaginary part, $\operatorname{Im} \lambda \neq 0$. We will show $\lambda \in \operatorname{res}(A)$.

The first step is to show $A-\lambda$ is injective. For all $f \in D(A)$,

$$
\|(A-\lambda) f\|^{2}=\|A f\|^{2}-2(\operatorname{Re} \lambda)\langle f, A f\rangle+|\lambda|^{2}\|f\|^{2}
$$

and so

$$
\begin{align*}
\|(A-\lambda) f\|^{2} & \geq\|A f\|^{2}-2|\operatorname{Re} \lambda|\|f\|\|A f\|+|\lambda|^{2}\|f\|^{2} \\
& =(\|A f\|-|\operatorname{Re} \lambda|\|f\|)^{2}+|\operatorname{Im} \lambda|^{2}\|f\|^{2} \\
& \geq|\operatorname{Im} \lambda|^{2}\|f\|^{2} . \tag{18.1}
\end{align*}
$$

The last inequality implies that $A-\lambda$ is injective, using here that $|\operatorname{Im} \lambda|>0$. That is, $\operatorname{ker}(A-\lambda)=\{0\}$.

Selfadjointness $\left(A^{*}=A\right)$ now gives $\operatorname{ker}\left(A^{*}-\lambda\right)=0$, and so $\overline{\operatorname{Ran}(A-\lambda)}=$ $\mathcal{H}$ by Proposition 17.2. That is, $A-\lambda$ has dense range.

Next we show $\operatorname{Ran}(A-\lambda)=\mathcal{H}$. Let $g \in \mathcal{H}$. By density of the range, we may take a sequence $f_{n} \in D(A)$ such that $(A-\lambda) f_{n} \rightarrow g$. The sequence $f_{n}$ is Cauchy, in view of (18.1). Hence the sequence $\left(f_{n},(A-\lambda) f_{n}\right)$ is Cauchy in $\mathcal{H} \times \mathcal{H}$, and so converges to ( $f, g$ ) for some $f \in \mathcal{H}$. Note each ordered pair $\left(f_{n},(A-\lambda) f_{n}\right)$ lies in the graph of $A-\lambda$, and this graph is closed by Proposition 17.4 (relying here on selfadjointness again). Therefore ( $\mathrm{f}, \mathrm{g}$ ) belongs to the graph of $A-\lambda$, and so $g \in \operatorname{Ran}(A-\lambda)$. Thus $A-\lambda$ has full range.

To summarize: we have shown $A-\lambda$ is injective and surjective, and so it has an inverse operator

$$
(A-\lambda)^{-1}: \mathcal{H} \rightarrow \mathrm{D}(A) \subset \mathcal{H}
$$

This inverse is bounded with

$$
\left\|(A-\lambda)^{-1} \mathrm{~g}\right\| \leq|\operatorname{Im} \lambda|^{-1}\|\mathrm{~g}\|, \quad \forall \mathrm{g} \in \mathcal{H}
$$

by taking $f=(A-\lambda)^{-1} g$ in estimate (18.1). The proof is thus complete.

## Characterizing the spectrum

We will characterize the spectrum in terms of approximate eigenfunctions. Given a number $\lambda \in \mathbb{C}$ and a sequence $w_{n} \in D(A)$, consider three conditions:
(W1) $\left\|(A-\lambda) w_{n}\right\|_{\mathcal{H}} \rightarrow 0$ as $n \rightarrow \infty$,
(W2) $\left\|w_{\mathfrak{n}}\right\|_{\mathcal{H}}=1$,
(W3) $w_{n} \rightharpoonup 0$ weakly in $\mathcal{H}$ as $n \rightarrow \infty$.
(We considered these conditions in Chapter 15 for the special case of the Laplacian).

Condition (W1) says $w_{n}$ is an "approximate eigenfunction", and condition (W2) simply normalizes the sequence. These conditions characterize the spectrum, for a selfadjoint operator.

Theorem 18.3. If A is selfadjoint then

$$
\operatorname{spec}(A)=\left\{\lambda \in \mathbb{C}:(W 1) \text { and (W2) hold for some sequence } w_{n} \in D(A)\right\}
$$

Proof. " $\supset$ " Assume (W1) and (W2) hold for $\lambda$, and that $A-\lambda$ has an inverse defined on $\mathcal{H}$. Then for $f_{n}=(A-\lambda) w_{n}$ we find

$$
\frac{\left\|(A-\lambda)^{-1} f_{n}\right\|_{\mathcal{H}}}{\left\|f_{n}\right\|_{\mathcal{H}}}=\frac{\left\|w_{n}\right\|_{\mathcal{H}}}{\left\|(A-\lambda) w_{n}\right\|_{\mathcal{H}}} \rightarrow \infty
$$

as $n \rightarrow \infty$, by (W1) and (W2). Thus the inverse operator is not bounded, and so $\lambda \in \operatorname{spec}(\mathcal{A})$.
" $\subset$ " Assume $\lambda \in \operatorname{spec}(A)$, so that $\lambda$ is real by Theorem 18.2. If $\lambda$ is an eigenvalue, say with normalized eigenvector f , then we simply choose $w_{\mathrm{n}}=\mathrm{f}$ for each $\mathfrak{n}$, and (W1) and (W2) hold trivially.

Suppose $\lambda$ is not an eigenvalue. Then $\lambda-\lambda$ is injective, hence so is $(A-\lambda)^{*}$, which equals $A-\lambda$ by selfadjointness of $A$ and reality of $\lambda$. Thus $\operatorname{ker}\left((A-\lambda)^{*}\right)=\{0\}$, and so $\operatorname{Ran}(A-\lambda)$ is dense in $\mathcal{H}$ by Proposition 17.2.

Injectivity ensures that $(\mathcal{A}-\lambda)^{-1}$ exists on $\operatorname{Ran}(\lambda-\lambda)$. If it is unbounded there, then we may choose a sequence $f_{n} \in \operatorname{Ran}(A-\lambda)$ with $\left\|(A-\lambda)^{-1} f_{n}\right\|_{\mathcal{H}}=$ 1 and $\left\|f_{n}\right\|_{\mathcal{H}} \rightarrow 0$. Letting $w_{n}=(A-\lambda)^{-1} f_{n}$ gives (W1) and (W2) as desired. Suppose on the other hand that $(A-\lambda)^{-1}$ is bounded on $\operatorname{Ran}(A-\lambda)$. Then the argument in the proof of Theorem 18.2 shows that $\operatorname{Ran}(\mathcal{A}-\lambda)=\mathcal{H}$, which means $\lambda$ belongs to the resolvent set, and not the spectrum. Thus this case cannot occur.

## Discrete and continuous spectra

Define the discrete spectrum
$\operatorname{spec}_{\text {disc }}(A)$
$=\{\lambda \in \operatorname{spec}(A): \lambda$ is an isolated eigenvalue of $A$ having finite multiplicity $\}$,
where "isolated" means that some neighborhood of $\lambda$ in the complex plane intersects $\operatorname{spec}(A)$ only at $\lambda$. By "multiplicity" we mean the geometric multiplicity (dimension of the eigenspace); if $A$ is not selfadjoint then we should use instead the algebraic multiplicity [HislopSigal].

Next define the continuous spectrum

$$
\begin{aligned}
& \text { spec }_{\text {cont }}(A) \\
& =\left\{\lambda \in \mathbb{C}:(\mathrm{W} 1),(\mathrm{W} 2) \text { and (W3) hold for some sequence } w_{n} \in D(A)\right\} .
\end{aligned}
$$

The continuous spectrum lies within the spectrum, by Theorem 18.3. The characterization in that theorem required only (W1) and (W2), whereas the continuous spectrum imposes in addition the "weak convergence" condition (W3).

A Weyl sequence for $A$ and $\lambda$ is a sequence $w_{n} \in D(A)$ such that (W1), (W2) and (W3) hold. Thus the preceding definition says the continuous spectrum consists of $\lambda$-values for which Weyl sequences exist.

The continuous spectrum can contain eigenvalues that are not isolated ("imbedded eigenvalues") or which have infinite multiplicity.

A famous theorem of Weyl says that for selfadjoint operators, the entire spectrum is covered by the discrete and continuous spectra.

Theorem 18.4. If $A$ is selfadjoint then

$$
\operatorname{spec}(A)=\operatorname{spec}_{\text {disc }}(A) \cup \operatorname{spec}_{\text {cont }}(A) .
$$

(Further, the discrete and continuous spectra are disjoint.)
We omit the proof. See [HislopSigal, Theorem 7.2].

## Applications to Schrödinger operators

The continuous spectrum of the Laplacian $-\Delta$ equals $[0, \infty)$, and the spectrum contains no eigenvalues, as we saw in Chapter 15.

The hydrogen atom too has continuous spectrum $[0, \infty)$, with its Schrödinger operator $\mathrm{L}=-\Delta-2 /|x|$ on $\mathbb{R}^{3}$ having domain $\mathrm{H}^{2}\left(\mathbb{R}^{\mathrm{d}}\right) \subset \mathrm{L}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)$; see [Taylor, Section 8.7]. The discrete spectrum $\left\{-1 / n^{2}: n \geq 1\right\}$ of the hydrogen atom was stated in Chapter 3.

As the hydrogen atom example suggests, potentials vanishing at infinity generate continuous spectrum that includes all nonnegative numbers:

Theorem 18.5. Assume $\mathrm{V}(\mathrm{x})$ is real-valued, continuous, and vanishes at infinity $(\mathrm{V}(\mathrm{x}) \rightarrow 0$ as $|\mathrm{x}| \rightarrow \infty)$.

Then the Schrödinger operator $-\Delta+\mathrm{V}$ is selfadjoint (with domain $\mathrm{H}^{2}\left(\mathbb{R}^{\mathrm{d}}\right) \subset$ $\mathrm{L}^{2}\left(\mathbb{R}^{\mathrm{d}}\right)$ ) and has continuous spectrum $=[0, \infty)$.

For a proof see [HislopSigal, Corollary 14.10], where a stronger theorem is proved that covers also the Coulomb potential $-2 /|x|$ for the hydrogen atom. Note the Coulomb potential vanishes at infinity but is discontinuous at the origin, where it blows up. The stronger version of the theorem requires (instead of continuity and vanishing at infinity) that for each $\varepsilon>0$, the potential $\mathrm{V}(\mathrm{x})$ be decomposable as $\mathrm{V}=\mathrm{V}_{2}+\mathrm{V}_{\infty}$ where $\mathrm{V}_{2} \in \mathrm{~L}^{2}$ and $\left\|\mathrm{V}_{\infty}\right\|_{L^{\infty}}<\varepsilon$. This decomposition can easily be verified for the Coulomb potential, by "cutting off" the potential near infinity.

Theorem 18.5 implies that any isolated eigenvalues of L must lie on the negative real axis (possibly accumulating at 0 ). For example, the -2 sech $^{2}$ potential in Chapter 16 generates a negative eigenvalue at -1 .

## Connection to generalized eigenvalues and eigenfunctions

Just as the discrete spectrum is characterized by eigenfunctions in $L^{2}$, so the full spectrum is characterized by existence of a generalized eigenfunctions that grows at most polynomially at infinity.

Theorem 18.6. Assume $\mathrm{V}(\mathrm{x})$ is real-valued and bounded on $\mathbb{R}^{\mathrm{d}}$. Then the Schrödinger operator $-\Delta+\mathrm{V}$ has spectrum

$$
\begin{aligned}
& \operatorname{spec}(-\Delta+\mathrm{V})= \\
& \quad \text { closure of }\{\lambda \in \mathbb{C}:(-\Delta+\mathrm{V}) \mathbf{u}=\lambda \mathfrak{u} \text { for some polynomially bounded } \mathfrak{u}\} .
\end{aligned}
$$

We omit the proof; see [GustafsonSigal, Theorem 5.22].

## Further reading

A wealth of information on spectral theory, especially for Schrödinger operators, can be found in the books [GustafsonSigal, HislopSigal, ReedSimon2, ReedSimon4].

## Chapter 19

## Discrete spectrum revisited

## Goal

To fit the discrete spectral Theorem 4.1 (from Part I of the course) into the spectral theory of selfadjoint operators and, in particular, to prove the absence of continuous spectrum in that situation.

## Discrete spectral theorem

The discrete spectral Theorem 4.1 concerns a symmetric, elliptic, bounded sesquilinear form $\mathfrak{a}(u, v)$ on an infinite dimensional Hilbert space $\mathcal{K}$, where $\mathcal{K}$ imbeds compactly and densely into the Hilbert space $\mathcal{H}$. The theorem guarantees existence of an ONB for $\mathcal{H}$ consisting of eigenvectors of $a$ :

$$
a\left(u_{j}, v\right)=\gamma_{j}\left\langle u_{j}, v\right\rangle_{\mathcal{H}} \quad \forall v \in \mathcal{K},
$$

where the eigenvalues satisfy

$$
0<\gamma_{1} \leq \gamma_{2} \leq \gamma_{3} \leq \cdots \rightarrow \infty
$$

We want to interpret these eigenvalues as the discrete spectrum of some selfadjoint, densely defined linear operator on $\mathcal{H}$. By doing so, we will link the discrete spectral theory in Part I of the course with the spectral theory of unbounded operators in Part II.

Our tasks are to identify the operator $A$ and its domain, to prove $A$ is symmetric, to determine the domain of the adjoint, to conclude selfadjointness, and finally to show that the spectrum of $A$ consists precisely of the eigenvalues $\gamma_{j}$.

## Operator $A$ and its domain

In the proof of Theorem 4.1 we found a bounded, selfadjoint linear operator B: $\mathcal{H} \rightarrow \mathcal{K} \subset \mathcal{H}$ with eigenvalues $1 / \gamma_{j}$ and eigenvectors $u_{j}$ :

$$
\mathrm{B} u_{j}=\frac{1}{\gamma_{j}} u_{j}
$$

We showed $B$ is injective (meaning its eigenvalues are nonzero). Notice $B$ has dense range because its eigenvectors $\mathfrak{u}_{j}$ span $\mathcal{H}$.
(Aside. This operator B relates to the sesquilinear form a by satisfying $\mathrm{a}(\mathrm{Bf}, v)=\langle\mathrm{f}, v\rangle_{\mathcal{H}}$ for all $v \in \mathcal{K}$. We will not need that formula below.)

Define

$$
\mathrm{A}=\mathrm{B}^{-1}: \operatorname{Ran}(\mathrm{B}) \rightarrow \mathcal{H}
$$

Then $A$ is a linear operator, and its domain

$$
\mathrm{D}(\mathrm{~A})=\operatorname{Ran}(\mathrm{B})
$$

is dense in $\mathcal{H}$.

## Symmetry of A

Let $u, v \in D(A)$. Then

$$
\begin{aligned}
\langle A u, v\rangle_{\mathcal{H}} & =\langle A u, B A v\rangle_{\mathcal{H}} & & \text { since } B A=\mathrm{Id}, \\
& =\langle B A u, A v\rangle_{\mathcal{H}} & & \text { since } B \text { is selfadjoint, } \\
& =\langle u, A v\rangle_{\mathcal{H}} & & \text { since } B A=\mathrm{Id} .
\end{aligned}
$$

## Domain of the adjoint

First we show $\mathrm{D}(A) \subset \mathrm{D}\left(A^{*}\right)$. Let $u \in \mathrm{D}(A)$. For all $v \in \mathrm{D}(A)$ we have

$$
\begin{array}{rlr}
\left|\langle u, A v\rangle_{\mathcal{H}}\right| & =\left|\langle A u, v\rangle_{\mathcal{H}}\right| & \text { by symmetry } \\
& \leq\|A u\|_{\mathcal{H}}\|v\|_{\mathcal{H}} &
\end{array}
$$

Hence the functional $v \mapsto\langle u, A v\rangle_{\mathcal{H}}$ is bounded on $D(A)$ with respect to the $\mathcal{H}$-norm, so that $u$ belongs to the domain of the adjoint $A^{*}$.

Next we show $\mathrm{D}\left(A^{*}\right) \subset \mathrm{D}(A)$. Let $\mathfrak{u} \in \mathrm{D}\left(A^{*}\right) \subset \mathcal{H}$. We have

$$
\left|\langle u, A v\rangle_{\mathcal{H}}\right| \leq \text { (const.) }\|v\|_{\mathcal{H}} \quad \forall v \in \mathrm{D}(\mathcal{A})=\operatorname{Ran}(\mathrm{B})
$$

Writing $v=\mathrm{Bg}$ gives

$$
\left|\langle u, g\rangle_{\mathcal{H}}\right| \leq(\text { const. })\|\mathrm{Bg}\|_{\mathcal{H}} \quad \forall \mathrm{g} \in \mathcal{H} .
$$

One can express $u$ in terms of the ONB as $u=\sum_{j} d_{j} u_{j}$. Fix $J \geq 1$ and choose $g=\sum_{j=1}^{J} \gamma_{j}^{2} d_{j} u_{j} \in \mathcal{H}$, so that $\mathrm{Bg}=\sum_{j=1}^{J} \gamma_{j} d_{j} u_{j}$. We deduce from the last inequality that

$$
\sum_{j=1}^{J} \gamma_{j}^{2}\left|\mathrm{~d}_{j}\right|^{2} \leq(\text { const. })\left(\sum_{j=1}^{J} \gamma_{j}^{2}\left|\mathrm{~d}_{j}\right|^{2}\right)^{1 / 2},
$$

and so

$$
\sum_{j=1}^{J} \gamma_{j}^{2}\left|\mathrm{~d}_{\mathrm{j}}\right|^{2} \leq(\text { const. })^{2}
$$

Letting $\mathrm{J} \rightarrow \infty$ implies that

$$
\sum_{j} \gamma_{j}^{2}\left|d_{j}\right|^{2} \leq(\text { const. })^{2}
$$

and so the sequence $\left\{\gamma_{j} d_{j}\right\}$ belongs to $\ell^{2}$. Put $f=\sum_{j} \gamma_{j} d_{j} u_{j} \in \mathcal{H}$. Then $B f=\sum_{j} d_{j} u_{j}=u$, and so $u \in \operatorname{Ran}(B)=D(A)$, as desired.

## Selfadjointness, and discreteness of the spectrum

Theorem 19.1. A is selfadjoint, with domain

$$
D(A)=\operatorname{Ran}(B)=\left\{\sum_{j} \gamma_{j}^{-1} c_{j} u_{j}:\left\{c_{j}\right\} \in \ell^{2}\right\}
$$

Furthermore, $\operatorname{spec}(A)=\operatorname{spec}_{\text {disc }}(A)=\left\{\gamma_{j}: \mathfrak{j} \geq 1\right\}$.
Proof. We have shown above that $A$ is symmetric and $D\left(A^{*}\right)=D(A)$, which together imply that $A$ is selfadjoint.

We will show that if

$$
\lambda \in \mathbb{C} \backslash\left\{\gamma_{1}, \gamma_{2}, \gamma_{3}, \ldots\right\}
$$

then $A-\lambda$ is invertible, so that $\lambda$ belongs to the resolvent set. Thus the spectrum consists of precisely the eigenvalues $\gamma_{j}$. Note each eigenvalue has finite
multiplicity by Theorem 4.1, and is isolated from the rest of the spectrum; hence $A$ has purely discrete spectrum.

The inverse of $A-\lambda$ can be defined explicitly, as follows. Define a bounded operator $\mathrm{C}: \mathcal{H} \rightarrow \mathcal{H}$ on $\mathrm{f}=\sum_{j} \mathrm{c}_{\mathrm{j}} \mathrm{u}_{\mathrm{j}} \in \mathcal{H}$ by

$$
\mathrm{Cf}=\sum_{j}\left(\gamma_{j}-\lambda\right)^{-1} c_{j} u_{j}
$$

where we note that $\left(\gamma_{j}-\lambda\right)^{-1}$ is bounded for all $\mathfrak{j}$, and in fact approaches 0 as $\mathfrak{j} \rightarrow \infty$, because $\left|\gamma_{j}-\lambda\right|$ is never zero and tends to $\infty$ as $\mathfrak{j} \rightarrow \infty$. This new operator has range $\operatorname{Ran}(C)=\operatorname{Ran}(B)$, because $\left(\gamma_{j}-\lambda\right)^{-1}$ is comparable to $\gamma_{j}^{-1}$ (referring here to the characterization of $\operatorname{Ran}(B)$ in Theorem 19.1). Thus $\operatorname{Ran}(\mathrm{C})=\mathrm{D}(\mathrm{A})$.

Clearly $(A-\lambda) C f=f$ by definition of $A$, and so $\operatorname{Ran}(A-\lambda)=\mathcal{H}$. Similarly one finds that $C(A-\lambda) u=u$ for all $u \in D(A)$. Thus $C$ is the inverse operator of $A-\lambda$. Because $C$ is bounded on all of $\mathcal{H}$ we conclude $A-\lambda$ is invertible, according to the definition in Chapter 17.

## Example: Laplacian on a bounded domain

To animate the preceding theory, let us consider the Laplacian on a bounded domain $\Omega \subset \mathbb{R}^{\mathrm{d}}$, with Dirichlet boundary conditions. We work with the Hilbert spaces

$$
\mathcal{H}=\mathrm{L}^{2}(\Omega), \quad \mathcal{K}=\mathrm{H}_{0}^{1}(\Omega)
$$

and the sesquilinear form

$$
\mathrm{a}(\mathrm{u}, v)=\int_{\Omega} \nabla \mathfrak{u} \cdot \nabla v \mathrm{dx}+\int_{\Omega} \mathfrak{u v} \mathrm{dx}=\langle\mathfrak{u}, v\rangle_{\mathrm{H}^{1}}
$$

which in Chapter 5 gave eigenfunctions satisfying $(-\Delta+1) u=(\lambda+1) u$ weakly. In this setting, $u=\operatorname{Bf}$ means that $(-\Delta+1) u=f$ weakly. Note $B: L^{2}(\Omega) \rightarrow H_{0}^{1}(\Omega)$, and recall that $A=B^{-1}$.

Proposition 19.2. The domain of the operator $\mathcal{A}$ contains $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$, and

$$
A=-\Delta+1
$$

on $\mathrm{H}^{2}(\Omega) \cap \mathrm{H}_{0}^{1}(\Omega)$.
Furthermore, if $\partial \Omega$ is smooth then $\mathrm{D}(\mathrm{A})=\mathrm{H}^{2}(\Omega) \cap \mathrm{H}_{0}^{1}(\Omega)$, in which case $A=-\Delta+1$ on all of its domain.

Proof. For all $u \in H^{2}(\Omega) \cap H_{0}^{1}(\Omega), v \in H_{0}^{1}(\Omega)$, we have

$$
\begin{aligned}
\langle\mathfrak{u}, v\rangle_{\mathrm{H}^{1}} & =\langle-\Delta \mathfrak{u}+\mathfrak{u}, v\rangle_{\mathrm{L}^{2}} & & \text { by parts } \\
& =\langle\mathrm{B}(-\Delta \mathfrak{u}+\mathfrak{u}), v\rangle_{\mathrm{H}^{1}} & & \text { by definition of } \mathrm{B} .
\end{aligned}
$$

Since both $u$ and $B(-\Delta u+u)$ belong to $H_{0}^{1}(\Omega)$, and $v \in H_{0}^{1}(\Omega)$ is arbitrary, we conclude from above that $u=B(-\Delta u+u)$. Therefore $u \in \operatorname{Ran}(B)=$ $D(A)$, and so $H^{2}(\Omega) \cap H_{0}^{1}(\Omega) \subset D(A)$.

Further, we find $A \mathfrak{u}=-\Delta u+u$ because $B=A^{-1}$, and so

$$
A=-\Delta+1 \quad \text { on } H^{2}(\Omega) \cap H_{0}^{1}(\Omega)
$$

Finally we note that if $\partial \Omega$ is $C^{2}$-smooth then by elliptic regularity the weak solution $u$ of $(-\Delta+1) u=f$ belongs to $H^{2}(\Omega)$, so that $\operatorname{Ran}(B) \subset$ $H^{2}(\Omega) \cap H_{0}^{1}(\Omega)$. Thus

$$
\mathrm{D}(\mathrm{~A})=\mathrm{H}^{2}(\Omega) \cap \mathrm{H}_{0}^{1}(\Omega)
$$

when $\partial \Omega$ is smooth enough. In that case $A=-\Delta+1$ on all of its domain.

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