The Pólya–Szegő Conjecture on Polygons: A Numerical Approach

By

LOGAN REED

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Written under the direction of

Siqi Fu

And approved by

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Dr. One

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Dr. Two

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Dr. Three

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Dr. Four

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TEST PAGE

THESIS ABSTRACT

Write Abstract

TODO SOMETHING ABOUT THE ABSTRACT THAT IS ABOUT THIS

LONG OR SO

by LOGAN REED

Thesis Director:

Siqi Fu

The topic that I chose to explore for this thesis is a study of the eigenvalues of the Dirichlet Laplacian on a two dimensional domain and . . .

Acknowledgment

I would be remiss if I did not take a moment to express appreciation for all of the people who have helped me through this process.

Contents

Chapter 1

Introduction

1.1 Physical Motivation

Physical drums consist of a rigid shell with a membrane which produces sound when hit. A similar thing can be "created" in a pure mathematical setting by studing specific partial differential equations over a closed region. Specifically, the frequencies of the drum membrane corresponds to the eigenvalues of the Dirichlet Laplacian. This construction makes it possible to "hear" drums where the shape of the drumhead is any closed and simple curve. To do this, we begin with the shape of our drumhead, which is a region of the real plane bounded by piecewise smooth curves. Since we wish to emulate the physical properties of a drum, we want to define some system that models the vibration of the drum membrane which produces the sound. This is done using the wave equation over our boundary, with the boundary condition that the function is 0 on the boundary. The Dirichlet Laplacian allows us to model the physical vibration of the drum, and we can calculate its eigenvalues to find the fundamental frequency and overtones. Thus, by modeling the vibration of the drumhead using the Dirichlet Laplacian and studying its properties, we can produce "sound" via the eigenvalues.

In 1877, Lord Rayleigh conjectured the following [6]

> If the area of a membrane be given, there must evidently be some form of boundary for which the pitch (of the principal tone) is the gravest possible, and this form can be no other than the circle.

In 1923, Faber published a proof which was followed by an independent proof by Krahn in 1925 [4].

Theorem 1.1.1 (Faber-Krahn)**.** *Let c be a positive number and B the ball with volume c. Then,*

$$\lambda_1(B) = \min\left\{\lambda_1(\Omega), \Omega \text{ open subset of } \mathbb{R}^N, |\Omega| = c\right\}.$$

From the Faber-Krahn inequality, we know that for any drumhead with a given area, the circle is the one with the lowest tone. In 1951, Polya and Szego conjectured that a similar statement holds for drumheads with a polygonal shape [5]. This conjecture has been shown to be true for 3 and 4 sided polygons, but remains unproven for any other number of sides.

There are two main hurdles that are halting progress on this conjecture. The first is that the tools that were used to prove both Lord Rayleighs conjecture as well as the small cases for the Polya-Szego conjecture are not available when the number of sides is greater than four. The main tool that is used is called Steiner Symmetrization, and when there are more than four sides this symmetrization method creates additional sides at each step.

The purpose of this paper is to show a specific method for running numerical approximations to suggest that this conjecture is indeed true. This is done using a method based on fundamental solutions [3]. Specifically, we consider all functions that satisfy the Laplace Equation and then solve for the linear coefficients using the boundary conditions. Once this is done we use gradient descent to find the polygon with the minimum first eigenvalue, which is equivalent to the first fundamental tone of the drum.

Consider a homogeneous elastic drumhead, or membrane, stretched over a rigid frame. We will represent the frame as a domain $\Omega \subset \mathbb{R}^2$. Take the function $u(x, y, t)$ to be the vertical displacement of the membrane from its resting position. Then for any disk $D \subset \Omega$, Newton's second law of motion states that

$$\int_{\partial D} T \frac{\partial u}{\partial \mathbf{n}} \, dS = \int_D \rho u_{tt} \, dA$$

where T is the constant tension, ρ is the density constant, and $\mathbf{n}$ is the outward normal of the boundary. By the divergence theorem, we have

$$\int_D T \Delta u \, dA = \int_D \rho u_{tt} \, dA$$

where Δ is the Laplace operator. From this we can get the wave equation on Ω

$$u_{tt} = c^2 \Delta u$$

where we define u to be 0 on the boundary and where $c = \sqrt{T/\rho}$. We can solve this wave equation using $u(x, y, t) = T(t)V(x, y)$ which gives us

$$\frac{T''}{c^2 T} = \frac{\Delta V}{V} = -\lambda$$

and finally we have reduced our problem to the Dirichlet Laplacian

$$\Delta V = -\lambda V$$

where V on the boundary is zero.

Chapter 2

Background

2.1 Measure Theory

Definition 2.1.1. *A Measure space is a pair (X, A) where X is a non-empty set and A is a σ-algebra of subsets of X. That is, A satisfies the following*

1. *$\emptyset \in A$*

2. *For countably many $A_j \in A$, $\cup A_j \in A$*

3. *If $B \in A$, then $X - B \in A$.*

Elements in A are called *measurable sets*. A *measure* μ on (X, A) is a non-negative function $\mu : A \to [0, \infty]$ such that $\mu(\emptyset) = 0$ and $\mu(\cup A_j) = \Sigma \mu(A_j)$ for any countably many, mutually disjoint $A_j \in A$. μ is said to be *finite* if $\mu(X) < \infty$, and μ is σ-finite if X is a countable union of sets in A with finite measures. A property is said to hold *almost everywhere* on a set A if it holds on A save a subset with zero measure.

A function $f : X \to [-\infty, \infty]$ is *measurable* if $\{x \in X : f(x) < \alpha\}$ is measurable for all $\alpha \in \mathbb{R}$. A *simple function* is a function of the form

$$f = \sum_{j=1}^{m} a_j \chi_{A_j},$$

where χ_S is the characteristic function on the set S, $a_j \in \mathbb{R}$, and $A_j \in A$.

Theorem 2.1.1 (Simple Function Approximation Theorem)**.** *Let $f : X \to [-\infty, \infty]$ be a measurable function. If f is non-negative, then there exists an increasing sequence of simple functions ϕ_j such that $0 \leq \phi_j \leq f$ and $\lim_{j \to \infty} \phi_j(x) = f(x)$. If f is bounded, then there exists a sequence of simple functions ϕ_j such that $\phi_j \to f$ uniformly on X.*

Proof.

Write Proof: Fu Notes and Evans Appendix

$\square$

We will assume our measure spaces (X, A, μ) are complete. That is, if $B \subset N$ and $\mu(N) = 0$ then $B \in A$. The *integration* with respect to μ is defined in the following way. We first define the integral for non-negative simple functions. Let $\phi = \sum_{j=1}^{m} a_j \chi_{A_j} \geq 0$, and define

$$\int_X \phi \, d\mu = \sum_{j=1}^{m} a_j \mu(A_j)$$

where we use the convention that $0 \cdot \infty = 0$. We define the integral for non-negative measurable functions as

$$\int_X f \, d\mu = \sup \left\{ \int_X \phi \, d\mu; 0 \leq \phi \leq f, \phi \text{ simple} \right\}.$$

For a measurable function $f : X \to [-\infty, \infty]$ we write $f^+ = \max\{f, 0\}$ and $f^- = \max\{-f, 0\}$, and we define

$$\int_X f \, d\mu = \int_X f^+ \, d\mu - \int_X f^- \, d\mu$$

given at least one of the integrals on the right hand side is finite. When f is complex-valued we define the integral by integrating the real and complex parts seperately. When $\int_X |f| \, d\mu < \infty$ we say f is *integrable*.

Definition 2.1.2. *If $0 < p < \infty$ and if f is a complex measurable function on X, define*

$$\|f\|_p = \left\{ \int_X |f|^p \, d\mu \right\}^{\frac{1}{p}}$$

and let $L^p(\mu)$ consist of all f for which $\|f\|_p < \infty$.

For $f : X \to \mathbb{C}$, the *support* of f is defined as $\text{supp}(f) = \overline{\{x \in X; f(x) \neq 0\}}$. Denote by $C_c(X)$ the family of continuous functions on X with compact support and by $C_0(X)$ the family of continuous functions that vanish at infinity.

2.2 Functional Analysis

Definition 2.2.1. *A complex linear space $\mathbb{H}$ is called a normed linear space if there exists a map $||\cdot|| : \mathbb{H} \to \mathbb{R}^+$ such that for any $x, y \in \mathbb{H}$ and $\lambda \in \mathbb{C}$,*

1. *$||\lambda x|| = |\lambda| ||x||$*

2. *$||x + y|| \leq ||x|| + ||y||$*

3. *$||x|| \geq 0$, and $||x|| = 0$ if and only if $x = 0$*

We call this map a norm

Definition 2.2.2. *A Banach space X is a complete, normed linear space.*

Definition 2.2.3. *We say X is seperable if X contains a countable dense subset.*

Definition 2.2.4. *A complex linear space $\mathbb{H}$ is called an inner product space with inner product $\langle \cdot, \cdot \rangle : \mathbb{H} \times \mathbb{H} \to \mathbb{C}$ if for any $x, y, z \in \mathbb{H}$ and $\lambda \in \mathbb{C}$,*

1. *$\langle \lambda x, y \rangle = \lambda \langle x, y \rangle$*

2. *$\langle x, y \rangle = \overline{\langle y, x \rangle}$*

3. *$\langle x + y, z \rangle = \langle x, z \rangle + \langle y, z \rangle$*

4. *$\langle x, x \rangle \geq 0$, and $\langle x, x \rangle = 0$ if and only if $x = 0$.*

For an inner product $\langle, \rangle$, the associated *norm* is $||u|| := \langle u, u \rangle^{\frac{1}{2}}$ for $u \in \mathbb{H}$. One could verify both of the definitions of a *norm* via the Cauchy-Schwarz inequality. We say that two elements $u, v \in \mathbb{H}$ are *orthogonal* if $\langle u, v \rangle = 0$. A *countable basis* $\{w_k\}_{k=1}^{\infty} \subset \mathbb{H}$ is *orthonormal* if the elements are pairwize orthogonal and the norm of each element is one.

Definition 2.2.5. *A Hilbert space $\mathbb{H}$ is a Banach space endowed with an inner product which generates the norm.*

For the remainder of this paper all Hilbert spaces will be assumed to be seperable.

Let X, Y be real Banach spaces.

Definition 2.2.6. *A mapping $A : X \to Y$ is a linear operator provided*

$$A(au + bv) = aAu + bAv$$

for all $u, v \in X$ and $a, b \in \mathbb{R}$.

Definition 2.2.7. *A linear operator $A : X \to Y$ is bounded if*

$$||A|| := \sup \{||Au||; ||u|| \leq 1\} \leq \infty.$$

Definition 2.2.8. *A linear operator $A : X \to Y$ is closed if whenever $u_k \to u$ in X and $Au_k \to v$ in Y, then $Au = v$*

finish once I've done the important proofs

2.3   PDEs

Definition 2.3.1 (Laplacian)**.**

$$-\Delta u := -\sum_{i=1}^{N} \frac{\partial^2 u}{\partial x_i^2}.$$

Definition 2.3.2. *The Dirichlet Laplacian is the Laplace Operator subject to Dirichlet boundary conditions. That is, we call u a solution to the Dirichlet Laplacian if u is a solution to*

$$\begin{cases} \Delta u = \lambda u & in \ \Omega \\ \ \ u = 0 & on \ \partial\Omega \end{cases}.$$

Should I write in this level or more akin to Henrots book. e.g. do I talk about the solution of Dirichlet Laplacian as actually being the unique solution to a variational problem. Because it is self contained there are proofs on both ends of the spectrum

2.4 Tools

Definition 2.4.1 (Schwarz Rearrangement)**.** *For any measurable set ω in $\mathbb{R}^N$, we denote by ω^* the ball of same volume as ω. If u is a non-negative measurable function defined on a measurable set Ω and vanishing on its boundary $\partial\Omega$, we denote by $\Omega(c) = \{x \in \Omega \,|\, u(x) \geq c\}$ its level sets. The Schwarz rearrangement of u is the function u^* defined on Ω^* by*

$$u^*(x) = \sup\{c/x \in \Omega(c)^*\}.$$

Without loss of generality, we fix the hyperplane of symmetry to be $x_N = 0$. Let $N \geq 2$ and $\Omega \subset \mathbb{R}^N$ be a measurable set. We denote by Ω' the projection of Ω on $\mathbb{R}^{N-1}$, and for every $x' \in \mathbb{R}^{N-1}$ we denote by $\Omega(x')$ the projection of Ω with $\{x'\} \times \mathbb{R}$.

Definition 2.4.2 (Steiner Symmetrization)**.** *Let $\Omega \subset \mathbb{R}^N$ be measurable. Then the set*

$$\Omega^* := \left\{x = (x', x_N) : -\frac{1}{2}|\Omega(x')| < x_N < \frac{1}{2}|\Omega(x')|, x' \in \Omega'\right\}$$

is the Steiner symmetrization of Ω with respect to the hyperplane $x_N = 0$.

Theorem 2.4.1. *Let Ω be a measurable set and u be a non-negative measurable function defined on Ω and vanishing on its boundary $\partial\Omega$. Let ϕ be any measurable function defined on $\mathbb{R}^+$ with values in $\mathbb{R}$, then*

$$\int_\Omega \phi(u(x)) \, \mathrm{d}x = \int_{\Omega^*} \phi(u^*(x)) \, \mathrm{d}x.$$

Theorem 2.4.2 (Pólya's Inequality)**.** *Let Ω be an open set and u a non-negative function belonging to the Sobolev space $H_0^1(\Omega)$. Then $u^* \in H_0^1(\Omega^*)$ and*

$$\int_\Omega |\nabla u(x)|^2 \, \mathrm{d}x \geq \int_{\Omega^*} |\nabla u^*(x)|^2 \, \mathrm{d}x.$$

Chapter 3

Eigenvalues of the Dirichlet Laplacian

3.1 Definition

Definition 3.1.1 (Rayleigh Quotient)**.** *For an operator L, we define the Rayleigh quotient to be*

$$R_L[v] := \frac{\sum_{i,j=1}^{N} \int_\Omega a_{ij}(x) \frac{\partial v}{\partial x_i} \frac{\partial v}{\partial x_j} \, \mathrm{d}x + \int_\Omega a_0(x) v^2(x) \, \mathrm{d}x}{\int_\Omega v(x)^2 \, \mathrm{d}x}.$$

This is used to express the first eigenvalue of the Dirichlet Laplacian in the following way

$$\lambda_1(\Omega) = \min_{v \in H_0^1(\Omega), v \neq 0} \frac{\int_\Omega |\nabla v(x)|^2 \, \mathrm{d}x}{\int_\Omega v(x)^2 \, \mathrm{d}x}.$$

Theorem 3.1.1. *Let Ω be a bounded open set. We assume that $\lambda_k'(\Omega)$ is simple. Then, the functions $t \to \lambda_k(t), t \to u_t \in L^2(\mathbb{R}^N)$ are differentiable at $t = 0$ with*

$$\lambda_k'(0) := -\int_\Omega \mathrm{div}(|\nabla u|^2 V) \, \mathrm{d}x.$$

If, moreover, Ω is of class C^2 or if Ω is convex, then

$$\lambda_k'(0) := -\int_\Omega \left(\frac{\partial u}{\partial n}\right)^2 V.n \, \mathrm{d}\sigma$$

and the derivative u' of u_t is the solution of

$$\begin{cases} -\Delta u' = \lambda_k u' + \lambda_k' u & \text{in} \Omega \\ \\ u' = -\frac{\partial u}{\partial n} V.n & \text{on} \partial\Omega \\ \\ \int_\Omega u u' \, \mathrm{d}\sigma = 0. \end{cases}$$

Theorem 3.1.2. *Each eigenvalue is real. Furthermore, if we repeat each eigenvalue according to its (finite) multiplicity, we have $\Sigma = \{\lambda_k\}_{k=1}^{\infty}$ where Σ is the set of eigenvalues, $0 < \lambda_1 \leq \lambda_2 \leq \lambda_3 \leq \ldots$, and $\lambda_k \to \infty$ as $k \to \infty$.*

Proof.

Pg 335 Evans. Ask Dr. Fu about elementary proof for first statement

3.2 Known Results

Rework Sections for this chapter. Known Results isn't descriptive

Theorem 3.2.1. *The Laplacian is invariant under orthogonal transformations.*

Proof. Let $f \in C_c^2(\mathbb{R}^N)$ and let A be an orthogonal $n \times n$ matrix over $\mathbb{R}$. Also, let $x = (x_1, x_2, \ldots, x_N)$. Since A is orthogonal, $\sum_{j=1}^{N} a_{ij} a_{kj} = \delta_{ik}$ where δ_{ik} is the Kronecker Delta function. So we have

$$(f \circ A)(x) = f\left(\sum_{i=1}^{N} a_{1i} x_i, \ldots, \sum_{i=1}^{N} a_{di} x_i\right).$$

Take $z_i = g_i(x_1, x_2, \ldots, x_N) = \sum_{k=1}^{N} a_{ik} x_k$. From a direct application of the chain rule we obtain

$$\frac{d}{dx_j}(f \circ A)(x) = \sum_{k=1}^{N} a_{kj} \cdot (\partial_k f)(\sum_{i=1}^{N} a_{1i} x_i, \ldots, \sum_{i=1}^{N} a_{di} x_i).$$

Further, by taking $\partial_k f$ in place of f, we obtain

$$\frac{d^2}{dx_j^2}(f \circ A)(x) = \sum_{k=1}^{N} a_{kj} \sum_{\ell=1}^{N} a_{\ell j}(\partial_\ell \partial_k f)(\sum_{i=1}^{N} a_{1i} x_i, \ldots, \sum_{i=1}^{N} a_{di} x_i).$$

With all of these pieces in place, we have the following

$$\Delta(f \circ A)(x) = \sum_{j=1}^{N} \frac{d^2}{dx_j^2}(f \circ A)(x)$$

$$= \sum_{j=1}^{N} \sum_{k=1}^{N} a_{kj} \sum_{\ell=1}^{N} a_{\ell j}(\partial_\ell \partial_k f)(\sum_{i=1}^{N} a_{1i}x_i, \ldots, \sum_{i=1}^{N} a_{di}x_i)$$

$$= \sum_{k,\ell=1}^{N} \left(\sum_{j=1}^{N} a_{kj}a_{\ell j} \right) (\partial_\ell \partial_k f)(\sum_{i=1}^{N} a_{1i}x_i, \ldots, \sum_{i=1}^{N} a_{di}x_i)$$

$$= \sum_{k,\ell=1}^{N} \delta_{k,\ell}(\partial_\ell \partial_k f)(\sum_{i=1}^{N} a_{1i}x_i, \ldots, \sum_{i=1}^{N} a_{di}x_i)$$

$$= \sum_{k=1}^{N} (\partial_k^2 f)(\sum_{i=1}^{N} a_{1i}x_i, \ldots, \sum_{i=1}^{N} a_{di}x_i)$$

$$= (\Delta f)(\sum_{i=1}^{N} a_{1i}x_i, \ldots, \sum_{i=1}^{N} a_{di}x_i)$$

$$= (\Delta f)(Ax)$$

$$= ((\Delta f) \circ A)(x).$$

Hence $\Delta(f \circ A)(x) = ((\Delta f) \circ A)(x)$ and so f is invariant under orthogonal transformations.

$\square$

We will use this result thoughout the rest of the proofs without explicit reference, especially when using translations and rotations.

Let $k > 0$ and H_k be a homothety of origin α and ratio k. That is, $H_k(x) := kx$. For a function u defined on Ω, we define the function $H_k u$ on $H_k(\Omega)$ by $H_k u(x) := u(x/k)$. Since $H_k \circ \Delta = k^2 \Delta \circ H_k$, we have

$$\lambda_n(H_k(\Omega)) = \frac{\lambda_n(\Omega)}{k^2}.$$

should I show more steps?

Theorem 3.2.2. *The minimization problems* $\min\{\lambda_n(\Omega); |\Omega| = c\}$ *as well as* $\min\{|\Omega|^{2/N}\lambda_n(\Omega)\}$ *are equivalent. That is, there exists a bijective correspon-*

dence between the solutions of these two problems.

Proof.

pg 8/9 of Henrot

$\square$

Theorem 3.2.3 (Faber-Krahn)**.** *Let c be a positive number and B the ball with volume c. Then,*

$$\lambda_1(B) = \min\left\{\lambda_1(\Omega), \Omega \text{ open subset of } \mathbb{R}^N, |\Omega| = c\right\}.$$

Proof. This proof is a straightforward application of Schwarz rearrangement (2.4.1) [5]. Let Ω be a bounded open set of measure c and $\Omega^* = B$ be the ball of the same volume. Let u_1 be a en eigenfunction with eigenvalue $\lambda_1(\Omega)$ and u_1^* its Schwarz rearrangement. Using 2.4.1 we have

$$\int_{\Omega^*} u_1^*(x)^2 \, \mathrm{d}x = \int_{\Omega} u_1(x)^2 \, \mathrm{d}x.$$

Further, using 2.4.2 we have

$$\int_{\Omega^*} |\nabla u_1^*(x)|^2 \, \mathrm{d}x \leq \int_{\Omega} |\nabla u_1(x)|^2 \, \mathrm{d}x.$$

Using Rayleigh quotients (3.1.1) we get the following

$$\lambda_1(\Omega^*) \leq \frac{\int_{\Omega} |\nabla u_1^*(x)|^2 \, \mathrm{d}x}{\int_{\Omega} u_1^*(x)^2 \, \mathrm{d}x}.$$

$$\lambda_1(\Omega) = \frac{\int_{\Omega} |\nabla u_1(x)|^2 \, \mathrm{d}x}{\int_{\Omega} u_1(x)^2 \, \mathrm{d}x}.$$

Using the previous two statements yields the desired results. $\square$

3.3 Continuity

For the purpose of the paper, we will only be considering continuity with variable domains.

I need to add essentially all of 2.3.3 in Henrots book pg 28

Look into sources for proofs

3.4 Polygons

Note P_N is the class of plane polygons with at most N edges.

Theorem 3.4.1. *Let $a > 0$ and $N \in \mathbb{N}$ be fixed. Then the problem*

$$\min \{\lambda_1(\Omega), \Omega \in P_N, |\Omega| = a\}$$

has a solution.

47 henrot

Proof. We will use the direct method of calculus of variations. Let Ω_n be a minimizing sequence in P_N for λ_1. Assume

finish proof

$\square$

Theorem 3.4.2. *Let $M \in \mathbb{N}$ and Ω be a polygon with M edges. Then Ω cannot be a (local) minimum for $|\Omega|\lambda_1(\Omega)$ in the class P_{M+1}.*

Theorem 3.4.3 (Pólya)**.** *The equilateral triangle has the least first eigenvalue among all triangles of given area. The square has the least first eigenvalue among all quadrilaterals of given area.*

Proof. This proof uses the same argument as the proof for the Faber-Krahn Theorem, but we will now use Steiner symmetrization in lieu of the Schwarz rearrangement. As the Steiner symmetrization shares the properties 2.4.1 and 2.4.2, we know that any Steiner symmetrization will not increase the first eigenvalue.

We will construct a sequence of Steiner symmetrizations that makes a triangular domain converge to an equilateral triangle. Let a_n, h_n, and A_n be the base, height, and one of the base's incident angles of the triangle T_n that we obtain at step n. Then we have

$$\frac{h_n}{a_{n+1}} = \frac{h_{n+1}}{a_n} = \sin A_n.$$

Denote the ratio $x_n = \frac{h_n}{a_n}$. Then we have

$$x_{n+1} = \frac{\sin^2 A_n}{x_n} = \frac{\sin^2(tan^{-1}(2x_n))}{x_n} = \frac{4x_n}{1 + 4x_n^2}.$$

Thus we have constructed the sequence $x_{n+1} = \frac{4x_n}{1+4x_n^2}$. This will converge to the fixed point of $f(x) = \frac{4x}{1+4x^2}$, which is $\frac{\sqrt{3}}{2}$.

$$\frac{4x}{1 + 4x^2} = x$$
$$x\left(4x^2 - 3\right) = 0$$

and so $x = \frac{\sqrt{3}}{2}$ is the fixed point of f.

One can use elementary geometry to find that for an equilateral triangle with side length a, the height h is $\frac{\sqrt{3}}{2}a$. So $\frac{h}{a} = \frac{\sqrt{3}}{2}$, and thus our sequence converges to the value characteristic of equilateral triangles. Moreover, by Sverak's Theorem, the sequence of triangles γ-converges to the equilateral triangle which we will denote by T_e. Then, for an initial triangle domain T, we have shown

$$\lambda_1(T_e) = \lim \lambda_1(T_n) \le \lambda_1(T).$$

Write proof for quadrilaterals

☐

Explain Issue with proof method for cases N ¿= 5

Theorem 3.4.4. *For $n \ge 3$ the regular polygon with n sides is an extreme point for the first eigenvalue of the Dirichlet Laplace operator among polygons with n sides and a fixed area.*

pg 56 of bogosel paper

Proof. By 3.2.2, our problem is equivalent (up to homothety) to solving the

problem

$$\min_{P \in P_n} \lambda_1(P) + |P|.$$

finish

□

3.5 Examples

Build from Henrot pg 10

We begin with a simple one dimensional case. Let $\Omega = (0, L)$. Solving the differential equation

$$\begin{cases} -u'' = \lambda u & x \in \Omega \\ u(0) = u(L) = 0 \end{cases},$$

we find that the only non-trivial solutions are

$$\lambda_n = \frac{n^2 \pi^2}{L^2}, \ u_n = \sin\left(\frac{n\pi x}{L}\right), \ n \geq 1.$$

$$\Omega = \{(x, y) : 0 < x < a, 0 < y < b, a, b \in \mathbb{R}\}.$$

Finish Rect and add Tri and maybe disk

Chapter 4

A Numerical Approach

4.1 Overview

Copy this section from Final Essay from last year

4.2 Method of Fundamental Solutions

We will use the method of fundamental solutions (MFS) to compute the eigen-values of a given polygon. The following construction is based on a similar method by Alves and Antunes [1]

Our goal is to numerically solve the Helmholtz equation with Dirichlet boundaries

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

We will consider the group of functions which satisfy $-\Delta u = \lambda u$ that are of the form

$$
u = a_1\phi_1^\lambda + \ldots + a_N\phi_N^\lambda,
$$

where $\phi_i^\lambda, i = 1, 2, \ldots, M$ are fundamental radial solutions of $-\Delta u = \lambda u$ with singularities laying outside of Ω. Let (y_i) be the singularities of ϕ_i^λ outside of Ω.

To find the coefficients $a_1, \ldots, a_N$ we impose the Dirichlet boundary con-dition on a discretization of $\partial\Omega$. Let (x_i) be a discretization of $\partial\Omega$, and let $x_{N+1} \in \Omega$. This leads to a system of equations

$$
\begin{cases}
u(x_i) = 0 & \text{if } 1 \leq i \leq N \\
u(x_i) = 1 & \text{if } i = N + 1
\end{cases}
$$

Note that the equation when $i = N + 1$ is used to guarantee that $u(x) \not\equiv 0$ [2].

Obviously we are interested when the system has non-trivial solutions. This occurs when the matrix $A_\lambda = (\phi_i^\lambda(x_j))_{i,j=1}^N$ is singular. As this shows the existence of an eigenfunction, we can find eigenvalues using the determinate of the matrix A_λ. Specifically, we can find the eigenvalues of Ω on some interval I by locating the values $\lambda \in I$ where $\det A_\lambda = 0$. Once we have found an

eigenvalue, we can solve the system to find a corresponding eigenfunction.

To apply MFS to our specific problem, we need to find suitable radial functions as well as (x_i) and (y_i).

First, we will find suitable radial functions. Let $\phi := x(r)$ be a radial function in polar coordinates. Then Helmholtz's equation becomes

$$-x'' - \frac{1}{r}x' = \lambda x$$
$$r^2 x'' + rx' + r^2 \lambda x = 0.$$

Substituting in $s = \sqrt{\lambda} r$ we have

$$s^2 y'' + sy' + s^2 y = 0,$$

where $y(s) = x(r)$. Note this is a specific case of Bessel's differential equation. Thus, our radial fundamental solutions can be Bessel functions of order 0. We choose to use the Hankel function of the first kind with order 0 as it is the most efficient computationally.

Definition 4.2.1. *We define the Bessel function of the first kind with order 0 in the following way*

$$J_0(x) = \frac{1}{\pi} \int_0^\pi \cos(x \sin \tau) \, d\tau.$$

We define the Bessel function of the second kind with order 0 in the following way

$$Y_0(x) = \frac{4}{\pi^2} \int_0^{\frac{1}{2}\pi} \cos(x \cos \tau) \left(e + \ln \left(2x \sin^2 \tau \right) \right) \, d\tau.$$

Finally, we define the Hankel function (of the first kind) with order 0 as

$$H_0(x) = J_0(x) + iY_0(x).$$

Thus our fundamental solutions will be of the form $\phi_i^\lambda = H_0(\sqrt{\lambda}|x - y_i|)$.

Should I go into how to choose $(x_i), (y_i)$ or refer to various papers and summarize?

4.3 Optimization

In the previous section we outlined the method of fundamental solutions, a method to calculate the eigenvalues of our equation. In this section we will outline a method to calculate the derivative of the eigenvalue and use gradient descent to find extremum.

From 3.1.1, the derivative of an eigenvalue is given by

$$\lambda'_k(0) := -\int_\Omega \left(\frac{\partial u}{\partial n}\right)^2 V.n \, \mathrm{d}\sigma.$$

Maybe derive formula? IDK how involved it is.

As we are taking the derivative with respect to the domain, we will begin by defining vector fields that allow us to write the derivative with respect to geometric parameters. We will find particular vector fields V which allow us to compute the derivative with respect to the coordinates of the vertices. Fix a vertex and label it v_0. Next, label the remaining vertices $v_1, v_2, \ldots, v_{N-1}$ going around the polygon counterclockwise. That is, for a vertex v_i the adjacent vertices should be v_{i-1}, v_{i+1} modulo N. Finally, take (x_i, y_i) to be the coordinates of v_i.

To find the derivative of λ_1 with respect to x_i we make a perturbation of v_i with $(1, 0)$. This induces a perturbation of the adjacent edges of the boundary, which we will denote as $E_{i-1,i}$ and $E_{i,i+1}$. For our particular case V will have the following form on the boundary

$$\begin{cases} L_{i-1,i}(x, y) & (x, y) \in E_{i-1,i} \\ L_{i,i+1}(x, y) & (x, y) \in E_{i,i+1} \\ 0 & \text{otherwise} \end{cases}$$

where $L_{j,k} : E_{j,k} \to [0, 1]$ is the following affine function

$$
L_{j,k}(x, y) =
\begin{cases}
(x_k - x_j)^{-1}(x - x_i) & \text{if } x_i \neq x_{i+1} \\
\\
0 & \text{otherwise}
\end{cases}
$$

Denote the outer normal of the edge $E_{i,i+1}$ by $n_{j,j+1} = (n_{j,j+1}^1, n_{j,j+1}^2)$. Then we can rewrite the derivative of the fundamental eigenvalue as

$$
\frac{d\lambda_1}{dx_i} = - \int_{E_{i-1,i}} L_{i-1,i} \left(\frac{\partial u}{\partial n} \right)^2 n_{i-1,i}^1 \, d\sigma - \int_{E_{i,i+1}} L_{i+1,i} \left(\frac{\partial u}{\partial n} \right)^2 n_{i,i+1}^1 \, d\sigma.
$$

Likewise, we can find the derivative with respect to the y value as

$$
\frac{d\lambda_1}{dx_{2i}} = - \int_{E_{i-1,i}} L_{i-1,i} \left(\frac{\partial u}{\partial n} \right)^2 n_{i-1,i}^2 \, d\sigma - \int_{E_{i,i+1}} L_{i+1,i} \left(\frac{\partial u}{\partial n} \right)^2 n_{i,i+1}^2 \, d\sigma.
$$

Following these computations, we have all of the pieces needed to optimize the fundamental eigenvalue using gradient descent.

Give Numerical Results

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