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# Internal and External Harmonic Functions in FlatRing Coordinates 

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# INTERNAL AND EXTERNAL HARMONIC FUNCTIONS IN FLAT-RING COORDINATES 

by

Lijuan Bi

A Dissertation Submitted in Partial Fulfillment of the Requirements for the Degree of

Doctor of Philosophy in
Mathematics
at
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# ABSTRACT <br> INTERNAL AND EXTERNAL HARMONIC FUNCTIONS IN FLAT-RING COORDINATES 

by

## Lijuan Bi

The University of Wisconsin-Milwaukee, 2018
Under the Supervision of Professor Hans Volkmer

The goal of this dissertation is to derive expansions for a fundamental solution of Laplace's equation in flat-ring coordinates in three-dimensional Euclidean space. These expansions are in terms of harmonic functions in the interior and the exterior of two different types of regions, "flat rings" and "peanuts" according to their shapes. We solve Laplace's equation in the interior and the exterior of these regions using the method of separation of variables. The internal and external "flat-ring" and "peanut" harmonic functions are expressed in terms of Lamé functions.
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## Chapter 1

## Introduction

It has been known for a long time that Laplace's equation in cartesian coordinates on $\mathbb{R}^{3}$

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial^{2} u}{\partial y^{2}}+\frac{\partial^{2} u}{\partial z^{2}}=0 \tag{1.0.1}
\end{equation*}
$$

can be solved in a quadric or a cyclidic coordinate system on $\mathbb{R}^{3}$ using separation of variables methods. A cyclidic coordinate system in $\mathbb{R}^{3}$ has three coordinates, each coordinate is obtained by a transformation from cartesian coordinates in $\mathbb{R}^{3}$, and the coordinate surfaces are cyclides. A cyclide in general is a fourth-order surface, which is a natural extension of a quadric surface. A quadric surface in general is a second-order surface.

Bôcher listed quadric and cyclidic coordinate systems where Laplace's equation can be solved using separation of variables in his book "Ueber die Reihenentwicklungen der Potentialtheorie" [9] published in 1894. We can also find these quadric and cyclidic coordinates systems in Miller's Book ([4], Table 14 on page 164, Table 17 on page 210). The quadric coordinates are cartesian, cylindrical, parabolic cylindrical, elliptic cylindrical, spherical, prolate spheroidal, oblate spheroidal, parabolic, paraboloidal, conical, and ellipsoidal coordinates. The cyclidic coordinates are flat-ring cyclide, flat-disk cyclide, bi-cyclide, cap-cyclide, and 3 -cyclide coordinates.

Using separation of variables methods, a solution of (1.0.1) can be written in the form

$$
u(x, y, z)=R(\alpha, \beta, \lambda) v_{1}(\alpha) v_{2}(\beta) v_{3}(\lambda)
$$

where $\alpha, \beta, \lambda$ are quadric or cyclidic coordinates, $R(\alpha, \beta, \lambda)$ is a fixed and known function, and $v_{1}(\alpha), v_{2}(\beta), v_{3}(\lambda)$ are solutions of a second order linear homogeneous ordinary differential equation, respectively. We have $R(\alpha, \beta, \lambda)=1$ if $\alpha, \beta, \lambda$ are quadric coordinates, and $R(\alpha, \beta, \lambda) \neq 1$, if $\alpha, \beta, \lambda$ are cyclidic coordinates. For example, we solve the (1.0.1) in a well-known quadric coordinate system, the spherical coordinate system,

$$
\begin{equation*}
x=r \cos \phi \sin \theta, y=r \sin \phi \sin \theta, z=r \cos \theta \tag{1.0.2}
\end{equation*}
$$

$$
r>0,-\pi<\phi<\pi, 0<\theta<\pi
$$

The equation (1.0.1) has solutions of the form

$$
u(x, y, z)=u_{1}(r) u_{2}(\theta) u_{3}(\phi)
$$

The $u_{1}, u_{2}, u_{3}$ have to satisfy the following equations, respectively:

$$
\begin{gather*}
r^{2} u_{1}^{\prime \prime}+2 r u_{1}^{\prime}-n(n+1) u_{1}=0,  \tag{1.0.3}\\
u_{2}^{\prime \prime}+\cot \theta u_{2}^{\prime}+\left[n(n+1)-\frac{m^{2}}{\sin ^{2} \theta}\right] u_{2}=0,  \tag{1.0.4}\\
u_{3}^{\prime \prime}+m^{2} u_{3}=0 \tag{1.0.5}
\end{gather*}
$$

where $m, n$ are separation parameters. A general solution of (1.0.3) is

$$
u_{1}(r)=c_{1} r^{n}+c_{2} r^{-n-1}
$$

where $c_{1}$ and $c_{2}$ are constants. A general solution of (1.0.4) is

$$
u_{2}(\theta)=c_{1} P_{n}^{m}(\cos \theta)+Q_{n}^{m}(\cos \theta)
$$

where $P_{n}^{m}$ and $Q_{n}^{m}$ are the associated Legendre functions of the first kind and the second kind, respectively. We consider the simplest cyclidic coordinate system, the toroidal coordinate system,

$$
x=\frac{\sinh \sigma \cos \phi}{\cosh \sigma-\cos \psi}, y=\frac{\sinh \sigma \sin \phi}{\cosh \sigma-\cos \psi}, z=\frac{\sin \psi}{\cosh \sigma-\cos \psi} .
$$

The coordinate surfaces $\sigma=$ constant are tori, i.e.

$$
\left(1+x^{2}+y^{2}+z^{2}\right)^{2}=4\left(x^{2}+y^{2}\right) \operatorname{coth}^{2} \sigma
$$

The coordinate surfaces $\psi=$ constant are spherical bowls, i.e.

$$
(z-\cot \psi)^{2}+x^{2}+y^{2}=\frac{1}{\sin ^{2} \psi}
$$

The coordinate surfaces $\phi=$ constant are planes, i.e.

$$
x \sin \phi=y \cos \phi .
$$

We assume

$$
u(x, y, z)=\sqrt{\cosh \sigma-\cos \psi} u_{1}(\sigma) u_{2}(\phi) u_{3}(\psi)
$$

We find $u_{1}, u_{2}, u_{3}$ satisfies each of the following equations, respectively:

$$
\begin{gather*}
\frac{1}{\sinh \sigma}\left(\sinh \sigma u_{1}^{\prime}\right)^{\prime}-\left(n^{2}-\frac{1}{4}+\frac{m^{2}}{\sinh ^{2} \sigma}\right) u_{1}=0  \tag{1.0.6}\\
u_{2}^{\prime \prime}+n^{2} u_{2}=0  \tag{1.0.7}\\
u_{3}^{\prime \prime}+m^{2} u_{3}=0 \tag{1.0.8}
\end{gather*}
$$

where $m, n$ are separation parameters. A general solution of (1.0.6) is

$$
u_{1}(\sigma)=c_{1} P_{n-\frac{1}{2}}^{m}(\cosh \sigma)+c_{2} Q_{n-\frac{1}{2}}^{m}(\cosh \sigma),
$$

where $c_{1}$ and $c_{2}$ are constants, and $P_{n}^{m}$ and $Q_{n}^{m}$ are the associated Legendre functions of the first kind and the second kind, respectively. Special functions we obtain from solving Laplace's equation in a quadric coordinate system via separation of variables are functions like Bessel functions, Legendre functions, Lamé polynomials, etc, which are some of the most classical special functions. For a cyclidic coordinate, special functions are functions like Legendre functions, Lamé periodic functions, Lamé-Wangerin functions, solutions of an ordinary equation with five elementary singularities (see [6]), etc.

The Helmholtz equation on $\mathbb{R}^{3}$,

$$
\frac{\partial^{2} u}{\partial x^{2}}+\frac{\partial^{2} u}{\partial y^{2}}+\frac{\partial^{2} u}{\partial z^{2}}+k^{2} u=0
$$

where $k^{2}$ is a real number, can be solved via separation of variables in the same quadric coordinate systems that Laplace's equation can be solved using separation of variables. However, the Helmholtz equation cannot be solved in the cyclidic coordinate systems listed above using separation of variables.

It is also long known that a fundamental solution of Laplace's equation in $\mathbb{R}^{3}$, i.e.

$$
\left.\frac{1}{\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|}=\left(x-x^{\prime}\right)^{2}+\left(y-y^{\prime}\right)^{2}+\left(z-z^{\prime}\right)^{2}\right)^{-\frac{1}{2}}
$$

where $\mathbf{r}=(x, y, z), \mathbf{r}^{\prime}=\left(x^{\prime}, y^{\prime}, z^{\prime}\right)$, and $\mathbf{r} \neq \mathbf{r}^{\prime}$, can be expressed using solutions from separation of variables (see [10], [11], [4], and [12]). For example, we consider spherical coordinates. The internal spherical harmonics are

$$
G_{n}^{m}(x, y, z)=r^{n} P_{n}^{m}(\cos \theta) e^{i m \phi},-n \leq m \leq n
$$

where $P_{n}^{m}$ is an associated Legendre function of the first kind (Ferrer's function). $G_{n}^{m}$ is a harmonic function in $\mathbb{R}^{3}$. The external spherical harmonics are

$$
H_{n}^{m}(x, y, z)=r^{-n-1} P_{n}^{m}(\cos \theta) e^{i m \phi},-n \leq m \leq n .
$$

$H_{n}^{m}$ is a harmonic function in $\mathbb{R}^{3}$ except the origin. An expansion of a fundamental solution is

$$
\begin{equation*}
\frac{1}{\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|}=\sum_{n=0}^{\infty} \sum_{m=-n}^{n} \frac{(n-m)!}{(n+m)!}\left(G_{n}^{m}(\mathbf{r}) \overline{H_{n}^{m}\left(\mathbf{r}^{\prime}\right)}\right) \tag{1.0.9}
\end{equation*}
$$

where $\|\mathbf{r}\|<\left\|\mathbf{r}^{\prime}\right\|$. As another example, we can consider toroidal coordinates. The internal toroidal harmonics are

$$
G_{n}^{m}(x, y, z)=\sqrt{\cosh \sigma-\cos \psi} Q_{n-\frac{1}{2}}^{m}(\cosh \sigma) e^{i n \psi} e^{i m \phi}
$$

where $m, n \in \mathbb{Z}$ and $Q_{n-\frac{1}{2}}^{m}$ is the associated Legendre function of the second kind. $G_{n}^{m}$ is harmonic in $\mathbb{R}^{3}$ except for the $z$-axis. The external toroidal harmonics are

$$
H_{n}^{m}(x, y, z)=\sqrt{\cosh \sigma-\cos \psi} P_{n-\frac{1}{2}}^{m}(\cosh \sigma) e^{i n \psi} e^{i m \phi}
$$

where $m, n \in \mathbb{Z}$ and $P_{n-\frac{1}{2}}^{m}$ is the associated Legendre function of the first kind. $H_{n}^{m}$ is harmonic in $\mathbb{R}^{3}$ except for the unit circle $z=0, x^{2}+y^{2}=1$.

In this dissertation we consider flat-ring coordinates which is listed as number 15 in Miller's list [9, page 210] (see [4]). The coordinates $\rho, \mu, \phi$ in algebraic form compose an orthogonal coordinate system in $\mathbb{R}^{3}$ with rotational symmetry. They vary according to $1<\rho<a<\mu<\infty$, where $a$ is a given number. For our convenience, we use flat-ring coordinates $\alpha, \beta, \phi$ in transcendental form, which is a transformation of the algebraic form. The surface $\alpha=$ constant looks similar to a flat-ring. The surface $\beta=$ constant looks similar to a peanut. We derive expansions analogous to (1.0.9) for flat-ring coordinate system in terms of flat-ring harmonics and peanut harmonics, respectively. As far as we know, the expansions are given for the first time.

The outline of this dissertation is as follows. In chapter 2 we discuss flat-ring coordinates. In chapter 3 we solve Laplace's equation in flat-ring coordinates using the method of separation of variables. The solutions are in terms of Lamé functions. In chapter 4 we collect some facts about Lamé functions that will be used in the later chapters. In chapter 5 we derive internal and external flat-ring harmonics, solve the Dirichlet problem and provide expansions for a fundamental solution of Laplace's equation in terms of internal and external flat-ring harmonics. In chapter 6 we derive internal and external peanut harmonics, solve the Dirichlet problem and provide an expansion for a fundamental solution of Laplace's equation in terms of internal and external peanut harmonics.

This introduction is based on Cohl's manuscript (see [2]), the introduction of [6], and the introduction of [7].

## Chapter 2

## Flat-Ring Coordinates

### 2.1 Flat-Ring Coordinates in Algebraic Form

Flat-ring coordinates $\rho, \mu, \phi$ form an orthogonal coordinate system in $\mathbb{R}^{3}$ with rotational symmetry. According to ([4], system 15, page 210), it is given by

$$
\begin{equation*}
x=\frac{\cos \phi}{r}, \quad y=\frac{\sin \phi}{r}, \quad z=\frac{1}{r}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}}, \tag{2.1.1}
\end{equation*}
$$

where

$$
\begin{equation*}
r=\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}+\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}} \tag{2.1.2}
\end{equation*}
$$

The number $a>1$ is given. The coordinates $\rho, \mu, \phi$ vary according to

$$
1<\rho<a<\mu<\infty .
$$

The square roots appearing (2.1.1) and (2.1.2) are understood to be positive.
It is easier to visualize in two dimensions. So we set $\phi=0$ to remove rotational symmetry and look at the two-dimensional coordinate system as a slice. Now

$$
\begin{equation*}
x=\frac{1}{r}, \quad z=\frac{1}{r}\left(\frac{((\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}} \tag{2.1.3}
\end{equation*}
$$

where

$$
\begin{equation*}
r=\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}+\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}} \tag{2.1.4}
\end{equation*}
$$

We will show the mapping $(x, z)=f(\rho, \mu)$ according to (2.1.3) and (2.1.4) is bijective.

Proposition 1. If $A=(1, a) \times(a, \infty), Q_{0}=\left\{(x, z): x>0, z>0, x^{2}+z^{2}<1\right\}$, and $(x, z)=f(\rho, \mu)$ according to (2.1.3) and (2.1.4), we obtain that $f$ maps $A$ to $Q_{0}$ and

$$
\begin{equation*}
\frac{\left(x^{2}+z^{2}+1\right)^{2}}{s}-\frac{\left(x^{2}+z^{2}-1\right)^{2}}{s-1}-\frac{4 z^{2}}{s-a}=0, \tag{2.1.5}
\end{equation*}
$$

where $s=\mu$ or $\rho$.
Proof. First we show that $f$ maps $A$ to $Q_{0}$. Let $(\rho, \mu) \in A$ and $(x, z)=f(\rho, \mu)$ according to (2.1.3) and (2.1.4). Since the square roots in (2.1.3) and (2.1.4) are understood to be positive, we have $x>0$ and $z>0$. We wish to show $x^{2}+z^{2}<1$. According to (2.1.4), we have

$$
r=\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}+\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}}
$$

We find

$$
\begin{equation*}
\frac{1}{r}=\frac{1}{\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}+\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}}}=\frac{\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}-\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}}}{\left(\frac{\mu a-\mu \rho+\rho a-a}{a(a-1)}\right)} . \tag{2.1.6}
\end{equation*}
$$

We also calculate

$$
\begin{align*}
x^{2}+z^{2}-1 & =\frac{1}{r^{2}}+\frac{1}{r^{2}}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)-1  \tag{2.1.7}\\
& =\frac{1}{r^{2}}\left(\frac{\mu a-\mu \rho+\rho a-a}{a(a-1)}\right)-1 . \tag{2.1.8}
\end{align*}
$$

Using the results from (2.1.6) and (2.1.8), we obtain

$$
\begin{aligned}
\frac{x^{2}+z^{2}-1}{\frac{1}{r}} & =\frac{\frac{1}{r^{2}}\left(\frac{\mu a-\mu \rho+\rho a-a}{a(a-1)}\right)-1}{\frac{1}{r}} \\
& =\frac{1}{r}\left(\frac{\mu a-\mu \rho+\rho a-a}{a(a-1)}\right)-r \\
& =-2\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}}<0 .
\end{aligned}
$$

It is clear that

$$
\begin{equation*}
x^{2}+z^{2}-1<0, \tag{2.1.9}
\end{equation*}
$$

because

$$
\frac{1}{r}>0
$$

Hence $f$ maps $A$ to $Q_{0}$. Second, we wish to show equation (2.1.5) is true.
We consider

$$
\begin{aligned}
& r^{2}\left[\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\mu}-\frac{\left(x^{2}+z^{2}-1\right)^{2}}{\mu-1}-\frac{4 z^{2}}{\mu-a}\right] \\
= & \frac{\left(\frac{x^{2}+z^{2}+1}{\frac{1}{r}}\right)^{2}}{\mu}-\frac{\left(\frac{x^{2}+z^{2}-1}{\frac{1}{r}}\right)^{2}}{\mu-1}-4 \frac{a-\rho}{a(a-1)} \\
= & \frac{\left(2\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}\right)^{2}}{\mu}-\frac{\left(2\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}}\right)^{2}}{\mu-1}-4 \frac{a-\rho}{a(a-1)} \\
= & 4\left[\frac{\rho}{a}-\frac{\rho-1}{a-1}-\frac{a-\rho}{a(a-1)}\right]=0 .
\end{aligned}
$$

So it is clear that

$$
\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\mu}-\frac{\left(x^{2}+z^{2}-1\right)^{2}}{\mu-1}-\frac{4 z^{2}}{\mu-a}=0 .
$$

With a similar calculation, we have

$$
\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\rho}-\frac{\left(x^{2}+z^{2}-1\right)^{2}}{\rho-1}-\frac{4 z^{2}}{\rho-a}=0 .
$$

Proposition 2. If $A=(1, a) \times(a, \infty), Q_{0}=\left\{(x, z): x, z>0, x^{2}+z^{2}<1\right\}$, and $(x, z)=$ $f(\rho, \mu)$ according to (2.1.3) and (2.1.4), then $f: A \rightarrow Q_{0}$ is bijective.

Proof. Let $(x, z) \in Q_{0}$. Assume

$$
\begin{equation*}
F(s)=s(s-1)(s-a)\left[\frac{\left(x^{2}+z^{2}+1\right)^{2}}{s}-\frac{\left(x^{2}+z^{2}-1\right)^{2}}{s-1}-\frac{4 z^{2}}{s-a}\right], \tag{2.1.10}
\end{equation*}
$$

The coefficient of $s^{2}$ is

$$
\begin{equation*}
\left(x^{2}+z^{2}+1\right)^{2}-\left(x^{2}+z^{2}-1\right)^{2}-4 z^{2}=4 x^{2}>0 \tag{2.1.11}
\end{equation*}
$$

Moreover,

$$
\begin{aligned}
& \lim _{s \rightarrow 0} F(s)=a\left(x^{2}+z^{2}+1\right)^{2}>0, \\
& \lim _{s \rightarrow 1} F(s)=(a-1)\left(x^{2}+z^{2}-1\right)^{2}>0, \\
& \lim _{s \rightarrow a} F(s)=-a(a-1) 4 z^{2}<0
\end{aligned}
$$

If $F(s)=0$, we can find two real solutions, say $\rho$ and $\mu$, such that $1<\rho<a$ and $\mu>a$ according to (2.1.11) and the inequalities above. Then we obtain

$$
\begin{equation*}
F(s)=4 x^{2}(s-\rho)(s-\mu) \tag{2.1.12}
\end{equation*}
$$

where $s$ is a real number. According to (2.1.12), we get

$$
\begin{aligned}
& \lim _{s \rightarrow 0} F(s)=4 x^{2} \rho \mu \\
& \lim _{s \rightarrow 1} F(s)=4 x^{2}(1-\rho)(1-\mu), \\
& \lim _{s \rightarrow a} F(s)=4 x^{2}(a-\rho)(a-\mu)
\end{aligned}
$$

We match the above corresponding limits respectively and we get

$$
\begin{gather*}
a\left(x^{2}+z^{2}+1\right)^{2}=4 x^{2} \rho \mu  \tag{2.1.13}\\
(a-1)\left(x^{2}+z^{2}-1\right)^{2}=4 x^{2}(1-\rho)(1-\mu)  \tag{2.1.14}\\
-a(a-1) 4 z^{2}=4 x^{2}(a-\rho)(a-\mu) \tag{2.1.15}
\end{gather*}
$$

From equation (2.1.13), we have

$$
\begin{equation*}
\frac{\rho \mu}{a}=\frac{\left(x^{2}+z^{2}+1\right)^{2}}{4 x^{2}} \tag{2.1.16}
\end{equation*}
$$

From equation (2.1.14), we obtain

$$
\begin{equation*}
\frac{(\rho-1)(\mu-1)}{a-1}=\frac{\left(x^{2}+z^{2}-1\right)^{2}}{4 x^{2}} \tag{2.1.17}
\end{equation*}
$$

We define

$$
\begin{equation*}
r=\left(\frac{\rho \mu}{a}\right)^{\frac{1}{2}}+\left(\frac{(\mu-1)(\rho-1)}{\mu-1}\right)^{\frac{1}{2}} \tag{2.1.18}
\end{equation*}
$$

then we substitute equation (2.1.16) and equation (2.1.17) in equation (2.1.18) to get

$$
r=\frac{x^{2}+z^{2}+1}{2 x}-\frac{x^{2}+z^{2}-1}{2 x}=\frac{1}{x}
$$

Hence,

$$
\begin{equation*}
x=\frac{1}{r} \tag{2.1.19}
\end{equation*}
$$

Then plug equation (2.1.19) in equation (2.1.15), we have

$$
z=\frac{1}{r}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}} .
$$

Therefore, $f$ is surjective. Suppose $f(\rho, \mu)=f(\bar{\rho}, \bar{\mu})=(x, z)$. By Proposition 1, $\rho, \mu, \bar{\rho}, \bar{\mu}$ are all roots of $F(s)=0$ with $F$ defined in (2.1.12). Since $(\rho, \mu) \in A,(\bar{\rho}, \bar{\mu}) \in A$, this implies $(\rho, \mu)=(\bar{\rho}, \bar{\mu})$. Thus, $f$ is injective. Hence, $f$ is bijective.

If we change the signs of the square roots in (2.1.3) and (2.1.4) we obtain coordinates for other regions in $\mathbb{R}^{2}$ according to Table 1.

|  | $\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}$ | $\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}}$ | $\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}}$ |
| :--- | :---: | :---: | :---: |
| $x>0, z>0, x^{2}+z^{2}<1$ | + | + | + |
| $x>0, z>0, x^{2}+z^{2}>1$ | + | - | + |
| $x>0, z<0, x^{2}+z^{2}>1$ | + | - | - |
| $x>0, z<0, x^{2}+z^{2}<1$ | + | + | - |

Table 2.1: Signs of roots in regions of $(x, z)$-plane

Some coordinate lines are shown in Figure 2.1. Note that the boundary of the region is given by the quarter circle $\rho=1$, the vertical segment $\mu=\infty$, and two horizontal segments $[0, b],[b, 1]$ represented by $\rho=a$ and $\mu=a$ respectively, where $b=\sqrt{a}-\sqrt{a-1}=$ $(\sqrt{a}+\sqrt{a-1})^{-1}$.


Figure 2.1: Coordinate lines of system (2.1.3) for $a=2$

We calculate the metric coefficients in this chapter. We will use them in the later chapters. Proposition 3. The metric coefficients $h_{\rho}$ and $h_{\mu}$ for (2.1.3) and (2.1.4) are given by

$$
\begin{aligned}
& h_{\rho}^{2}:=\left(\frac{\partial x}{\partial \rho}\right)^{2}+\left(\frac{\partial z}{\partial \rho}\right)^{2} \\
& =\frac{\mu-\rho}{4(a-\rho) \rho(\rho-1)} \frac{1}{r^{2}} \\
& =\frac{1}{16}\left(-\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\rho^{2}}+\frac{\left(x^{2}+z^{2}-1\right)^{2}}{(\rho-1)^{2}}+\frac{4 z^{2}}{(\rho-a)^{2}}\right) \\
& h_{\mu}^{2}:=\left(\frac{\partial x}{\partial \mu}\right)^{2}+\left(\frac{\partial z}{\partial \mu}\right)^{2} \\
& =\frac{\mu-\rho}{4(\mu-a) \mu(\mu-1)} \frac{1}{r^{2}} \\
& =\frac{1}{16}\left(-\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\mu^{2}}+\frac{\left(x^{2}+z^{2}-1\right)^{2}}{(\mu-1)^{2}}+\frac{4 z^{2}}{(\mu-a)^{2}}\right)
\end{aligned}
$$

Proof. Let

$$
J=\left(\begin{array}{ll}
\frac{\partial x}{\partial \rho} & \frac{\partial x}{\partial \mu} \\
\frac{\partial z}{\partial \rho} & \frac{\partial z}{\partial \mu}
\end{array}\right),
$$

where $(x, z)=f(\rho, \mu)$ according to (2.1.3) and (2.1.4). First, we wish to show that

$$
J^{T} J=\left(\begin{array}{cc}
h_{\rho}^{2} & 0 \\
0 & h_{\mu}^{2}
\end{array}\right)
$$

where

$$
\begin{aligned}
h_{\rho}^{2} & =\frac{\mu-\rho}{4(a-\rho) \rho(\rho-1)} \frac{1}{r^{2}} \\
h_{\mu}^{2} & =\frac{\mu-\rho}{4(\mu-a) \mu(\mu-1)} \frac{1}{r^{2}} .
\end{aligned}
$$

We calculate

$$
\begin{equation*}
\frac{\partial r}{\partial \rho}=\frac{1}{2}\left(\frac{\mu}{a \rho}\right)^{\frac{1}{2}}+\frac{1}{2}\left(\frac{\mu-1}{(a-1)(\rho-1)}\right)^{\frac{1}{2}} \tag{2.1.20}
\end{equation*}
$$

which can be used for the rest of this proposition. Let

$$
\begin{equation*}
L_{1}=\frac{1}{r} \frac{\partial r}{\partial \rho}\left(1+\frac{(\mu-a)(a-\rho)}{a(a-1)}\right) . \tag{2.1.21}
\end{equation*}
$$

We have

$$
\begin{aligned}
h_{\rho}^{2} & =\left(-\frac{1}{r^{2}} \frac{\partial r}{\partial \rho}\right)^{2} \\
& +\left(-\frac{1}{r^{2}}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}}\left(\frac{\partial r}{\partial \rho}\right)-\frac{1}{2} \frac{1}{r}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{-\frac{1}{2}}\left(\frac{\mu-a}{a(a-1)}\right)\right)^{2} \\
& =\frac{1}{r^{2}}\left(\frac{1}{r^{2}}\left(\frac{\partial r}{\partial \rho}\right)^{2}\left(1+\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)+\frac{1}{r}\left(\frac{\mu-a}{a(a-1)}\right)\left(\frac{\partial r}{\partial \rho}\right)+\frac{1}{4}\left(\frac{a(a-1)}{(\mu-a)(a-\rho)}\right)\left(\frac{\mu-a}{a(a-1)}\right)^{2}\right) \\
& =\frac{1}{r^{2}}\left(1+\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{-1}\left(\left(L_{1}+\frac{1}{2}\left(\frac{\mu-a}{a(a-1)}\right)^{2}+\frac{1}{4}\left(\frac{\mu-a}{a(a-1)(a-\rho)}\right)\right)\right. \\
& =\frac{\mu-\rho}{4(a-\rho) \rho(\rho-1)} \frac{1}{r^{2}} .
\end{aligned}
$$

Similarly, we have

$$
h_{\mu}^{2}=\frac{\mu-\rho}{4(\mu-a) \mu(\mu-1)} \frac{1}{r^{2}} .
$$

Substitute (2.1.6) and (2.1.20) in (2.1.21) and simplify it then we have

$$
L_{1}=-\frac{1}{2} \frac{\mu-a}{a(a-1)}+\frac{1}{2}\left(\frac{\mu(\mu-1)}{a(a-1)}\right)^{\frac{1}{2}}\left(\left(\frac{\rho}{\rho-1}\right)^{\frac{1}{2}}-\left(\frac{\rho-1}{\rho}\right)^{\frac{1}{2}}\right)
$$

Similarly, we have

$$
L_{2}:=\frac{1}{r} \frac{\partial r}{\partial \mu}\left(1+\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)=\frac{1}{2} \frac{a-\rho}{a(a-1)}+\frac{1}{2}\left(\frac{\rho(\rho-1)}{a(a-1)}\right)^{\frac{1}{2}}\left(\left(\frac{\mu}{\mu-1}\right)^{\frac{1}{2}}-\left(\frac{\mu-1}{\mu}\right)^{\frac{1}{2}}\right) .
$$

$L_{1}$ and $L_{2}$ will be used in the following calculation. The off-diagonal element in $J^{T} J$ multiplied by $r^{2}$ defined in (2.1.4) is

$$
\begin{aligned}
& r^{2}\left(\frac{\partial x}{\partial \rho} \frac{\partial x}{\partial \mu}+\frac{\partial z}{\partial \rho} \frac{\partial z}{\partial \mu}\right)=\frac{\frac{\partial r}{\partial \rho}}{r} \frac{\frac{\partial r}{r}}{\frac{\partial r}{r}}+ \\
& \left(\frac{\frac{\partial r}{\partial \rho}}{-r}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}}-\frac{1}{2}\left(\frac{\mu-a}{(a-\rho) a(a-1)}\right)^{\frac{1}{2}}\right)\left(\frac{\frac{\partial r}{\partial \mu}}{-r}\left(\frac{(\mu-a)(a-\rho)}{a(a-1)}\right)^{\frac{1}{2}}+\frac{1}{2}\left(\frac{a-\rho}{(\mu-a) a(a-1)}\right)^{\frac{1}{2}}\right) \\
& =0
\end{aligned}
$$

We can compute $J^{T} J$ in a slightly different way as follows. Let

$$
K=\left(\begin{array}{ll}
\frac{\partial \rho}{\partial x} & \frac{\partial \rho}{\partial z} \\
\frac{\partial \mu}{\partial x} & \frac{\partial \mu}{\partial z}
\end{array}\right)=J^{-1} .
$$

Then

$$
J^{T} J=\left(K K^{T}\right)^{-1} .
$$

Let

$$
\begin{aligned}
& F_{s}:=\frac{\left(x^{2}+z^{2}+1\right)^{2}}{s}-\frac{\left(x^{2}+z^{2}-1\right)^{2}}{s-1}-\frac{4 z^{2}}{s-a}, \\
& P_{s}:=\frac{\partial F_{s}}{\partial s}=-\frac{\left(x^{2}+z^{2}+1\right)^{2}}{s^{2}}+\frac{\left(x^{2}+z^{2}-1\right)^{2}}{(s-1)^{2}}+\frac{4 z^{2}}{(s-a)^{2}} \\
& Q_{s}:=\frac{x^{2}+z^{2}+1}{s}-\frac{x^{2}+z^{2}-1}{s-1} \\
& \tilde{Q}_{s}:=\frac{x^{2}+z^{2}+1}{s}-\frac{x^{2}+z^{2}-1}{s-1}, \text { where } s=\mu \text { or } \rho .
\end{aligned}
$$

In these definitions $x, z, \rho$, and $\mu$ are independent variables.
Let $(x, z)=f(\rho, \mu)$. We have

$$
\begin{aligned}
& \frac{\partial F_{s}}{\partial s} \frac{\partial s}{\partial x}=P_{\mu} \frac{\partial s}{\partial x}=\frac{4 x\left(x^{2}+z^{2}+1\right)}{s}-\frac{4 x\left(x^{2}+z^{2}-1\right)}{s-1}=4 x Q_{s} \\
& \frac{\partial F_{s}}{\partial s} \frac{\partial s}{\partial z}=P_{s} \frac{\partial s}{\partial z}=\frac{4 x\left(x^{2}+z^{2}+1\right)}{s}-\frac{4 x\left(x^{2}+z^{2}-1\right)}{s-1}-4 z \frac{2}{s-a}=4 x\left(Q_{s}-\frac{2}{s-a}\right) .
\end{aligned}
$$

We can obtain

$$
\begin{aligned}
P_{s} \frac{\partial s}{\partial x} & =4 x Q_{s} \\
P_{s} \frac{\partial s}{\partial z} & =4 x\left(Q_{s}-\frac{2}{s-a}\right)
\end{aligned}
$$

We have

$$
\begin{aligned}
& \frac{P_{s} P_{t}}{16} k_{i j} \\
= & \frac{P_{s} P_{t}}{16}\left(\frac{\partial s}{\partial x} \frac{\partial t}{\partial x}+\frac{\partial s}{\partial z} \frac{\partial t}{\partial z}\right) \\
= & x^{2} Q_{s} Q_{t}+z^{2}\left(Q_{s}-\frac{2}{s-a}\right)\left(Q_{t}-\frac{2}{t-a}\right) \\
= & \frac{1}{2} Q_{s} F_{t}+\frac{1}{2} Q_{t} F_{s}+\frac{F_{t}-F_{s}}{t-s} \text { where } s \neq t, s=\rho \text { or } \mu, \text { and } t=\rho \text { or } \mu .
\end{aligned}
$$

Therefore,

$$
x^{2} Q_{s} Q_{t}+z^{2}\left(Q_{s}-\frac{2}{s-a}\right)\left(Q_{t}-\frac{2}{t-a}\right)=\frac{1}{2} Q_{s} F_{t}+\frac{1}{2} Q_{t} F_{s}+\frac{F_{t}-F_{s}}{t-s}, \text { if } s \neq t
$$

By taking limit $s \rightarrow t$, we obtain

$$
x^{2} Q_{t}{ }^{2}+z^{2}\left(Q_{t}-\frac{2}{t-a}\right)^{2}=Q_{s} F_{t}+P_{t}
$$

Therefore,

$$
\frac{P_{s} P_{t}}{16} k_{i j}=x^{2} Q_{s} Q_{t}+z^{2}\left(Q_{s}-\frac{2}{s-a}\right)\left(Q_{t}-\frac{2}{t-a}\right)=\left\{\begin{array}{c}
P_{s}, \text { if } s=t \\
0, \text { if } s \neq t
\end{array}\right.
$$

Hence,

$$
\begin{aligned}
& {h_{\rho}}^{2}=\frac{\mu-\rho}{4(a-\rho) \rho(\rho-1)} \frac{1}{R^{2}}=\frac{1}{16}\left(-\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\rho^{2}}+\frac{\left(x^{2}+z^{2}-1\right)^{2}}{(\rho-1)^{2}}+\frac{4 z^{2}}{(\rho-a)^{2}}\right) \\
& {h_{\mu}}^{2}=\frac{\mu-\rho}{4(\mu-\rho) \mu(\mu-1)} \frac{1}{R^{2}}=\frac{1}{16}\left(-\frac{\left(x^{2}+z^{2}+1\right)^{2}}{\mu^{2}}+\frac{\left(x^{2}+z^{2}-1\right)^{2}}{(\mu-1)^{2}}+\frac{4 z^{2}}{(\mu-a)^{2}}\right)
\end{aligned}
$$

### 2.2 Flat-Ring Coordinates In Transcendental Form

Miller [4] also considers the following transcendental form of flat-ring coordinates. We write $a=k^{-2}$ with $k \in(0,1)$, and set

$$
\begin{equation*}
\rho=\operatorname{sn}^{2}(\beta, k), \quad \mu=\operatorname{sn}^{2}(\alpha, k) \tag{2.2.1}
\end{equation*}
$$

where $\operatorname{sn}(z, k)$ denotes the Jacobian elliptic function of modulus $k$ and $z$ is a complex number. We will also use the Jacobian elliptic functions $\operatorname{cn}(z, k), \operatorname{dn}(z, k)$, the complementary modulus $k^{\prime}=\sqrt{1-k^{2}}$, and the complete elliptic integrals $K, K^{\prime} . K$ and $K^{\prime}$ are the complete elliptic integrals. More details about Jacobian elliptic functions are given in the appendix.

The interval $\beta \in\left(K-i K^{\prime}, K\right)$ is mapped bijectively onto the interval $\rho \in(1, a)$. The interval $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right)$ is mapped bijectively onto the interval $\mu \in(a, \infty)$. Substituting (2.2.1) into (2.1.1) and (2.1.2), we obtain

$$
\begin{equation*}
x=\frac{\cos \phi}{R}, \quad y=\frac{\sin \phi}{R}, \quad z=\frac{i \operatorname{dn}(\alpha, k) \operatorname{dn}(\beta, k)}{k^{\prime} R}, \tag{2.2.2}
\end{equation*}
$$

where

$$
R=k \operatorname{sn}(\alpha, k) \operatorname{sn}(\beta, k)+\frac{k}{k^{\prime}} \operatorname{cn}(\alpha, k) \operatorname{cn}(\beta, k) .
$$

The coordinate surface $\alpha=\alpha_{0}$ is given by

$$
\begin{equation*}
\frac{\left(x^{2}+y^{2}+z^{2}+1\right)^{2}}{\operatorname{sn}^{2}\left(\alpha_{0}, k\right)}+\frac{\left(x^{2}+y^{2}+z^{2}-1\right)^{2}}{\operatorname{cn}^{2}\left(\alpha_{0}, k\right)}+\frac{4 k^{2} z^{2}}{\operatorname{dn}^{2}\left(\alpha_{0}, k\right)}=0, \tag{2.2.3}
\end{equation*}
$$

and the coordinate surface $\beta=\beta_{0}$ is given by

$$
\begin{equation*}
\frac{\left(x^{2}+y^{2}+z^{2}+1\right)^{2}}{\operatorname{sn}^{2}\left(\beta_{0}, k\right)}+\frac{\left(x^{2}+y^{2}+z^{2}-1\right)^{2}}{\operatorname{cn}^{2}\left(\beta_{0}, k\right)}+\frac{4 k^{2} z^{2}}{\operatorname{dn}^{2}\left(\beta_{0}, k\right)}=0 . \tag{2.2.4}
\end{equation*}
$$

The metric coefficients of (2.2.2) are

$$
\begin{equation*}
h_{\alpha}=h_{\beta^{\prime}}=\frac{k\left(\operatorname{sn}^{2}(\alpha, k)-\operatorname{sn}^{2}(\beta, k)\right)^{1 / 2}}{R}, \quad h_{\phi}=\frac{1}{R}, \tag{2.2.5}
\end{equation*}
$$

where $\beta=K+i \beta^{\prime}$ and $\beta^{\prime}$ is real.
When $\phi=0$, we have the planar coordinate system

$$
\begin{equation*}
x=\frac{1}{R}, \quad z=\frac{i \operatorname{dn}(\alpha, k) \operatorname{dn}(\beta, k)}{k^{\prime} R} . \tag{2.2.6}
\end{equation*}
$$

The transcendental form of flat-ring coordinates has the advantage that it can be extended to a coordinate system for almost the whole space $\mathbb{R}^{3}$. We can achieve this in three different ways.

### 2.2.1 First variant of transcendental flat-ring coordinates

We take $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right)$ and $\beta \in\left(K-i K^{\prime}, K+3 i K^{\prime}\right)$. Then it can be shown that (2.2.6) establishes a bijective real-analytic map between $(\alpha, \beta)$ and the region

$$
Q_{1}=\{(x, z): x>0\} \backslash\left\{(x, 0): 0<x \leq b^{-1}=\frac{1+k^{\prime}}{k}\right\}
$$

$Q_{1}$ is the right-hand half-plane with a cut from $x=0$ to $x=b^{-1}$ along the $x$-axis. The range of $\alpha, \beta$ in regions of the $(x, z)$-plane is shown in Table 2.2. Some coordinate lines of this coordinate system are depicted in Figure 2.2. The first variant of flat-ring coordinates will be useful when we discuss internal harmonics associated with coordinate surfaces $\alpha=\alpha_{0}$. When we discuss internal harmonics associated with coordinate surfaces $\beta=\beta_{0}$, then the following second variant will be more convenient.

|  | $\alpha$ | $\beta$ |
| :---: | :---: | :---: |
| $x>0, z>0, x^{2}+z^{2}<1$ | $\left(i K^{\prime}, K+i K^{\prime}\right)$ | $\left(K-i K^{\prime}, K\right)$ |
| $x>0, z>0, x^{2}+z^{2}>1$ | $\left(i K^{\prime}, K+i K^{\prime}\right)$ | $\left(K, K+i K^{\prime}\right)$ |
| $x>0, z<0, x^{2}+z^{2}>1$ | $\left(i K^{\prime}, K+i K^{\prime}\right)$ | $\left(K+i K^{\prime}, K+2 i K^{\prime}\right)$ |
| $x>0, z<0, x^{2}+z^{2}<1$ | $\left(i K^{\prime}, K+i K^{\prime}\right)$ | $\left(K+2 i K^{\prime}, K+3 i K^{\prime}\right)$ |

Table 2.2: Range of $\alpha, \beta$ in regions of ( $x, z$ )-plane (first variant)

### 2.2.2 Second variant of transcendental flat-ring coordinates

Now we take $\alpha \in\left(i K^{\prime}, 2 K+i K^{\prime}\right)$ and $\beta \in\left(K-i K^{\prime}, K+i K^{\prime}\right)$. Then (2.2.6) establishes a bijective real-analytic map between $(\alpha, \beta)$ and the region

$$
Q_{2}=\{(x, z): x>0\} \backslash\left\{(x, 0): 0<x \leq b=\frac{k}{1+k^{\prime}} \text { or } x \geq b^{-1}=\frac{k}{1-k^{\prime}}\right\} .
$$

$Q_{2}$ is the right-hand half-plane with cuts along the $x$-axis from $x=0$ to $x=\frac{k}{1+k^{\prime}}$ and from $\frac{k}{1-k^{\prime}}$ to $\infty$ The range of $\alpha, \beta$ in regions of the $(x, z)$-plane is shown in Table 2.3.

|  | $\alpha$ | $\beta$ |
| :---: | :---: | :---: |
| $x>0, z>0, x^{2}+z^{2}<1$ | $\left(i K^{\prime}, K+i K^{\prime}\right)$ | $\left(K-i K^{\prime}, K\right)$ |
| $x>0, z>0, x^{2}+z^{2}>1$ | $\left(i K^{\prime}, K+i K^{\prime}\right)$ | $\left(K, K+i K^{\prime}\right)$ |
| $x>0, z<0, x^{2}+z^{2}>1$ | $\left(K+i K^{\prime}, 2 K+i K^{\prime}\right)$ | $\left(K, K+i K^{\prime}\right)$ |
| $x>0, z<0, x^{2}+z^{2}<1$ | $\left(K+i K^{\prime}, 2 K+i K^{\prime}\right)$ | $\left(K-i K^{\prime}, K\right)$ |

Table 2.3: Range of $\alpha, \beta$ in regions of ( $x, z$ )-plane (second variant)


Figure 2.2: Coordinate lines $\alpha=\alpha_{i}, \alpha_{1}=0.3 K+i K^{\prime}, \alpha_{2}=0.5 K+i K^{\prime}, \alpha_{3}=0.7 K+i K^{\prime}$, and $\beta=\beta_{i}, \beta_{1}=K-0.5 i K^{\prime}, \beta_{2}=K, \beta_{3}=K+0.5 i K^{\prime}, \beta_{4}=K+1.5 i K^{\prime}, \beta_{5}=K+2 i K^{\prime}$, $\beta_{6}=K+2.5 i K^{\prime}$ of system (2.2.6) for $a=2$.

### 2.2.3 Third variant of transcendental flat-ring coordinates

We take $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right)$ and $\beta \in\left(K-3 i K^{\prime}, K+i K^{\prime}\right)$. Then (2.2.6) establishes a bijective real-analytic map between $(\alpha, \beta)$ and the region

$$
Q_{3}=\{(x, z): x>0\} \backslash\left\{(x, 0): 0<x \leq b=\frac{k}{1+k^{\prime}}\right\} .
$$

### 2.2.4 Reflection

Let $x>0, z>0$ with coordinates $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right), \beta \in\left(K-i K^{\prime}, K+i K^{\prime}\right)$ according to (2.2.6). These coordinates are the same in all three variants. Then the coordinates $\tilde{\alpha}, \beta$ of
the point $(x,-z)$ are

$$
\begin{array}{lll}
\tilde{\alpha}=\alpha, \quad \tilde{\beta}=2 K+2 i K^{\prime}-\beta & \text { for the first variant } \\
\tilde{\alpha}=2 K+2 i K^{\prime}-\alpha, \quad \tilde{\beta}=\beta & \text { for the second variant, } \\
\tilde{\alpha}=\alpha, \quad \tilde{\beta}=2 K-2 i K^{\prime}-\beta & \text { for the third variant. }
\end{array}
$$

We see that the reflection $(x, z) \mapsto(x,-z)$ is represented by reflections in the coordinates but in each variant in a different way. We will use these reflections when we derive harmonic functions later.

## Chapter 3

## $R$-Separation of Variables

We wish to solve Laplace's equation in flat-ring coordinates by using $R$-separation of variables. We state one theorem that we will use in the last two chapters. Let $U(x, y, z)$ be a function defined in the region

$$
\tilde{Q}_{0}=\left\{(x \cos \phi, x \sin \phi, z):(x, z) \in Q_{0}, \phi \in(-\pi, \pi)\right\} .
$$

where

$$
Q_{0}=\left\{(x, z): x>0, z>0, x^{2}+z^{2}<1\right\}
$$

In this region, we have flat-ring coordinates $\rho \in(1, a), \mu \in(a, \infty)$, and $\phi \in(-\pi, \pi)$ as we showed in chapter 2.

Lemma 4. Suppose

$$
U(x, y, z)=\left(x^{2}+y^{2}\right)^{-1 / 4} u_{1}(\rho) u_{2}(\mu) u_{3}(\phi),
$$

$\rho \in(1, a), \mu \in(a, \infty)$, and $\phi \in(-\pi, \pi)$. Assume $u_{1}, u_{2}, u_{3}$ are not identically zero. Then $U$ is harmonic if and only if there are separation constants $m, h$ such that $u_{1}$ and $u_{2}$ satisfy the algebraic Lamé equation

$$
\frac{d^{2} u}{d s^{2}}+\frac{1}{2}\left(\frac{1}{s}+\frac{1}{s-1}+\frac{1}{s-a}\right) \frac{d u}{d s}+\frac{a h-\left(m^{2}-\frac{1}{4}\right) s}{4 s(s-1)(s-a)} u=0
$$

and $u_{3}$ satisfies

$$
\frac{d^{2} u_{3}}{d \phi^{2}}+m^{2} u_{3}=0
$$

Note that $u_{1}$ and $u_{2}$ satisfy the same differential equation but on different intervals $(1, a)$ and $(a, \infty)$, respectively.

Proof. We consider cylindrical coordinates

$$
x=r \cos \phi, \quad y=r \sin \phi, \quad z=z .
$$

Let $U(x, y, z)=u(r, z, \phi)$. Then the equation

$$
\Delta U=0
$$

transforms to

$$
\begin{equation*}
\frac{\partial^{2} u}{\partial r^{2}}+\frac{1}{r} \frac{\partial u}{\partial r}+\frac{\partial^{2} u}{\partial z^{2}}+\frac{1}{r^{2}} \frac{\partial^{2} u}{\partial \phi^{2}}=0 . \tag{3.0.1}
\end{equation*}
$$

We wish to separate the variable $\phi$ from the other two variables $r$ and $z$. We assume

$$
\begin{equation*}
u(r, z, \phi)=v(r, z) u_{3}(\phi) \tag{3.0.2}
\end{equation*}
$$

We plug (3.0.2) in (3.0.1) then we have $u_{3}$ is a solution of

$$
\frac{d^{2} u_{3}}{d \phi^{2}}+m^{2} u_{3}=0
$$

where $m$ is a constant. It follows that

$$
\begin{equation*}
\frac{\partial^{2} v}{\partial r^{2}}+\frac{1}{r} \frac{\partial v}{\partial r}+\frac{\partial^{2} v}{\partial z^{2}}=\frac{m^{2}}{r^{2}} v . \tag{3.0.3}
\end{equation*}
$$

In order to separate the variables $r$ and $z$, we set

$$
v(r, z)=r^{-\frac{1}{2}} \tilde{v}(r, z) .
$$

We calculate

$$
\begin{equation*}
\frac{\partial v}{\partial r}=-\frac{1}{2} r^{\frac{-3}{2}} \tilde{v}+r^{\frac{-1}{2}} \frac{\partial \tilde{v}}{\partial r} \tag{3.0.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial^{2} v}{\partial r^{2}}=\frac{3}{4} r^{\frac{-5}{2}} \tilde{v}-r^{\frac{-3}{2}} \frac{\partial \tilde{v}}{\partial r}+r^{\frac{-1}{2}} \frac{\partial^{2} \tilde{v}}{\partial r^{2}} \tag{3.0.5}
\end{equation*}
$$

From (3.0.4), we have

$$
\begin{equation*}
\frac{1}{2} r^{\frac{-3}{2}} \tilde{v}=-\frac{\partial v}{\partial r}+r^{\frac{-1}{2}} \frac{\partial \tilde{v}}{\partial r} \tag{3.0.6}
\end{equation*}
$$

Multiplying (3.0.6) both sides by $-\frac{3}{2 r}$, we have

$$
\begin{equation*}
\frac{3}{4} r^{\frac{-5}{2}} \tilde{v}=-\frac{3}{2 r}\left(\frac{\partial v}{\partial r}+r^{\frac{-1}{2}} \frac{\partial \tilde{v}}{\partial r}\right) \tag{3.0.7}
\end{equation*}
$$

Substituting (3.0.7) in (3.0.5) and rearranging the terms, we obtain

$$
\begin{equation*}
\frac{\partial^{2} v}{\partial r^{2}}+\frac{1}{r} \frac{\partial v}{\partial r}=\frac{1}{4} r^{\frac{-5}{2}} \tilde{v}+r^{\frac{-1}{2}} \frac{\partial^{2} \tilde{v}}{\partial r^{2}} \tag{3.0.8}
\end{equation*}
$$

We substitute (3.0.8) in (3.0.3) then we have

$$
\begin{equation*}
\frac{\partial^{2} \tilde{v}}{\partial r^{2}}+\frac{\partial^{2} \tilde{v}}{\partial z^{2}}=\frac{m^{2}-\frac{1}{4}}{r^{2}} \tilde{v} \tag{3.0.9}
\end{equation*}
$$

We change variables by setting

$$
\begin{equation*}
\tilde{v}(r, z)=w(\rho, \mu) \tag{3.0.10}
\end{equation*}
$$

with flat-ring coordinates $\rho$ and $\mu$. The metric coefficients are

$$
h_{\rho}=\frac{\sqrt{\mu-\rho}}{R f_{1}(\rho)}
$$

and

$$
h_{\mu}=\frac{\sqrt{\mu-\rho}}{R f_{2}(\mu)},
$$

where

$$
\begin{aligned}
& f_{1}(\rho)=2 \sqrt{(a-\rho) \rho(\rho-1)} \\
& f_{2}(\mu)=2 \sqrt{(\mu-a) \mu(\mu-1)}
\end{aligned}
$$

and

$$
R=\left(\frac{\mu \rho}{a}\right)^{\frac{1}{2}}+\left(\frac{(\mu-1)(\rho-1)}{a-1}\right)^{\frac{1}{2}} .
$$

Equation (3.0.9) can be transformed to

$$
\begin{equation*}
\frac{1}{h_{\rho} h_{\mu}}\left[\frac{\partial}{\partial \rho}\left(\frac{h_{\rho}}{h_{\mu}} \frac{\partial w}{\partial \rho}\right)+\frac{\partial}{\partial \mu}\left(\frac{h_{\mu}}{h_{\rho}} \frac{\partial w}{\partial \mu}\right)\right]=R^{2}\left(m^{2}-\frac{1}{4}\right) w . \tag{3.0.11}
\end{equation*}
$$

We plug $h_{\rho}, h_{\mu}, f_{1}$, and $f_{2}$ from above in (3.0.11) and simplify it then we have

$$
\begin{equation*}
f_{1}(\rho) \frac{\partial}{\partial \rho}\left(f_{1}(\rho) \frac{\partial w}{\partial \rho}\right)+f_{2}(\mu) \frac{\partial}{\partial \mu}\left(f_{2}(\mu) \frac{\partial w}{\partial \mu}\right)=\left(m^{2}-\frac{1}{4}\right)(\mu-\rho) w . \tag{3.0.12}
\end{equation*}
$$

Now we separate the variables $\rho$ and $\mu$ by setting

$$
\begin{equation*}
w(\rho, \mu)=u_{1}(\rho) u_{2}(\mu) . \tag{3.0.13}
\end{equation*}
$$

We plug (3.0.13) in (3.0.12) and rearrange the terms and it follows that

$$
\begin{align*}
& \frac{f_{1}(\rho) \frac{\mathrm{d}}{\mathrm{~d} \rho}\left(f_{1}(\rho) \frac{\mathrm{d} u_{1}}{\mathrm{~d} \rho}\right)+\left(m^{2}-\frac{1}{4}\right) \rho u_{1}}{u_{1}}  \tag{3.0.14}\\
= & \frac{-f_{2}(\mu) \frac{\mathrm{d}}{\mathrm{~d} \mu}\left(f_{2}(\mu) \frac{\mathrm{d} u_{2}}{\mathrm{~d} \mu}\right)+\left(m^{2}-\frac{1}{4}\right) \mu u_{2}}{u_{2}}=h \tag{3.0.15}
\end{align*}
$$

where $h$ is a constant. Then we have

$$
\begin{equation*}
\frac{d^{2} u_{1}}{d \rho^{2}}+\frac{1}{2}\left(\frac{1}{\rho}+\frac{1}{\rho-1}+\frac{1}{\rho-a}\right) \frac{d u_{1}}{d \rho}+\frac{a h-\left(m^{2}-\frac{1}{4}\right) \rho}{4 \rho(\rho-1)(\rho-a)} u_{1}=0 \tag{3.0.16}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{d^{2} u_{2}}{d \mu^{2}}+\frac{1}{2}\left(\frac{1}{\mu}+\frac{1}{\mu-1}+\frac{1}{\mu-a}\right) \frac{d u_{2}}{d \mu}+\frac{a h-\left(m^{2}-\frac{1}{4}\right) \mu}{4 \mu(\mu-1)(\mu-a)} u_{2}=0 . \tag{3.0.17}
\end{equation*}
$$

Let $U(x, y, z)$ be a function defined in the region

$$
\tilde{Q}_{1}=\left\{(x \cos \phi, x \sin \phi, z):(x, z) \in Q_{1}, \phi \in(-\pi, \pi)\right\}
$$

where

$$
Q_{1}=\{(x, z): x>0\} \backslash\left\{(x, 0): 0<x \leq b^{-1}=\frac{1+k^{\prime}}{k}\right\}
$$

Note that $\tilde{Q}_{1}$ consists of all of $\mathbb{R}^{3}$ with the exception of a disk in the $x y$-plane centered at the origin with radius $b^{-1}$, and the half-plane $y=0, x \leq 0$. In this region, we have the first variant of transcendental flat-ring coordinates $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right), \beta \in\left(K-i K^{\prime}, K+3 i K^{\prime}\right)$, $\phi \in(-\pi, \pi)$.
Theorem 5. Suppose that $U$ has the form

$$
\begin{equation*}
U(x, y, z)=\left(x^{2}+y^{2}\right)^{-1 / 4} v_{1}(\beta) v_{2}(\alpha) v_{3}(\phi) \tag{3.0.18}
\end{equation*}
$$

where $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right), \beta \in\left(K-i K^{\prime}, K+3 i K^{\prime}\right)$, and $\phi \in(-\pi, \pi)$. Assume $v_{1}, v_{2}, v_{3}$ are not identically zero. Then $U$ is harmonic if and only if there are separation constants $m, h$ such that $v_{1}$ and $v_{2}$ both satisfy the transcendental Lamé equation

$$
\begin{equation*}
\frac{d^{2} v}{d z^{2}}+\left(h-\left(m^{2}-\frac{1}{4}\right) k^{2} \operatorname{sn}^{2}(z, k)\right) v=0 \tag{3.0.19}
\end{equation*}
$$

and $v_{3}$ satisfies

$$
\begin{equation*}
\frac{d^{2} v_{3}}{d \phi^{2}}+m^{2} v_{3}=0 \tag{3.0.20}
\end{equation*}
$$

Proof. It follows that $v_{3}$ satisfies

$$
\begin{equation*}
\frac{d^{2} v_{3}}{d \phi^{2}}+m^{2} v_{3}=0 \tag{3.0.21}
\end{equation*}
$$

from the above lemma. Substituting $\rho=s n^{2}(\beta, k)$, where $\beta \in\left(K-i K^{\prime}, K+3 i K^{\prime}\right)$ in (3.0.16) and simplifying it, we have

$$
\begin{equation*}
\frac{d^{2} v}{d z^{2}}+\left(h-\left(m^{2}-\frac{1}{4}\right) k^{2} \operatorname{sn}^{2}(z, k)\right) v=0 \tag{3.0.22}
\end{equation*}
$$

Similarly, substituting $\mu=s n^{2}(\alpha, k)$ where $\alpha \in\left(i K^{\prime}, K+i K^{\prime}\right)$ in (3.0.17) and simplifying it, we obtain

$$
\begin{equation*}
\frac{d^{2} v}{d z^{2}}+\left(h-\left(m^{2}-\frac{1}{4}\right) k^{2} \mathrm{sn}^{2}(z, k)\right) v=0 \tag{3.0.23}
\end{equation*}
$$

## Chapter 4

## The Lamé Equation

### 4.1 The Lamé Equation

We collect some known results [13] on the Lamé equation

$$
\begin{equation*}
\frac{d^{2} w}{d z^{2}}+\left(h-\nu(\nu+1) k^{2} \operatorname{sn}^{2}(z, k)\right) w=0 . \tag{4.1.1}
\end{equation*}
$$

We assume that $0<k<1$ and $\nu \geq-\frac{1}{2}$. The function $\operatorname{sn}^{2}(z, k)$ is even and has period $2 K$.

### 4.2 Lamé Periodic Fuctions

The problem we are going to consider in chapter 5 is to find values of the spectral parameter $h$ such that (4.1.1) has nontrivial $4 K$-periodic solutions. This eigenvalue problem is treated in $[13,15.5 .1]$ and $[21$, Theorem 1.2]. The eigenvalue problem splits into four eigenvalue problems with separated boundary conditions

| boundary conditions | eigenvalues | eigenfunctions | period |
| :--- | :--- | :--- | :--- |
| $w(0)=w(K)=0$ | $b_{\nu}^{2 n+2}\left(k^{2}\right)$ | $\operatorname{Es}_{\nu}^{2 n+2}\left(z, k^{2}\right)$ | $w(z+2 K)=w(z)$ |
| $w^{\prime}(0)=w(K)=0$ | $b_{\nu}^{2 n+1}\left(k^{2}\right)$ | $\operatorname{Es}_{\nu}^{2 n+1}\left(z, k^{2}\right)$ | $w(z+2 K)=-w(z)$ |
| $w(0)=w^{\prime}(K)=0$ | $a_{\nu}^{2 n+1}\left(k^{2}\right)$ | $\operatorname{Ec}_{\nu}^{2 n+1}\left(z, k^{2}\right)$ | $w(z+2 K)=-w(z)$ |
| $w^{\prime}(0)=w^{\prime}(K)=0$ | $a_{\nu}^{2 n}\left(k^{2}\right)$ | $\operatorname{Ec}_{\nu}^{2 n}\left(z, k^{2}\right)$ | $w(z+2 K)=w(z)$ |

where $n \in \mathbb{N}_{0}$ in each case denotes the number of zeros of the eigenfunction in the interval $(0, K)$. We will normalize the eigenfunctions $w=$ Es and $w=$ Ec such that

$$
\int_{0}^{K} w(t)^{2} d t=1
$$

Since $\operatorname{sn}^{2}(z, k)$ also has the imaginary period $2 i K^{\prime}$, we can consider a second eigenvalue problem where $h$ has to be determined in such a way that (4.1.1) admits nontrivial solutions with period $4 i K^{\prime}$. This problem can be reduced to the previous one by the substitutions

$$
z^{\prime}=i\left(z-K-i K^{\prime}\right), \quad h^{\prime}=\nu(\nu+1)-h .
$$

Then we obtain

$$
\begin{equation*}
\frac{d^{2} w}{d z^{\prime 2}}+\left(h^{\prime}-\nu(\nu+1) k^{\prime 2} \operatorname{sn}^{2}\left(z^{\prime}, k^{\prime}\right)\right) w=0 \tag{4.2.1}
\end{equation*}
$$

This is the same equation as (4.1.1) but with $h^{\prime}, k^{\prime}$ in place of $h, k$, respectively. Therefore, the $4 i K^{\prime}$-periodic solutions of (4.1.1) are

$$
\begin{aligned}
\mathrm{Ec}_{\nu}^{\prime n}\left(z, k^{2}\right) & :=\operatorname{Ec}_{\nu}^{n}\left(z^{\prime}, k^{\prime 2}\right) \\
\operatorname{Es}_{\nu}^{\prime n+1}\left(z, k^{2}\right) & :=\operatorname{Es}_{\nu}^{n+1}\left(z^{\prime}, k^{\prime 2}\right),
\end{aligned}
$$

with corresponding eigenvalues

$$
a_{\nu}^{\prime n}\left(k^{2}\right)=\nu(\nu+1)-a_{\nu}^{n}\left(k^{2}\right), \quad b_{\nu}^{\prime n+1}\left(k^{2}\right)=\nu(\nu+1)-b_{\nu}^{n+1}\left(k^{\prime 2}\right),
$$

where $n \in \mathbb{N}_{0} . \mathrm{Ec}^{\prime}{ }_{\nu}^{n}$ have the period $2 i K^{\prime}$ if $n$ is even, and $4 i K^{\prime}$ if $n$ is odd. $\mathrm{Es}^{\prime}{ }_{\nu}^{n+1}$ have the period $2 i K^{\prime}$ if $n$ is odd, and $4 i K^{\prime}$ if $n$ is even.

### 4.3 Frobenius Solutions

The Lamé equation (4.1.1) has a regular singularity at $z=i K^{\prime}$ with exponents $\nu+1,-\nu$. The Lamé function of the second kind $\mathrm{Fc}_{\nu}^{n}\left(z, k^{2}\right), 0<\operatorname{Re} z<2 K$, is defined as the Frobenius solution of (4.1.1) with $h=a_{\nu}^{n}\left(k^{2}\right)$ belonging to the exponent $\nu+1$. It is known that $\mathrm{Ec}_{\nu}^{n}$ and $\mathrm{Fc}_{\nu}^{n}$ are linearly independent. Similarly, $\mathrm{Fs}_{\nu}^{n+1}$ is the Frobenius solution of (3.0.19) with $h=b_{\nu}^{n+1}$ which belongs to the exponent $\nu+1$ at $z=i K^{\prime}$. Similarly, we define $\mathrm{Fc}^{\prime \prime}{ }_{\nu}$ and $\mathrm{Fs}_{\nu}^{\prime n+1}{ }^{\nu}$.

We normalize the Lamé functions of the second kind such that

$$
\begin{equation*}
W\left[\mathrm{Ec}_{\nu}^{n}, \mathrm{Fc}_{\nu}^{n}\right]=W\left[\mathrm{Es}_{\nu}^{n+1}, \mathrm{Fs}_{\nu}^{n+1}\right]=1 \tag{4.3.1}
\end{equation*}
$$

where $W$ is the Wronskian.

### 4.4 Lamé-Wangerin Functions

The problem we consider in chapter 6 is to find values of the spectral parameter $h$ such that (4.1.1) has a nontrivial solution $w(z)$ that $\operatorname{sn}(z, k)^{\frac{1}{2}} w(z)$ is bounded at $i K^{\prime}$ and $2 K+i K^{\prime}$.

The theory of this eigenvalue problem can be found in [13, 15.6] and [18, page 172].

$$
\begin{array}{l|l|l}
\text { boundary conditions } & \text { eigenvalues } & \text { eigenfunctions } \\
\hline \operatorname{sn}(z, k)^{\frac{1}{2}} w(z) \text { bounded at } z=i K^{\prime} \text { and } z=2 K+i K^{\prime} & c_{\nu}^{n}\left(k^{2}\right) & W_{\nu}^{n}\left(z, k^{2}\right)
\end{array}
$$

where $n \in \mathbb{N}_{0}$. $W_{\nu}^{n}\left(z, k^{2}\right)$ has exactly $n$ zeros on the open interval $\left(i K^{\prime}, 2 K+i K^{\prime}\right)$. $W_{\nu}^{n}\left(z, k^{2}\right)$ is an even function when $n$ is even and $W_{\nu}^{n}\left(z, k^{2}\right)$ is an odd function when $n$ is odd. We normalize the eigenfunctions $w=W_{\nu}^{n}$ such that

$$
\int_{0}^{2 K} w^{2}\left(t+i K^{\prime}\right) d t=1
$$

It is more convenient to work on the Lamé equation by only considering real numbers. To do so, we can write

$$
\begin{equation*}
z=t+i K^{\prime}, \text { where } 0<t<2 K \tag{4.4.1}
\end{equation*}
$$

It is known that

$$
\begin{equation*}
s n\left(t+i K^{\prime}, k\right)=\frac{1}{k(\operatorname{sn}(t, k))} \tag{4.4.2}
\end{equation*}
$$

Substituting (4.4.1) and (4.4.2) in (4.1.1), we have

$$
\begin{equation*}
\frac{d^{2} w}{d t^{2}}+\left(h-\frac{\nu(\nu+1)}{s n^{2}(t, k)}\right) w=0, \text { where } 0<t<2 K \tag{4.4.3}
\end{equation*}
$$

More details about Lamé-Wangerin functions are given in the appendix.

## Chapter 5

## Flat-Ring Harmonic Functions

### 5.1 Internal Flat-Ring Harmonics

For internal flat-ring harmonics, we use the first variant of the flat-ring coordinates. We have $\alpha=t+i K^{\prime}$ with $0<t<K$. If $\alpha_{0}=t_{0}+i K^{\prime}$ for some fixed value $t_{0} \in(0, K)$, then $\alpha=\alpha_{0}$ describes a closed surface. The closed surface looks like a flat ring (see Figure 5.1). So we call this kind of closed surface "flat ring". Let


Figure 5.1: flat-ring, $a=2, \alpha=0.8 K+i K^{\prime}$

$$
\begin{equation*}
D_{1}=\left\{\mathbf{r} \in \mathbb{R}^{3}: \frac{\left(\|r\|^{2}+1\right)^{2}}{\operatorname{sn}^{2}\left(\alpha_{0}, k\right)}+\frac{\left(\|r\|^{2}-1\right)^{2}}{\mathrm{cn}^{2}\left(\alpha_{0}, k\right)}+\frac{4 k^{2} z^{2}}{\operatorname{dn}^{2}\left(\alpha_{0}, k\right)}>0\right\}, \tag{5.1.1}
\end{equation*}
$$

where $\alpha_{0} \in\left(i K^{\prime}, K+i K^{\prime}\right)$. $D_{1}$ describes the interior of a flat-ring for a fixed $\alpha_{0}$. Internal flat-ring harmonics are harmonic functions of the separated form (3.0.18) which are harmonic in the region described by $D_{1}$ for a fixed $\alpha_{0}$. We can consider the region forms by the union of the regions described by $D_{1}$ for each $\alpha_{0}$, which is all $\mathbb{R}^{3}$ except the $z$-axis. Suppose that $v_{1}$ and $v_{2}$ are solutions of the Lamé equation (3.0.19) defined on the strip $0<$ Real $v_{s}<2 K$, where $s=1,2$, and $v_{3}$ is a solution of (3.0.20). We use the first variant of the transcendental flat-ring coordinate system and $-\pi<\phi<\pi$ then $U(x, y, z)$ defined by (3.0.18) is a harmonic function in $\tilde{Q}_{1}$. We want this function to be harmonic on $\mathbb{R}^{3}$ except the $z$-axis. Clearly, we need $m \in \mathbb{Z}$, and we choose $u_{3}(\phi)=e^{i m \phi}$ (alternatively, we could use $\cos (m \phi), m=0,1,2, \ldots$ and $\sin (m \phi), m=1,2,3, \ldots)$. Then we have to require that the function $v_{1}(\beta) v_{2}(\alpha)$ is analytic in the right-hand half plane $x>0, z \in \mathbb{R}$. We know $v_{1}(\beta) v_{2}(\alpha)$ is always analytic in the quadrant $x>0, z>0$. When we use the third variant, we can analytically extend this function to the quadrant $x>0, z<0$ across the segment $b=\frac{k}{1-k^{\prime}}>x>0, z=0$. We want the first and third extension to be the same. In order to make the first and third extension the same, we need that $v_{1}$ to be periodic with period $4 i K^{\prime}$. Therefore, we take $v_{1}=\mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}$ or $v_{1}=\mathrm{Es}^{\prime \prime n+1}{ }_{|m|-\frac{1}{2}}$. When we use the second variant, we can analytically extend this function to the quadrant $x>0, z<0$ across the segment $b=\frac{k}{1+k^{\prime}}<x<b^{-1}=\frac{k}{1-k^{\prime}}$, $z=0$. We want the first and second extension to be the same. In order to make the first and second extension the same, we need that $v_{1}$ has the same parity with respect to $K+i K^{\prime}$ as $v_{2}$. Therefore, $v_{1}$ has to be a constant multiple of $v_{2}$. We can take $v_{1}=v_{2}$. We define internal flat-ring harmonics by

$$
\begin{align*}
\operatorname{Gc}_{m}^{n}(x, y, z) & =\left(x^{2}+y^{2}\right)^{-1 / 4} \mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\beta, k^{2}\right) \mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\alpha, k^{2}\right) e^{i m \phi},  \tag{5.1.2}\\
\operatorname{Gs}_{m}^{n+1}(x, y, z) & =\left(x^{2}+y^{2}\right)^{-1 / 4} \operatorname{Es}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\beta, k^{2}\right) \mathrm{Es}^{\prime n+1}{ }_{|m|-\frac{1}{2}}\left(\alpha, k^{2}\right) e^{i m \phi}, \tag{5.1.3}
\end{align*}
$$

where $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$. Note that the Lamé functions $\operatorname{Ec}^{\prime}(\zeta)$ and $\operatorname{Es}^{\prime}(\zeta)$ are analytic on the strip $0<\operatorname{Re} \zeta<2 K$, and that $\alpha$ and $\beta$ lie in this strip.

We collect some properties of internal harmonics in the following theorem. We refer to the inversion at the unit sphere given by

$$
\sigma(\mathbf{r})=\|\mathbf{r}\|^{-2} \mathbf{r}
$$

Theorem 6. The internal flat-ring harmonics $\mathrm{Gc}_{m}^{n}$ and $\mathrm{Gs}_{m}^{n+1}$ are harmonic in all of $\mathbb{R}^{3}$ except for the $z$-axis. They have the following symmetry regarding inversion at the unit sphere:

$$
\begin{aligned}
\mathrm{Gc}_{m}^{n}(\sigma(\mathbf{r})) & =\|\mathbf{r}\| \mathrm{Gc}_{m}^{n}(\mathbf{r}), \\
\mathrm{Gs}_{m}^{n+1}(\sigma(\mathbf{r})) & =-\|\mathbf{r}\| \mathrm{Gs}_{m}^{n+1}(\mathbf{r})
\end{aligned}
$$

and regarding reflection at the xy-plane

$$
\begin{aligned}
\operatorname{Gc}_{m}^{n}(x, y,-z) & =(-1)^{n+1} \operatorname{Gc}_{m}^{n}(x, y, z), \\
\operatorname{Gs}_{m}^{n+1}(x, y,-z) & =(-1)^{n} \operatorname{Gs}_{m}^{n+1}(x, y, z)
\end{aligned}
$$

Proof. From our consideration at the beginning of this section we know that $U=$ Gc and $U=$ Gs are harmonic functions on $\mathbb{R}^{3}$ except for the $z$-axis and the two circles centered at the origin with radii $b^{-1}$ and $b$ on the $x y$-plane, respectively. The function $U$ is bounded in a neighborhood of the two circles. Therefore, the two circles are removable singularities of $U$ (see [3], Theorem XIII, page 271). Hence, $U$ is harmonic on $\mathbb{R}^{3}$ except for the $z$-axis.

The inversion $\sigma$ in the upper half-space $z>0$ is expressed in flat-ring coordinates by the reflection of $\beta$ at $K: \beta \mapsto 2 K-\beta$. The Lamé functions $\mathrm{Ec}^{\prime}$ remain unchanged under this reflection while the Lamé functions Es' change sign. Similarly, the reflection $z \mapsto-z$ is expressed by reflection of $\beta$ at $K+i K^{\prime}: \beta \mapsto 2\left(K+i K^{\prime}\right)-\beta$. The Lame functions $\mathrm{Ec}_{\nu}^{\prime n}$, $\mathrm{Es}^{\prime n}{ }_{\nu}$ remain unchanged under this reflection if $n$ is odd but change sign when $n$ is even.

### 5.2 The Dirichlet Problem

From the known orthogonality and completeness properties of the periodic Lamé functions, we obtain the following theorem.
Theorem 7. The system of functions

$$
(8 \pi)^{-1 / 2} \operatorname{Ec}^{\prime}{ }_{|m|-\frac{1}{2}}^{n}\left(\beta, k^{2}\right) e^{i m \phi},(8 \pi)^{-1 / 2} \operatorname{Es}_{|m|-\frac{1}{2}}^{\prime n+1}\left(\beta, k^{2}\right) e^{i m \phi}, \quad m \in \mathbb{Z}, n \in \mathbb{N}_{0}
$$

is an orthonormal basis in the Hilbert space

$$
H_{1}=L^{2}\left(\left(K-i K^{\prime}, K+3 i K^{\prime}\right) \times(-\pi, \pi)\right) .
$$

We use the internal flat-ring harmonics to solve the Dirichlet problem for harmonic functions in the region $D_{1}$ defined in (5.1.1).

Theorem 8. Let $f$ be a function defined on the boundary $\partial D_{1}$ of the region $D_{1}$. Suppose that $f$ is represented in flat-ring coordinates as

$$
\left(x^{2}+y^{2}\right)^{1 / 4} f(x, y, z)=g(\beta, \phi), \quad \beta=K+i \beta^{\prime}, \beta^{\prime} \in(-K, 3 K), \phi \in(-\pi, \pi)
$$

such that $g \in H_{1}$. For all $m \in \mathbb{Z}$ and $n \in \mathbb{N}_{0}$, define

$$
\begin{aligned}
c_{m}^{n} & :=\frac{1}{8 \pi \mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\alpha_{0}, k^{2}\right)} \int_{-\pi}^{\pi} \int_{-K}^{3 K} g(\beta, \phi) \mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}^{n}\left(\beta, k^{2}\right) e^{-i m \phi} d \beta^{\prime} d \phi \\
& =\frac{1}{8 \pi\left\{\mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\alpha_{0}, k^{2}\right)\right\}^{2}} \int_{\partial D_{1}} \frac{1}{h_{\beta^{\prime}}(\mathbf{r})} f(\mathbf{r}) \mathrm{Gc}_{-m}^{n}(\mathbf{r}) d S(\mathbf{r})
\end{aligned}
$$

and

$$
\begin{aligned}
d_{m}^{n+1} & :
\end{aligned}=\frac{1}{8 \pi \mathrm{Es}_{|m|-\frac{1}{2}}^{\prime n+1}\left(\alpha_{0}, k^{2}\right)} \int_{-\pi}^{\pi} \int_{-K}^{3 K} g(\beta, \phi) \mathrm{Es}^{\prime \prime}{ }_{|m|-\frac{1}{2}}^{n+1}\left(\beta, k^{2}\right) e^{-i m \phi} d \beta^{\prime} d \phi
$$

Then the function

$$
\begin{equation*}
u(\mathbf{r})=\sum_{m \in \mathbb{Z}} \sum_{n=0}^{\infty}\left(c_{m}^{n} \mathrm{Gc}_{m}^{n}(\mathbf{r})+d_{m}^{n+1} \mathrm{Gs}_{m}^{n+1}(\mathbf{r})\right) \tag{5.2.1}
\end{equation*}
$$

is harmonic in $D_{1}$ and it assumes the boundary values $f$ on $\partial D_{1}$ in the weak sense. The infinite series in (5.2.1) converges absolutely and uniformly in compact subsets of $D_{1}$.

Proof. Using flat-ring coordinates, we can write surface integrals over $\partial D_{1}$ as double integrals:

$$
\begin{aligned}
\int_{\partial D_{1}} f(\mathbf{r}) d S(\mathbf{r}) & =\int_{-K}^{3 K} \int_{-\pi}^{\pi} h_{\beta^{\prime}} h_{\phi}\left(x^{2}+y^{2}\right)^{-1 / 4} g(\beta, \phi) d \beta^{\prime} d \phi \\
& =\int_{-K}^{3 K} \int_{-\pi}^{\pi} h_{\beta^{\prime}}\left(x^{2}+y^{2}\right)^{1 / 4} g(\beta, \phi) d \beta^{\prime} d \phi
\end{aligned}
$$

with the metric coefficients $h_{\beta^{\prime}}, h_{\phi}$ given in (2.2.5). This shows that the two formulas given for $c_{m}^{n}$ and $d_{m}^{n+1}$ agree.

The rest of the proof is similar to the proof of Theorem 6.3 in [6].

### 5.3 External Flat-Ring Harmonics

External flat-ring harmonics are harmonic functions $U$ of the form (3.0.18) which are harmonic outside of all flat rings (5.1.1). Therefore, they are harmonic in $\mathbb{R}^{3}$ except for the annulus $b^{2} \leq x^{2}+y^{2} \leq b^{-2}$ in the $(x, y)$-plane. It is clear that $m$ must be an integer and arguing as at the beginning of Section 5.1, we see that $v_{1}$ must have period $4 i K^{\prime}$. Since the function $U$ has to be analytic along the $z$-axis, we require that $v_{2}$ is a solution of the Lamé equation (3.0.19) which is bounded (actually must converge to 0 ) when $\alpha$ approached $i K^{\prime}$. These are the Frobenius solutions of the Lamé equation that we mentioned in Section 4.2. Thus we define external flat-ring harmonics by

$$
\begin{align*}
\mathrm{Hc}_{m}^{n}(x, y, z) & =\left(x^{2}+y^{2}\right)^{-1 / 4} \mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\beta, k^{2}\right) \mathrm{Fc}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\alpha, k^{2}\right) e^{i m \phi},  \tag{5.3.1}\\
\operatorname{Hs}_{m}^{n+1}(x, y, z) & =\left(x^{2}+y^{2}\right)^{-1 / 4} \operatorname{Es}^{\prime \prime}{ }_{|m|-\frac{1}{2}}\left(\beta, k^{2}\right) \mathrm{Fs}^{\prime n+1}{ }_{|m|-\frac{1}{2}}\left(\alpha, k^{2}\right) e^{i m \phi}, \tag{5.3.2}
\end{align*}
$$

where $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$. We use the first variant of the transcendental flat-ring coordinates (we can also use the third variant but not the second one).

Theorem 9. The external flat-ring harmonics $\mathrm{Hc}_{m}^{n}, \mathrm{Hs}_{m}^{n+1}$ are harmonic in all of $\mathbb{R}^{3}$ except for the closed annulus in the xy-plane centered at the origin with inner radius $\frac{k}{1+k^{\prime}}$ and outer radius $\frac{k}{1-k^{\prime}}$. They have the same symmetry properties as the internal flat-ring harmonics; see Theorem 6.

Proof. From our consideration at the beginning of this chapter we know that the functions $U=\mathrm{Hc}_{m}^{n}$ and $U=\mathrm{Hs}_{m}^{n+1}$ are harmonic on $\mathbb{R}^{3}$ except for the $z$-axis and the annulus $b^{2} \leq$ $x^{2}+y^{2} \leq b^{-2}$ in the $(x, y)$-plane. Since $U$ is bounded the $z$-axis is a removable singularity (see [3], Theorem VI, page 335). Therefore, $U$ is harmonic on $\mathbb{R}^{3}$ except for the annulus $b^{2} \leq x^{2}+y^{2} \leq b^{-2}$ in the $(x, y)$-plane. (It appears that $U$ has a continuous extension on this annulus but not an analytic extension.)

We now show that external harmonics admit an integral representation in terms of internal harmonics.

Theorem 10. Let $\alpha_{0} \in\left(i K^{\prime}, K+i K^{\prime}\right)$, $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$, and let $\mathbf{r}^{\prime}$ be a point outside $\bar{D}_{1}$, where $D_{1}$ is given by (5.1.1). Then

$$
\begin{equation*}
\operatorname{Hc}_{m}^{n}\left(\mathbf{r}^{\prime}\right)=\frac{1}{\left\{\mathrm{Ec}_{|m|-\frac{1}{2}}^{\prime n}\left(\alpha_{0}, k^{2}\right)\right\}^{2}} \int_{\partial D_{1}} \frac{\mathrm{Gc}_{m}^{n}(\mathbf{r})}{h_{\alpha}(\mathbf{r}) 4 \pi\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|} d S(\mathbf{r}) \tag{5.3.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{Hs}_{m}^{n+1}\left(\mathbf{r}^{\prime}\right)=\frac{1}{\left\{\operatorname{Es}_{|m|-\frac{1}{2}}^{\prime \prime n+1}\left(\alpha_{0}, k^{2}\right)\right\}^{2}} \int_{\partial D_{1}} \frac{\mathrm{Gs}_{m}^{n+1}(\mathbf{r})}{h_{\alpha}(\mathbf{r}) 4 \pi\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|} d S(\mathbf{r}) \tag{5.3.4}
\end{equation*}
$$

Proof. Let $D$ be an open bounded subset of $\mathbb{R}^{3}$ with smooth boundary. For $u, v \in C^{2}(\bar{D})$, Green's formula states that

$$
\begin{equation*}
\int_{D}(u \Delta v-v \Delta u) d \mathbf{r}=\int_{\partial D}\left(u \frac{\partial v}{\partial \nu}-v \frac{\partial u}{\partial \nu}\right) d S \tag{5.3.5}
\end{equation*}
$$

where $\frac{\partial u}{\partial \nu}$ is the outward normal derivative of $u$ on the boundary $\partial D$ of $D$. We apply (5.3.5) to $D=D_{1}, u=G=\mathrm{Gc}_{m}^{n}$, and

$$
\begin{equation*}
v(\mathbf{r})=\frac{1}{4 \pi\left\|\mathbf{r}^{\prime}-\mathbf{r}\right\|} \tag{5.3.6}
\end{equation*}
$$

Since $\Delta u=\Delta v=0$ on $D_{1}$ we obtain

$$
\begin{equation*}
0=\int_{\partial D_{1}}\left(G \frac{\partial v}{\partial \nu}-v \frac{\partial G}{\partial \nu}\right) d S \tag{5.3.7}
\end{equation*}
$$

Let $B_{r}(\mathbf{q})$ denote the open ball centered at $\mathbf{q} \in \mathbb{R}^{3}$ with radius $r>0$. We apply (5.3.5) a second time with $D=B_{R}(\mathbf{0})-\bar{D}_{1}-B_{\epsilon}\left(\mathbf{r}^{\prime}\right)$ with large $R$ and small $\epsilon>0$. Choose $u=H=\mathrm{Hc}_{m}^{n}$ and $v$ as in (5.3.6). Note that $\Delta u=\Delta v=0$ on $D$. By a standard argument, taking the limit $\epsilon \rightarrow 0$, we obtain

$$
\begin{equation*}
H\left(\mathbf{r}^{\prime}\right)=\int_{\partial B_{R}(\mathbf{0})}\left(H \frac{\partial v}{\partial \nu}-v \frac{\partial H}{\partial \nu}\right) d S-\int_{\partial D_{1}}\left(H \frac{\partial v}{\partial \nu}-v \frac{\partial H}{\partial \nu}\right) d S \tag{5.3.8}
\end{equation*}
$$

where, in the second integral, $\frac{\partial}{\partial \nu}$ denotes again the derivative in the direction of the outward normal as in (5.3.7). As $R \rightarrow \infty$, the first integral on the right-hand side of (5.3.8) tends to 0 (see [8], page 109), so we obtain

$$
\begin{equation*}
H\left(\mathbf{r}^{\prime}\right)=-\int_{\partial D_{1}}\left(H \frac{\partial v}{\partial \nu}-v \frac{\partial H}{\partial \nu}\right) d S \tag{5.3.9}
\end{equation*}
$$

Set $E(\zeta)=\mathrm{Ec}^{\prime \prime}{ }_{|m|-\frac{1}{2}}(\zeta)$ and $F(\zeta)=\mathrm{Fc}^{\prime \prime}{ }_{|m|-\frac{1}{2}}(\zeta)$. We multiply (5.3.7) by $F\left(\alpha_{0}\right)$, then multiply (5.3.9) by $E\left(\alpha_{0}\right)$ and add these equations. Using the definitions (5.1.2) and (5.3.1) of internal and external flat-ring harmonics and cancelling terms, we find

$$
\begin{equation*}
E\left(\alpha_{0}\right) H\left(\mathbf{r}^{\prime}\right)=\int_{\partial D_{1}} v\left(E\left(\alpha_{0}\right) \frac{\partial H}{\partial \nu}-F\left(\alpha_{0}\right) \frac{\partial G}{\partial \nu}\right) d S . \tag{5.3.10}
\end{equation*}
$$

Since flat-ring coordinates are orthogonal, the normal derivative and the partial derivative with respect to $\alpha$ are related by

$$
\frac{\partial}{\partial \nu}=\frac{1}{h_{\alpha}} \frac{\partial}{\partial \alpha}
$$

where $h_{\alpha}$ is given in (2.2.5). Let $\mathbf{r}=(x, y, z) \in \partial D_{1}$ with flat-ring coordinates $\alpha, \beta, \phi$. Then

$$
\begin{aligned}
E\left(\alpha_{0}\right) \frac{\partial H}{\partial \nu}(\mathbf{r}) & -F\left(\alpha_{0}\right) \frac{\partial G}{\partial \nu}(\mathbf{r}) \\
= & E\left(\alpha_{0}\right) \frac{\partial\left(\left(x^{2}+y^{2}\right)^{-1 / 4}\right)}{\partial \nu} F\left(\alpha_{0}\right) E(\beta) e^{i m \phi} \\
& +E\left(\alpha_{0}\right)\left(x^{2}+y^{2}\right)^{-1 / 4} h_{\alpha}^{-1} F^{\prime}\left(\alpha_{0}\right) E(\beta) e^{i m \phi} \\
& \quad-F\left(\alpha_{0}\right) \frac{\partial\left(\left(x^{2}+y^{2}\right)^{-1 / 4}\right)}{\partial \nu} E\left(\alpha_{0}\right) E(\beta) e^{i m \phi} \\
& \quad-F\left(\alpha_{0}\right)\left(x^{2}+y^{2}\right)^{-1 / 4} h_{\alpha}^{-1} E^{\prime}\left(\alpha_{0}\right) E(\beta) e^{i m \phi} \\
= & h_{\alpha}^{-1}\left(x^{2}+y^{2}\right)^{-1 / 4}\left\{E\left(\alpha_{0}\right) F^{\prime}\left(\alpha_{0}\right)-E^{\prime}\left(\alpha_{0}\right) F\left(\alpha_{0}\right)\right\} E(\beta) e^{i m \phi} .
\end{aligned}
$$

We now use the Wronskian (4.3.1) and obtain

$$
\begin{equation*}
E\left(\alpha_{0}\right) \frac{\partial H}{\partial \nu}(\mathbf{r})-F\left(\alpha_{0}\right) \frac{\partial G}{\partial \nu}(\mathbf{r})=\frac{1}{h_{\alpha}(\mathbf{r}) E\left(\alpha_{0}\right)} G(\mathbf{r}) . \tag{5.3.11}
\end{equation*}
$$

When we substitute (5.3.11) in (5.3.10) we obtain (5.3.3). The proof of (5.3.4) is similar.

### 5.4 A Fundamental Solution

We obtain the expansion of (5.3.6) in internal and external flat-ring harmonics of the first kind by combining Theorems 8 and 10 .

Theorem 11. Let $\mathbf{r}, \mathbf{r}^{\prime} \in \mathbb{R}^{3}$ with flat-ring coordinates $\alpha, \alpha^{\prime} \in\left(i K^{\prime}, K+i K^{\prime}\right)$, respectively. If $\operatorname{Re} \alpha>\operatorname{Re} \alpha^{\prime}$ then

$$
\begin{equation*}
\frac{1}{\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|}=\frac{1}{2} \sum_{m \in \mathbb{Z}} \sum_{n=0}^{\infty}\left(\mathrm{Gc}_{m}^{n}(\mathbf{r}) \mathrm{Hc}_{-m}^{n}\left(\mathbf{r}^{\prime}\right)+\mathrm{Gs}_{m}^{n+1}(\mathbf{r}) \mathrm{Hs}_{-m}^{n+1}\left(\mathbf{r}^{\prime}\right)\right) \tag{5.4.1}
\end{equation*}
$$

Proof. We pick $\alpha_{0} \in\left(i K^{\prime}, K+i K^{\prime}\right)$ such that $\operatorname{Re} \alpha^{\prime}<\operatorname{Re} \alpha_{0}<\operatorname{Re} \alpha$, and consider the domain $D_{1}$ defined in (5.1.1). The function $f(\mathbf{q}):=\frac{1}{\left\|\mathbf{q}-\mathbf{r}^{\prime}\right\|}$ is harmonic on an open set containing $\bar{D}_{1}$. Therefore, by Theorem 8, we have

$$
\begin{equation*}
f(\mathbf{r})=\sum_{m \in \mathbb{Z}} \sum_{n=0}^{\infty}\left(c_{m}^{n} \mathrm{Gc}_{m}^{n}(\mathbf{r})+d_{m}^{n+1} \mathrm{Gs}_{m}^{n+1}(\mathbf{r})\right) \tag{5.4.2}
\end{equation*}
$$

where $c_{m}^{n}$ and $d_{m}^{n+1}$ can be evaluated by Theorem 10:

$$
c_{n}^{m}=\frac{1}{2} \mathrm{Hc}_{-m}^{n}\left(\mathbf{r}^{\prime}\right), \quad d_{m}^{n+1}=\frac{1}{2} \mathrm{Hs}_{-m}^{n+1}\left(\mathbf{r}^{\prime}\right) .
$$

We use $h_{\beta^{\prime}}=h_{\alpha}$. Thus we obtain (5.4.1).

## Chapter 6

## Peanut Harmonic Functions

### 6.1 External Peanut Harmonics

It is more convenient to consider external peanut harmonics first. We use the second variant of the flat-ring coordinates. We have $\beta=K+i t$ with $-K^{\prime}<t<K^{\prime}$. If $\beta_{0}=K+i t_{0}$ for some fixed value $t_{0} \in\left(-K^{\prime}, K^{\prime}\right)$, then $\beta=\beta_{0}$ describes a closed surface (if we add two points on the $z$-axis, where the coordinate system is not valid). The closed surface looks like a peanut (see Figure 6.1). So we call this kind of closed surface "peanut". The interior of


Figure 6.1: peanut, $a=2, \beta=K+0.7 i K^{\prime}$
a peanut is given by $t<t_{0}$ and the exterior by $t>t_{0}$ (if we add a disk or the exterior of a disk in the $(x, y)$-plane, where the coordinate system is not valid).

We also want to describe the region $D_{2}$ (one peanut) in Cartesian coordinates as we did for the region $D_{1}$ (one flat ring) earlier. We consider the surface given by

$$
\begin{equation*}
f(\mathbf{r}):=\frac{\left(\|\mathbf{r}\|^{2}+1\right)^{2}}{\operatorname{sn}^{2}\left(\beta_{0}, k\right)}+\frac{\left(\|\mathbf{r}\|^{2}-1\right)^{2}}{\operatorname{cn}^{2}\left(\beta_{0}, k\right)}+\frac{4 k^{2} z^{2}}{\operatorname{dn}^{2}\left(\beta_{0}, k\right)}=0 \tag{6.1.1}
\end{equation*}
$$

Then $f(\sigma(\mathbf{r}))=\|\mathbf{r}\|^{-4} f(\mathbf{r})$, where $\sigma(\mathbf{r})=\|\mathbf{r}\|^{-2} \mathbf{r}$ denote an inversion at the unit sphere. It follows that if $\mathbf{r}$ is on the surface, then the mirror point $\sigma(\mathbf{r})$ is on the surface. Therefore, the surface given by (6.1.1) has two parts, namely the surface $\beta=\beta_{0}$ and $\beta=\bar{\beta}_{0}$. Unless $t_{0}=0$, one surface is interior to the unit sphere and the other one is its mirror image and it lies exterior to the unit sphere. This is also clear from the observation that $f(\mathbf{r})$ remains the same if we replace $\beta_{0}$ by its conjugate. Between the surfaces we have $f(\mathbf{r})>0$ and interior to the smaller surface and exterior to the larger surface we have $f(\mathbf{r})<0$. We can check this by considering $\mathbf{r}=0$ and $\|\mathbf{r}\|=1$.

We consider three cases:

- If $t_{0}=0$ then the surface $\beta=\beta_{0}$ is the unit sphere and its interior $D_{2}$ is the open unit ball

$$
B=\{\mathbf{r}:\|r\|<1\}
$$

Notice that $f(\mathbf{r})$ is not well-defined in this case.

- If $-K^{\prime}<t_{0}<0$ then the surface $\beta=\beta_{0}$ is in $B$. Then we have

$$
D_{2}=B \cap\{\mathbf{r}: f(\mathbf{r})<0\} .
$$

- If $0<t_{0}<K^{\prime}$ then the surface $\beta=\beta_{0}$ is exterior to the unit sphere and

$$
D_{2}=B \cup\{\mathbf{r}: f(\mathbf{r})>0\} .
$$

External peanut harmonics are harmonic functions $U$ of the form (3.0.18) that are harmonic outside a peanut described by $D_{2}$ for a fixed $\beta$. We can consider the region formed by the union of the regions outside each peanut described by $D_{2}$, which is $\mathbb{R}^{3}$ except the disc centered at the origin with radius $b=\frac{k}{1+k^{\prime}}$ on the $x y$-plane. Suppose that $v_{1}$ and $v_{2}$ are solutions of the Lamé equation (3.0.19) defined on the strip $0<$ Real $v_{s}<2 K$, where $s=1,2$, and $v_{3}$ is a solution of (3.0.20). We use the second variant of the transcendental flat-ring coordinate system and $-\pi<\phi<\pi$ then $U(x, y, z)$ defined by (3.0.18) is a harmonic function in $\tilde{Q}_{2}$. We want this function to be harmonic on $\mathbb{R}^{3}$ except for the disc centered at the origin with radius $b=\frac{k}{1+k^{\prime}}$ on the $x y$-plane. Clearly, we need $m \in \mathbb{Z}$ and we choose $u_{3}(\phi)=e^{i m \phi}$ (alternatively, we could use $\cos (m \phi), m=0,1,2, \ldots$ and $\sin (m \phi)$, $m=1,2,3, \ldots)$. Then we have to require that the function $v_{1}(\beta) v_{2}(\alpha)$ is analytic in the
right-hand half plane $x>0, z \in \mathbb{R}$ except the segment between the origin and $b$ on the $x$-axis. As harmonic functions are smooth, the function $v_{2}$ has to be bounded at the points $i K^{\prime}$ and $2 K+i K^{\prime}$. Therefore, we take $v_{2}=W_{|m|-\frac{1}{2}}^{n}$. We know $v_{1}(\beta) v_{2}(\alpha)$ is always analytic in the quadrant $x>0, z>0$. When we use the first variant, we can analytically extend this function to the quadrant $x>0, z<0$ across the ray $x>b^{-1}=\frac{k}{1-k^{\prime}}, z=0$. When we use the second variant, we can analytically extend this function to the quadrant $x>0, z<0$ across the segment $b^{-1}>x>b, z=0$. We want the first and second extension to be the same. In order to make the first and second extension the same, we need $v_{1}$ to have the same parity with respect to $K+i K^{\prime}$ as $v_{2}$. Therefore, $v_{1}$ has to be a constant multiple of $v_{2}$. We can take $v_{1}=v_{2}$. Thus we define external peanut harmonics by

$$
\begin{equation*}
H_{m}^{n}(x, y, z)=\left(x^{2}+y^{2}\right)^{-1 / 4} W_{|m|-\frac{1}{2}}^{n}\left(\beta, k^{2}\right) W_{|m|-\frac{1}{2}}^{n}\left(\alpha, k^{2}\right) e^{i m \phi} . \tag{6.1.2}
\end{equation*}
$$

We collect some properties of external peanut harmonics in the following theorem.
Theorem 12. The external peanut harmonics $H_{m}^{n}$ are harmonic on $\mathbb{R}^{3}$ except for the disc centered at the origin with radius $b=\frac{k}{1+k^{\prime}}$ on the xy-plane. Regarding reflection at the xy-plane, we have

$$
H_{m}^{n}(x, y,-z)=(-1)^{n} H_{m}^{n}(x, y, z)
$$

Proof. From our consideration at the beginning of this section we know that $H_{m}^{n}$ is a harmonic function on $\mathbb{R}^{3}$ except for the disc centered at the origin with the radius $b$ on the $x y$-plane, $z$-axis, and the circle centered at the origin with radius $b^{-1}$ on the $x y$-plane. $H_{m}^{n}$ is bounded in a neighborhood of the circle. Therefore, the circle is a removable singularity of $H_{m}^{n}$ (see [3], Theorem XIII, page 271). Since $H_{m}^{n}$ is bounded in a neighborhood of $z$-axis, the $z$-axis is a removable singularity (see [3], Theorem VI, page 335). Hence, $H_{m}^{n}$ are harmonic on $\mathbb{R}^{3}$ except for the disc centered at the origin with radius $b=\frac{k}{1+k^{\prime}}$ on the $x y$-plane.

The reflection $z \mapsto-z$ is expressed by reflection of $\alpha$ at $K+i K^{\prime}: \alpha \mapsto 2\left(K+i K^{\prime}\right)-\alpha$. The Lamé-Wangerin function $W_{\nu}^{n}$ remains unchanged under this reflection if $n$ is even but changes sign if $n$ is odd.

### 6.2 Internal Peanut Harmonics

Internal peanut harmonics are harmonic functions $U$ of the form (3.0.18) which are harmonic inside a peanut described by $D_{2}$ for a fixed $\beta$. We can consider the region formed by the union of the regions inside each peanut described by $D_{2}$, which is $\mathbb{R}^{3}$ except for the infinite annulus with inner radius $b=\frac{k}{1-k^{\prime}}$ and outer radius $\infty$ in the $(x, y)$-plane. It is clear that $m$ must be an integer and arguing as at the beginning of section 6.1.

We define internal peanut harmonics by

$$
\begin{equation*}
G_{m}^{n}(x, y, z)=\left(x^{2}+y^{2}\right)^{-1 / 4} W_{|m|-\frac{1}{2}}^{n}\left(2 K-\beta, k^{2}\right) W_{|m|-\frac{1}{2}}^{n}\left(\alpha, k^{2}\right) e^{i m \phi} \tag{6.2.1}
\end{equation*}
$$

where $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$.
Theorem 13. The internal peanut harmonics $G_{m}^{n}$ are harmonic on $\mathbb{R}^{3}$ except for the infinite annulus with inner radius $b^{-1}=\frac{k}{1-k^{\prime}}$ and outer radius $\infty$ in the ( $x, y$ )-plane. They have the following symmetry regarding inversion at the unit sphere (Kelvin inversion):

$$
G_{m}^{n}(\sigma(\mathbf{r}))=\|\mathbf{r}\| H_{m}^{n}(\mathbf{r}),
$$

and regarding reflection at the $x y$-plane

$$
G_{m}^{n}(x, y,-z)=(-1)^{n} G_{m}^{n}(x, y, z)
$$

Proof. The inversion $\sigma$ for $x>0$ is expressed in the flat-ring coordinates by the reflection of $\beta$ at $K: \beta \rightarrow 2 K-\beta$. The region of all internal peanuts is under the inversion $\sigma$ of $D_{2}$ for all $\beta_{0}$.

The reflection $z \mapsto-z$ is expressed by reflection of $\alpha$ at $K+i K^{\prime}: \alpha \mapsto 2\left(K+i K^{\prime}\right)-\alpha$. The Lamé-Wangerin function $W_{\nu}^{n}$ remains unchanged under this reflection if $n$ is even but changes sign if $n$ is odd.

We now show that external harmonics admit an integral representation in terms of internal harmonics. Assume

$$
\begin{equation*}
W\left[E_{\nu}^{n}, F_{\nu}^{n}\right](z)=w, \tag{6.2.2}
\end{equation*}
$$

where $W$ is the Wronskian, $E_{\nu}^{n}(z)=W_{\nu}^{n}(z), F_{\nu}^{n}(z)=W_{\nu}^{n}(2 K-z)$, and $w$ is a constant.
Theorem 14. Let $\beta_{0} \in\left(K-i K^{\prime}, K+i K^{\prime}\right)$, $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$, and let $\mathbf{r}^{\prime}$ be a point outside $\bar{D}_{2}$, where $D_{2}$ is given in section 6.1. Then

$$
\begin{equation*}
H_{m}^{n}\left(\mathbf{r}^{\prime}\right)=\frac{w}{\left\{W_{|m|-\frac{1}{2}}^{n}\left(\beta_{0}, k^{2}\right)\right\}^{2}} \int_{\partial D_{2}} \frac{G_{m}^{n}(\mathbf{r})}{h_{\beta}(\mathbf{r}) 4 \pi\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|} d S(\mathbf{r}) \tag{6.2.3}
\end{equation*}
$$

Proof. Let $D$ be an open bounded subset of $\mathbb{R}^{3}$ with smooth boundary. For $u, v \in C^{2}(\bar{D})$, Green's formula states that

$$
\begin{equation*}
\int_{D}(u \Delta v-v \Delta u) d \mathbf{r}=\int_{\partial D}\left(u \frac{\partial v}{\partial \nu}-v \frac{\partial u}{\partial \nu}\right) d S \tag{6.2.4}
\end{equation*}
$$

where $\frac{\partial u}{\partial \nu}$ is the outward normal derivative of $u$ on the boundary $\partial D$ of $D$. We apply (6.2.4) to $D=D_{2}, u=G=G_{m}^{n}$, and

$$
\begin{equation*}
v(\mathbf{r})=\frac{1}{4 \pi\left\|\mathbf{r}^{\prime}-\mathbf{r}\right\|} \tag{6.2.5}
\end{equation*}
$$

Since $\Delta u=\Delta v=0$ on $D_{2}$ we obtain

$$
\begin{equation*}
0=\int_{\partial D_{2}}\left(G \frac{\partial v}{\partial \nu}-v \frac{\partial G}{\partial \nu}\right) d S \tag{6.2.6}
\end{equation*}
$$

Let $B_{r}(\mathbf{q})$ denote the open ball centered at $\mathbf{q} \in \mathbb{R}^{3}$ with radius $r>0$. We apply (6.2.4) a second time with $D=B_{R}(\mathbf{0})-\bar{D}_{2}-B_{\epsilon}\left(\mathbf{r}^{\prime}\right)$ with large $R$ and small $\epsilon>0$. Choose $u=H=H_{m}^{n}$ and $v$ as in (6.2.5). Note that $\Delta u=\Delta v=0$ on $D$. By a standard argument, taking the limit $\epsilon \rightarrow 0$, we obtain

$$
\begin{equation*}
H\left(\mathbf{r}^{\prime}\right)=\int_{\partial B_{R}(\mathbf{0})}\left(H \frac{\partial v}{\partial \nu}-v \frac{\partial H}{\partial \nu}\right) d S-\int_{\partial D_{2}}\left(H \frac{\partial v}{\partial \nu}-v \frac{\partial H}{\partial \nu}\right) d S \tag{6.2.7}
\end{equation*}
$$

where, in the second integral, $\frac{\partial}{\partial \nu}$ denotes again the derivative in the direction of the outward normal as in (6.2.6). As $R \rightarrow \infty$, the first integral on the right-hand side of (6.2.7) tends to 0 (see [8], page 109), so we obtain

$$
\begin{equation*}
H\left(\mathbf{r}^{\prime}\right)=-\int_{\partial D_{2}}\left(H \frac{\partial v}{\partial \nu}-v \frac{\partial H}{\partial \nu}\right) d S \tag{6.2.8}
\end{equation*}
$$

Set $E(\zeta)=W_{|m|-\frac{1}{2}}^{n}(\zeta)$ and $F(\zeta)=W_{|m|-\frac{1}{2}}^{n}(2 K-\zeta)$. We multiply (6.2.6) by $F\left(\beta_{0}\right)$, then multiply (6.2.8) by $E\left(\beta_{0}\right)$ and add these equations. Using the definitions (6.1.2) and (6.2.1) of internal and external flat-ring harmonics and cancelling terms, we find

$$
\begin{equation*}
E\left(\beta_{0}\right) H\left(\mathbf{r}^{\prime}\right)=\int_{\partial D_{2}} v\left(E\left(\beta_{0}\right) \frac{\partial H}{\partial \nu}-F\left(\beta_{0}\right) \frac{\partial G}{\partial \nu}\right) d S \tag{6.2.9}
\end{equation*}
$$

Since flat-ring coordinates are orthogonal, the normal derivative and the partial derivative with respect to $\beta$ are related by

$$
\frac{\partial}{\partial \nu}=\frac{1}{h_{\beta}} \frac{\partial}{\partial \beta},
$$

where $h_{\alpha}$ is given in (2.2.5). Let $\mathbf{r}=(x, y, z) \in \partial D_{2}$ with flat-ring coordinates $\alpha, \beta, \phi$. Then

$$
\begin{aligned}
E\left(\beta_{0}\right) \frac{\partial H}{\partial \nu}(\mathbf{r})- & F\left(\beta_{0}\right) \frac{\partial G}{\partial \nu}(\mathbf{r}) \\
= & E\left(\beta_{0}\right) \frac{\partial\left(\left(x^{2}+y^{2}\right)^{-1 / 4}\right)}{\partial \nu} F\left(\beta_{0}\right) E(\alpha) e^{i m \phi} \\
& +E\left(\beta_{0}\right)\left(x^{2}+y^{2}\right)^{-1 / 4} h_{\beta}^{-1} F^{\prime}\left(\beta_{0}\right) E(\alpha) e^{i m \phi} \\
& \quad-F\left(\beta_{0}\right) \frac{\partial\left(\left(x^{2}+y^{2}\right)^{-1 / 4}\right)}{\partial \nu} E\left(\beta_{0}\right) E(\alpha) e^{i m \phi} \\
& \quad-F\left(\beta_{0}\right)\left(x^{2}+y^{2}\right)^{-1 / 4} h_{\beta}^{-1} E^{\prime}\left(\beta_{0}\right) E(\alpha) e^{i m \phi} \\
= & h_{\beta}^{-1}\left(x^{2}+y^{2}\right)^{-1 / 4}\left\{E\left(\beta_{0}\right) F^{\prime}\left(\beta_{0}\right)-E^{\prime}\left(\beta_{0}\right) F\left(\beta_{0}\right)\right\} E(\alpha) e^{i m \phi}
\end{aligned}
$$

We now use the Wronskian (6.2.2) and obtain

$$
\begin{equation*}
E\left(\beta_{0}\right) \frac{\partial H}{\partial \nu}(\mathbf{r})-F\left(\beta_{0}\right) \frac{\partial G}{\partial \nu}(\mathbf{r})=\frac{w}{h_{\beta}(\mathbf{r}) E\left(\beta_{0}\right)} G(\mathbf{r}) \tag{6.2.10}
\end{equation*}
$$

When we substitute (6.2.10) in (6.2.9) we obtain (6.2.3).

### 6.3 The Dirichlet Problem

From the known orthogonality and completeness properties of the Lamé-Wangerin functions, we obtain the following theorem.
Theorem 15. The system of functions

$$
(8 \pi)^{-1 / 2} W_{|m|-\frac{1}{2}}^{n}\left(\alpha, k^{2}\right) e^{i m \phi}, \quad m \in \mathbb{Z}, n \in \mathbb{N}_{0}
$$

is an orthonormal basis in the Hilbert space

$$
H_{2}=L^{2}\left(\left(i K^{\prime}, 2 K+i K^{\prime}\right) \times(-\pi, \pi)\right) .
$$

where $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$.
We use the external peanut harmonics of the second kind to solve the Dirichlet problem for harmonic functions in the region $D_{2}$ defined in section 6.1.
Theorem 16. Let $f$ be a function defined on the boundary $\partial D_{2}$ of the region $D_{2}$. Suppose that $f$ is represented in flat-ring coordinates as

$$
\left(x^{2}+y^{2}\right)^{1 / 4} f(x, y, z)=g\left(\alpha^{\prime}, \phi\right), \quad \alpha^{\prime}=\alpha+i K^{\prime}, \quad \alpha \in(0,2 K), \quad \phi \in(-\pi, \pi)
$$

such that $g \in H_{2}$. For all $m \in \mathbb{Z}$ and $n \in \mathbb{N}_{0}$, we define

$$
\begin{aligned}
c_{m}^{n} & :
\end{aligned}=\frac{1}{8 \pi W_{|m|-\frac{1}{2}}^{n}\left(\beta_{0}, k^{2}\right)} \int_{-\pi}^{\pi} \int_{0}^{2 K} g\left(\alpha^{\prime}, \phi\right) W_{|m|-\frac{1}{2}}^{n}\left(\alpha^{\prime}, k^{2}\right) e^{-i m \phi} d \alpha d \phi .
$$

Then the function

$$
\begin{equation*}
u(\mathbf{r})=\sum_{m \in \mathbb{Z}} \sum_{n=0}^{\infty}\left(c_{m}^{n} G_{m}^{n}(\mathbf{r})\right) \tag{6.3.1}
\end{equation*}
$$

is harmonic in $D_{2}$ and it assumes the boundary values $f$ on $\partial D_{2}$ in the weak sense. The infinite series in (6.3.1) converges absolutely and uniformly in compact subsets of $D_{2}$.
Proof. Using flat-ring coordinates we can write surface integrals over $\partial D_{2}$ as double integrals:

$$
\begin{aligned}
\int_{\partial D_{2}} f(\mathbf{r}) d S(\mathbf{r}) & =\int_{0}^{2 K} \int_{-\pi}^{\pi} h_{\alpha} h_{\phi}\left(x^{2}+y^{2}\right)^{-1 / 4} g\left(\alpha^{\prime}, \phi\right) d \alpha d \phi \\
& =\int_{0}^{2 K} \int_{-\pi}^{\pi} h_{\alpha}\left(x^{2}+y^{2}\right)^{1 / 4} g\left(\alpha^{\prime}, \phi\right) d \alpha d \phi
\end{aligned}
$$

with the metric coefficients $h_{\alpha}, h_{\phi}$ given in (2.2.5). This shows that the formula given for $c_{m}^{n}$ agrees. The rest of the proof is similar to the proof of Theorem 6.3 in [6] .

### 6.4 A Fundamental Solution

We obtain the expansion of (6.2.5) in internal and external peanut harmonics by combining Theorems 16 and 14.

Theorem 17. Let $\mathbf{r}, \mathbf{r}^{\prime} \in \mathbb{R}^{3}$ with flat-ring coordinates $\beta, \beta^{\prime} \in\left(K-i K^{\prime}, K+i K^{\prime}\right)$, respectively. If $\operatorname{Re} \beta>\operatorname{Re} \beta^{\prime}$ then

$$
\begin{equation*}
\frac{1}{\left\|\mathbf{r}-\mathbf{r}^{\prime}\right\|}=\frac{1}{2} \sum_{m \in \mathbb{Z}} \sum_{n=0}^{\infty}\left(G_{m}^{n}(\mathbf{r}) H_{-m}^{n}\left(\mathbf{r}^{\prime}\right)\right), \tag{6.4.1}
\end{equation*}
$$

where $m \in \mathbb{Z}, n \in \mathbb{N}_{0}$.
Proof. We pick $\beta_{0} \in\left(K-i K^{\prime}, K+i K^{\prime}\right)$ such that $\operatorname{Re} \beta^{\prime}<\operatorname{Re} \beta_{0}<\operatorname{Re} \beta$, and consider the domain $D_{2}$ defined in section 6.1. The function $f(\mathbf{q}):=\frac{1}{\left\|\mathbf{q}-\mathbf{r}^{\prime}\right\|}$ is harmonic on an open set containing $\bar{D}_{2}$. Therefore, by Theorem 16, we have

$$
\begin{equation*}
f(\mathbf{r})=\sum_{m \in \mathbb{Z}} \sum_{n=0}^{\infty}\left(c_{m}^{n} G_{m}^{n}(\mathbf{r})\right), \tag{6.4.2}
\end{equation*}
$$

where $c_{m}^{n}$ can be evaluated by Theorem 14:

$$
c_{n}^{m}=\frac{1}{2} \mathrm{Hc}_{-m}^{n}\left(\mathbf{r}^{\prime}\right) .
$$

We used $h_{\beta^{\prime}}=h_{\alpha}$. Thus we obtain (6.4.1).

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## Appendix

## General Descriptions of the Functions sn $u$, cn $u$, dn $u$

Let $0<k^{2}<1$,

$$
\begin{aligned}
K(k):= & \int_{0}^{1} \frac{\mathrm{~d} x}{\sqrt{1-x^{2}} \sqrt{1-k^{2} x^{2}}}, \\
& K^{\prime}(k):=K\left(k^{\prime}\right)
\end{aligned}
$$

(I) The function $\mathrm{sn} u$ is a doubly-periodic function of $u$ with periods $4 K, 2 i K^{\prime}$. It is analytic except at the points congruent to $i K^{\prime}$ or to $2 K+i K^{\prime}\left(\bmod .4 K, 2 i K^{\prime}\right)$; these points are simple poles, the residues at the first set all being $k^{-1}$ and the residues at the second set all being $-k^{-1}$; and the function has a simple zero at all points congruent to $0\left(\bmod .2 K, 2 i K^{\prime}\right)$.
(II) The function cnu is a doubly-periodic function of $u$ with periods $4 K$ and $2 K+2 i K^{\prime}$. It is analytic except at points congruent to $i K^{\prime}$ or to $2 K+i K^{\prime}\left(\bmod .4 K, 2 K+2 i K^{\prime}\right.$; these points are simple poles, the residues at the the first sent being $-i k^{-1}$, and the residues at the second set being $i k^{-1}$; and the function has a simple zero at all points congruent to $K$ (mod. $2 K, 2 i K^{\prime}$ ).
(III) The function $\operatorname{dn} u$ is a doubly-periodic function of $u$ with periods $2 K$ and $4 i K^{\prime}$. It is analytic except at points congruent to $i K^{\prime}$ or to $3 i K^{\prime}$ (mod. $2 K, 4 i K^{\prime}$ ); these points are simple poles, the residues at the first set being $-i$, and the residues at the second set being $i$; and the function has a simple zero at all points congruent to $K+i K^{\prime}\left(\bmod .2 K, 2 i K^{\prime}\right)$.

It may be observed that $\operatorname{sn} u, \mathrm{cn} u$, and $\operatorname{dn} u$ are the only functions satisfying the above descriptions, respectively, see [22, page 504].

## Lamé-Wangerin Functions

## The eigenvalue problem

We consider the differential equation

$$
\begin{equation*}
-y^{\prime \prime}+q(x) y=\lambda y, \quad 0<x \leq b<\infty . \tag{6.4.3}
\end{equation*}
$$

We assume that $q:(0, b] \rightarrow \mathbb{R}$ is a continuous function. Additionally, we assume that $q(x)$ has the Laurent series representation

$$
\begin{equation*}
q(x)=\frac{1}{x^{2}} \sum_{k=0}^{\infty} q_{k} x^{k}, \quad 0<x<\epsilon, \tag{6.4.4}
\end{equation*}
$$

for some $\epsilon>0$. The differential equation (6.4.3) has a regular singularity at $x=0$ with indicial equation

$$
r(r-1)-q_{0}=0
$$

The solutions (called exponents) of this equation are

$$
r_{1}=\frac{1}{2}-\sqrt{\frac{1}{4}+q_{0}}, \quad r_{2}=\frac{1}{2}+\sqrt{\frac{1}{4}+q_{0}} .
$$

If $q_{0}=-\frac{1}{4}$ we have a double root $r_{1}=r_{2}=\frac{1}{2}$. If $q_{0}>-\frac{1}{4}$, there are two distinct real roots $r_{1}<\frac{1}{2}<r_{2}$. If $q_{0}<-\frac{1}{4}$, the roots $r_{1}, r_{2}$ are conjugate complex.

For every $\lambda \in \mathbb{C}$, equation (6.4.3) has a fundamental system $y_{1}(x, \lambda), y_{2}(x, \lambda)$ with the property

$$
\begin{aligned}
& y_{1}(x, \lambda)=x^{r_{1}} \sum_{k=0}^{\infty} c_{k} x^{k}+c y_{2}(x, \lambda) \ln x, \quad 0<x<\epsilon \\
& y_{2}(x, \lambda)=x^{r_{2}} \sum_{k=0}^{\infty} d_{k} x^{k}, \quad 0<x<\epsilon
\end{aligned}
$$

If $r_{2}-r_{1} \neq 0,1,2, \ldots$ then $c_{0}=d_{0}=1, c=0$. If $r_{1}=r_{2}$ then $c_{0}=0, d_{0}=1, c=-1$. If $r_{2}-r_{1}=m_{0}$ a positive integer, then $c_{0}=d_{0}=1, c_{m_{0}}=0$. This is a well-known result; see [18, Satz 1, page 147]. We have

$$
d_{0}=1, \quad d_{1}=\frac{q_{1}}{2 r_{2}}, \quad d_{2}=\frac{2 r_{2} q_{2}-2 \lambda r_{2}+q_{1}^{2}}{4 r_{2}\left(2 r_{2}+1\right)}
$$

In the following we assume that

$$
q_{0} \geq-\frac{1}{4}
$$

Let $0<\beta \leq \pi$. We call $\lambda \in \mathbb{C}$ an eigenvalue if

$$
\begin{equation*}
\cos \beta y_{2}(b, \lambda)=\sin \beta y_{2}^{\prime}(b, \lambda), \quad \prime=\frac{d}{d x} . \tag{6.4.5}
\end{equation*}
$$

We call $y_{2}(x, \lambda)$ a corresponding eigenfunction.

## Main results

In this section we argue essentially the same way as in the theory of regular Sturm-Liouville problems.

Lemma 1. For every $\lambda, \mu \in \mathbb{C}$ we have

$$
\lim _{x \rightarrow 0^{+}} y_{2}(x, \lambda) y_{2}^{\prime}(x, \mu)= \begin{cases}0 & \text { if } q_{0}>-\frac{1}{4} \\ \frac{1}{2} & \text { if } q_{0}=-\frac{1}{4}\end{cases}
$$

Proof. If $q_{0}>-\frac{1}{4}$ then $r_{2}>\frac{1}{2}$ and if $q_{0}=-\frac{1}{4}$ then $r_{2}=\frac{1}{2}$.
Lemma 2. If $\lambda_{1}, \lambda_{2}$ are two distinct real eigenvalues, then

$$
\int_{0}^{b} y_{2}\left(x, \lambda_{1}\right) y_{2}\left(x, \lambda_{2}\right) d x=0
$$

Proof. Let $u_{j}(x)=y_{2}\left(x, \lambda_{j}\right)$. For $0<\delta<b$ we obtain

$$
\begin{aligned}
\left(\lambda_{1}-\lambda_{2}\right) \int_{\delta}^{b} u_{1}(x) u_{2}(x) d x & =\left[u_{1}(x) u_{2}^{\prime}(x)-u_{1}^{\prime}(x) u_{2}(x)\right]_{\delta}^{b} \\
& =-u_{1}(\delta) u_{2}^{\prime}(\delta)+u_{1}^{\prime}(\delta) u_{2}(\delta) .
\end{aligned}
$$

By Lemma 1 , as $\delta \rightarrow 0^{+}$,

$$
\left(\lambda_{1}-\lambda_{2}\right) \int_{0}^{b} u_{1}(x) u_{2}(x) d x=0 .
$$

Since $\lambda_{1} \neq \lambda_{2}$ the desired statement follows.
Lemma 3. All eigenvalues are real.
Proof. Let $\lambda$ be a complex eigenvalue. Setting $u_{1}=y_{2}(x, \lambda)$ and $u_{2}=y_{2}(x, \bar{\lambda})$ we obtain arguing as in the proof of Lemma 2,

$$
2 i \operatorname{Im} \lambda \int_{0}^{b}\left|u_{1}(x)\right|^{2} d x=0
$$

Therefore, $\operatorname{Im} \lambda=0$.

We introduce the Prüfer angle $\theta(x, \lambda)$ and Prüfer radius $r(x, \lambda)>0$ for $0<x \leq b$ by setting

$$
y_{2}(x, \lambda)=r(x, \lambda) \sin \theta(r, \lambda), \quad y_{2}^{\prime}(x, \lambda)=r(x, \lambda) \cos \theta(r, \lambda) .
$$

The Prüfer angle is defined to be between 0 and $\pi / 2$ for small positive $x$ and by continuity otherwise. For small positive $x$ we have

$$
\theta(x, \lambda)=\arctan \frac{y_{2}(x \cdot \lambda)}{y_{2}^{\prime}(x, \lambda)}
$$

This shows that $\theta(x, \lambda)$ is an analytic function of $x$ in a neighborhood of 0 . We calculate

$$
\begin{equation*}
\theta(x)=a_{1} x+a_{2} x^{2}+a_{3} x^{3}+\ldots \tag{6.4.6}
\end{equation*}
$$

with

$$
\begin{gathered}
a_{1}=\frac{1}{r_{2}}, \quad a_{2}=-\frac{q_{1}}{2 r_{2}^{3}} \\
a_{3}=\frac{12 \lambda r_{2}^{3}-12 r_{2}^{3} q_{2}-8 r_{2}^{3}+9 q_{1}^{2} r_{2}-4 r_{2}^{2}+3 q_{1}^{2}}{12 r_{2}^{5}\left(2 r_{2}+1\right)}
\end{gathered}
$$

The Prüfer angle satisfies the first order differential equation

$$
\begin{equation*}
\theta^{\prime}=\cos ^{2} \theta+(\lambda-q(x)) \sin ^{2} \theta \tag{6.4.7}
\end{equation*}
$$

The following lemma collects properties of $\theta(x, \lambda)$. These properties are analogous to those of the Prüfer angle for regular Sturm-Liouville problems; see [19].
Lemma 4. (a) When $\theta\left(x_{0}, \lambda\right)=n \pi$ for some integer $n$ and $0 \leq x_{0} \leq b$, then $\theta(x, \lambda)<n \pi$ for $0<x<x_{0}$ and $\theta(x, \lambda)>n \pi$ for $x_{0}<x \leq b$.
(b) The function $\lambda \rightarrow \theta(x, \lambda)$ is continuous and strictly increasing for every $0<x \leq b$.
(c) $\lim _{\lambda \rightarrow-\infty} \theta(b, \lambda)=0$.
(d) $\lim _{\lambda \rightarrow \infty} \theta(b, \lambda)=\infty$.

Proof. (a) If $\theta\left(x_{0}, \lambda\right)=n \pi$ then (6.4.7) implies $\theta^{\prime}\left(x_{0}, \lambda\right)=1$.
(b) Let $\lambda_{1}<\lambda_{2}$. The expansion (6.4.6) shows that $\frac{d^{k}}{d x^{k}} \theta\left(0, \lambda_{1}\right)=\frac{d^{k}}{d x^{k}} \theta\left(0, \lambda_{2}\right)$ for $k=0,1,2$ but $\frac{d^{3}}{d x^{3}} \theta\left(0, \lambda_{1}\right)<\frac{d^{3}}{d x^{3}} \theta\left(0, \lambda_{2}\right)$. Therefore, we can choose $0<\delta<x$ such that $\theta\left(\delta, \lambda_{1}\right)<\theta\left(\delta, \lambda_{2}\right)$. Then $\theta\left(x, \lambda_{1}\right)<\theta\left(x, \lambda_{2}\right)$ follows from a standard differential inequality; see [16, Theorem 2.1, page 144].
(c) The expansion (6.4.6) shows that there is $0<\delta<b$ such that $0<\theta(\delta, 0)<\frac{\pi}{2}$. Let $\Theta(x, \lambda)$ be the solution of (6.4.7) with initial condition $\Theta(\delta, \lambda)=\frac{\pi}{2}$. It follows from (b) that $\theta(\delta, \lambda)<\frac{\pi}{2}$ for $\lambda \leq 0$. Therefore, $0<\theta(b, \lambda)<\Theta(b, \lambda)$ for $\lambda \leq 0$. It is known from regular Sturm-Liouville theory that $\lim _{\lambda \rightarrow-\infty} \Theta(b, \lambda)=0$. This proves (c).
(d) We choose $0<\delta<b$. Let $\Theta(x, \lambda)$ be the solution of (6.4.7) determined by $\Theta(\delta, \lambda)=0$. Then $\Theta(b, \lambda)<\theta(b, \lambda)$ for all real $\lambda$. It is known from regular Sturm-Liouville theory that $\lim _{\lambda \rightarrow \infty} \Theta(b, \lambda)=\infty$. This proves (d).

Theorem 18. The eigenvalues (solutions of (6.4.5)) form an increasing sequence

$$
\lambda_{0}<\lambda_{1}<\lambda_{2}<\ldots
$$

converging to infinity. The corresponding eigenfunctions $y_{2}\left(x, \lambda_{n}\right)$ have exactly $n$ zeros in the open interval $(0, b)$.

Proof. By Lemma 3, all eigenvalues are real. A real number $\lambda$ is an eigenvalue if and only if there exists $n=0,1,2, \ldots$ such that

$$
\begin{equation*}
\theta(b, \lambda)=\beta+n \pi . \tag{6.4.8}
\end{equation*}
$$

By Lemma 4, for every $n=0,1,2, \ldots$, equation (6.4.8) has exactly one solution $\lambda=\lambda_{n}$. By Lemma $4(\mathrm{a}), y_{2}\left(x, \lambda_{n}\right)$ has exactly $n$ zeros in $(0, b)$.

Let $C[0, b]$ be the vector space of continuous real-valued functions defined on $[0, b]$ endowed with the inner product

$$
\langle f, g\rangle=\int_{0}^{b} f(x) g(x) d x
$$

This is a pre-Hilbert space. We could also work in the Hilbert space $L^{2}(0, b)$ but we want to avoid usage of the Lebesgue integral in this section. An orthogonal sequence $\left\{f_{n}\right\}, f_{n} \neq 0$, in a pre-Hilbert space $H$ is called complete if the closed linear span of the sequence is $H$. The sequence is complete if and only if every function $f \in H$ can be written as a generalized Fourier series

$$
f=\sum_{n} \frac{\left\langle f, f_{n}\right\rangle}{\left\langle f_{n}, f_{n}\right\rangle} f_{n}
$$

with convergence in $H$.
By Lemma 2, the sequence $\left\{y_{2}\left(x, \lambda_{n}\right)\right\}_{n=0}^{\infty}$ is an orthogonal sequence in $C[0, b]$. We are going to show that this sequence is complete (this result is missing in [17].).

For $\lambda \in \mathbb{C}$ let $y_{3}(x, \lambda)$ be the solution of (6.4.3) determined by initial conditions

$$
y_{3}(b, \lambda)=\sin \beta, \quad y_{3}^{\prime}(b, \lambda)=\cos \beta .
$$

$\lambda$ is an eigenvalue if and only if $y_{2}$ and $y_{3}$ are linearly dependent. Choose $\lambda^{*} \in \mathbb{R}$ such that $\lambda^{*}$ is not an eigenvalue, and define the Green's function

$$
G(x, t)=W^{-1} \begin{cases}y_{2}(x) y_{3}(t), & 0<x<t \leq b \\ y_{2}(t) y_{3}(x), & 0<t \leq x \leq b\end{cases}
$$

where $y_{j}(x)=y_{j}\left(x, \lambda^{*}\right)$ for $j=2,3$, and $W$ is the Wronskian

$$
W=y_{3} y_{2}^{\prime}-y_{3}^{\prime} y_{2}
$$

$W$ is a nonzero constant.

Lemma 5. If we set $G(x, t)=0$ for $x=0$ or $t=0$, then $G:[0, b] \times[0, b] \rightarrow \mathbb{R}$ is continuous. Proof. We use the estimates

$$
\left|y_{2}(x)\right| \leq B_{2} x^{r_{2}}, \quad\left|y_{3}(x)\right| \leq B_{3} x^{r_{1}}|\ln x|, \quad \text { for } 0<x \leq b,
$$

where $B_{2}, B_{3}$ are constants. If $0<t \leq x \leq b$, we have

$$
\left|y_{2}(t) y_{3}(x)\right| \leq B_{2} B_{3} t^{r_{2}} x^{r_{1}} \leq B_{2} B_{3} x^{r_{1}+r_{2}}|\ln x|=B_{2} B_{3} x|\ln x| .
$$

This shows that $G$ is continuous on the triangle $0 \leq x \leq t \leq b$. Since $G(x, t)=G(t, x), G$ is continuous on $[0, b]^{2}$.

We consider the integral operator $T: C[0, b] \rightarrow C[0, b]$ defined by

$$
T f(x)=\int_{0}^{b} G(x, t) f(t) d t
$$

Since $G$ is continuous and $G(x, t)=G(t, x)$, we have the following expansion theorem; see [19].

Theorem 19. The operator $T$ has a sequence of nonzero eigenvalues $\left\{\mu_{n}\right\}$ with corresponding orthogonal eigenfunctions $f_{n}$ such that the closed linear span of $\left\{f_{n}\right\}$ contains the range of $T$.

Lemma 6. $\lambda$ is an eigenvalue in the sense of (6.4.5) if and only if $\left(\lambda-\lambda^{*}\right)^{-1}$ is an eigenvalue of $T$. The corresponding eigenfunctions are the same.

Proof. Let $f \in C[0, b]$ and set $y=T f$. Then, for $0<x \leq b$,

$$
y(x)=W^{-1}\left(\int_{0}^{x} y_{2}(t) f(t) d t\right) y_{3}(x)+W^{-1}\left(\int_{x}^{b} y_{3}(t) f(t) d t\right) y_{2}(x)
$$

Differentiating we obtain

$$
y^{\prime}(x)=W^{-1}\left(\int_{0}^{x} y_{2}(t) f(t) d t\right) y_{3}^{\prime}(x)+W^{-1}\left(\int_{x}^{b} y_{3}(t) f(t) d t\right) y_{2}^{\prime}(x) .
$$

Differentiating once more we find

$$
y^{\prime \prime}(x)=-f(x)+\left(q(x)-\lambda^{*}\right) y(x) .
$$

If $T f=\left(\lambda-\lambda^{*}\right)^{-1} f$ with $f \neq 0$ then setting $y=T f \neq 0$ we obtain $-y^{\prime \prime}+q(x) y=\lambda y$. We have

$$
y(b)=W^{-1}\left(\int_{0}^{b} y_{2}(t) f(t) d t\right) y_{3}(b), \quad y^{\prime}(b)=W^{-1}\left(\int_{0}^{b} y_{2}(t) f(t) d t\right) y_{3}^{\prime}(b) .
$$

Therefore, $y$ satisfies the boundary condition $\cos \beta y(b)=\sin \beta y^{\prime}(b)$. If $r_{1}<0$ then $y$ must be a constant multiple of $y_{2}\left(x, \lambda_{n}\right)$ for some $n$ because $y=T f$ is continuous on $[0, b]$. Therefore, $\lambda=\lambda_{n}$ for some $n=0,1,2 \ldots$ If $r_{1} \geq 0$ then $y(x)=O(x)+O\left(x^{r_{2}}\right)$ as $x \rightarrow 0^{+}$so we reach the same conclusion.

Conversely, we wish to show that $T y_{2}\left(x, \lambda_{n}\right)=\left(\lambda_{n}-\lambda^{*}\right)^{-1} y_{2}\left(x, \lambda_{n}\right)$. This follows from the equations

$$
\begin{aligned}
& \int_{0}^{x} y_{2}(t) y_{2}\left(t, \lambda_{n}\right) d t=\left(\lambda_{n}-\lambda^{*}\right)^{-1}\left(y_{2}\left(x, \lambda_{n}\right) y_{2}^{\prime}(x)-y_{2}^{\prime}\left(x, \lambda_{n}\right) y_{2}(x)\right) \\
& \int_{x}^{b} y_{3}(t) y_{2}\left(t, \lambda_{n}\right) d t=\left(\lambda_{n}-\lambda^{*}\right)^{-1}\left(-y_{2}\left(x, \lambda_{n}\right) y_{3}^{\prime}(x)+y_{2}^{\prime}\left(x, \lambda_{n}\right) y_{3}(x)\right)
\end{aligned}
$$

Theorem 20. The eigenfunctions $\left\{y_{2}\left(x, \lambda_{n}\right)\right\}_{n=0}^{\infty}$ are complete in $C[0, b]$.
Proof. It follows from Lemma 6 and Theorem 19 that the closed linear span of the eigenfunctions $y\left(x, \lambda_{n}\right)$ contains the range of $T$. Therefore, it is enough to show that the range of $T$ is dense in $C[0, b]$. Let $g \in C^{2}[0, b]$ with compact support in $(0, b)$. Then $f(x):=$ $-g^{\prime \prime}(x)+\left(q(x)-\lambda^{*}\right) g(x)$ is continuous on $[0, b]$ with compact support in $(0, b)$. The function $y=T f$ satisfies $f(x)=-y^{\prime \prime}(x)+\left(q(x)-\lambda^{*}\right) y(x)$, so $v(x):=g(x)-y(x)$ satisfies $-v^{\prime \prime}(x)+\left(q(x)-\lambda^{*}\right) v(x)=0$. Now $\cos \beta v(v)=\sin \beta v^{\prime}(b)$, and $y(x)$ is a constant multiple of $y_{2}(x)$ for $x$ close to 0 . Therefore, $g=y$ and $g$ is in the range of $T$. The set of functions $g \in C^{2}[0, b]$ with compact support in $(0, b)$ is dense in $C[0, b]$. The proof is complete.

## Connectios to Weyl's theory of singular Sturm-Liouville problems

We see that $y_{2} \in L^{2}\left(0, \frac{1}{2} \epsilon\right)$ for all $q_{0} \in \mathbb{R}$, and $y_{1} \in L^{2}\left(0, \frac{1}{2} \epsilon\right)$ if and only if $q_{0}<\frac{3}{4}$. From this observation we obtain the following lemma.

Lemma 7. Equation (6.4.3) is in the limit-point case at $x=0$ if and only if $q_{0} \geq \frac{3}{4}$.
Lemma 7 is a special case of [18, Satz 1, page 152].
The point $b$ is a regular end point. At this end point we pose the boundary condition

$$
\begin{equation*}
\cos \beta y(b)=\sin \beta y^{\prime}(b) \tag{6.4.9}
\end{equation*}
$$

If $q_{0} \geq \frac{3}{4}$ we do not need a boundary condition at $x=0$. If $-\frac{1}{4} \leq q_{0}<\frac{3}{4}$, we add a boundary condition that forces eigenfunctions to be constant multiples of $y_{2}$. Let $H=L^{2}(0, b)$. If $q_{0} \geq \frac{3}{4}$, we introduce the linear operator $A: H \supset D(A) \rightarrow H$ by $A y=-y^{\prime \prime}+q y$ with domain

$$
D(A)=\left\{y: y, y^{\prime} \in A C[\delta, b] \text { for all } 0<\delta<b, y,-y^{\prime \prime}+q y \in H, y \text { satisfies (6.4.9) }\right\} .
$$

If $-\frac{1}{4} \leq q_{0}<\frac{3}{4}$, we add the boundary condition

$$
\lim _{x \rightarrow 0^{+}}\left(y_{2}(x) y^{\prime}(x)-y_{2}^{\prime}(x) y(x)\right)=0 .
$$

It is known that $A$ is a self-adjoint operator. For the operator $T: H \rightarrow H$ defined in Section 2 we have

$$
T=\left(A-\lambda^{*}\right)^{-1}
$$

Since $T$ is compact, $A$ is a self-adjoint operator with compact resolvent. We also see that $A$ is bounded below. The eigenvalues of $A$ are exactly the numbers $\lambda_{n}, n=0,1,2, \ldots$

## Examples

Example 3: The Lamé-Wangerin functions are eigenfunctions of the differential equation

$$
-y^{\prime \prime}+\frac{\nu(\nu+1)}{\operatorname{sn}^{2}(x, k)} y=\lambda y
$$

Then $q_{0}=\nu(\nu+1), q_{1}=0, q_{2}=\nu(\nu+1) \frac{1}{3}\left(k^{2}+1\right)$. The boundary condition is $y(K)=0$ or $y^{\prime}(K)=0$.

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Bi, L.; Mukherjea, A. Poisson distributions: identification of parameters from the distribution of the maximum and a conjecture on the partial sums of the power series for $\exp (\mathrm{x})$. Statistics \& Probability Letters. 81 (2011), no. 5, 611-613.

Bi, L.; Mukherjea, A. Identification of parameters and the distribution of the minimum of the tri-variate normal. Statistics \& Probability Letters. 80 (2010), no. 23-24, 1819-826.

Interior and Exterior Flat-Ring Harmonics (in preparation)

## Selected Honors and Distinctions

Graduate Student Travel Award, 2017
$\$ 1,000$ to attend mathematics joint meetings in January 2018 at San Diego.

Research Excellence Award, 2016
$\$ 1,000$ per year for graduate student who exhibits outstanding research potential.

## Ernst Schwandt Teaching Award, 2016

$\$ 250$ award for Graduate Teaching Assistants who demonstrate outstanding teaching performance.
Only 4 Graduate Teaching Assistants out of 90 won this award in 2016.

## Chancellor's Fellowship Award, 2010-2015

$\$ 4,000$ per year for years 2010-2012, $\$ 3,000$ for 2013, $\$ 2,000$ for 2014, and $\$ 1,000$ for 2015 for
graduate students with exceptional academic records with high promise of future success.

