ANALYSIS OF FLUID FLOW IN THE ENTRANCE REGION OF A DUCT WITH AN ECCENTRIC ANNULAR CROSS-SECTION
A Dissertation
Presented to
the Faculty of the Department of Mathematics
University of Houston
Houston, Texas

# In Partial Fulfillment <br> of the Requirements for the Degree Doctor of Philosophy 

by
John Terry Wilson
December 1974

# ANALYSIS OF FLUID FLOW IN THE ENTRANCE REGION OF A DUCT WITH AN ECCENTRIC <br> ANNULAR CROSS-SECTION 

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#### Abstract

The purpose of this dissertation is to analyze the velocity development of laminar incompressible flow in the entrance region of a straight duct with an unchanging eccentric annular cross-section. A linearized version of the governing equations is solved under the assumptions that the velocity is zero on the duct wall and that the initial velocity profile is uniform across the crosssection. The analysis leads to a two-dimensional eigenvalue value problem which is then posed in the appropriate Hilbert space. Galerkin's method is shown to converge for, and then applied to, this eigenvalue problem.


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## CHAPTER 1

## INTRODUCTION

In a straight duct with an arbitrary but unchanging cross-section, the laminar flow of an incompressible fluid undergoes a velocity development from an initial profile at the entrance of the duct to a fully developed profile far downstream. The region where this development occurs is called the entrance region of the duct. Analytical expressions for the fully developed velocity have been obtained for ducts of various geometries. However, due to the nonlinearity of the governing equations, an analytical expression for the velocity in the entrance region has never been obtained even for the simplest geometries.

In 1964, Sparrow et. al. [1]* proposed a linearized version of the momentum equation and then proceeded, after making several assumptions, to solve the entrance region problem for two ducts: the parallel-plate channel and the circular tube. Since their solutions compare favorably with experimental results, their linearization procedure has been applied to ducts of other shapes. In particular,

[^0]Wiginton and Wendt [2] have generalized this procedure and have obtained results for the rectangular duct. A report by Shah and London [3] contains a summary of research in this and other areas.

Using the Sparrow linearization procedure, we propose to analyze the velocity development in the entrance region of a duct with an eccentric annular cross-section. This analysis leads to a two dimensional boundary value problem which we solve by the method of Galerkin. Since the Sparrow procedure is applicable to ducts of various shapes, we now state the problem and present part of the analysis for a duct with an arbitrary cross-section.

We consider steady state laminar incompressible flow in a straight duct with an arbitrary but unchanging cross-section. The duct axis is taken to be in the direction of the positive $z$-axis and the duct entrance is assumed to lie in the xy-plane (see Fig. 1)*. The equations which govern the fluid motion in the duct are the $z$-direction momentum equation and the continuity equation. For this type of flow, they simplify to

$$
\begin{equation*}
\bar{V} \cdot \operatorname{grad} w=-\frac{1}{\rho} \frac{\partial p}{\partial z}+\nu \nabla^{2} w \tag{1-1}
\end{equation*}
$$

and
*The figures for Chapters 1 and 2 are contained in Appendix A.

$$
\begin{equation*}
\operatorname{div} \overline{\mathrm{V}}=0 \tag{1-2}
\end{equation*}
$$

respectively (see Appendix B). In these equations, $\overline{\mathrm{V}}=(\mathrm{u}, \mathrm{v}, \mathrm{w})$ is the velocity vector, $\mathrm{p}(\mathrm{x}, \mathrm{y}, \mathrm{z})$ is the pressure, $\rho$ is the density, and $\nu$ is the kinematic viscosity. The symbol $\nabla^{2}$ represents the three dimensional Laplacian operator in $\mathrm{x}, \mathrm{y}$, and z . We now make the following assumptions:
(i) the term $\frac{\partial^{2} w}{\partial z^{2}}$ is negligible compared to

$$
\frac{\partial^{2} w}{\partial x^{2}}+\frac{\partial^{2} w}{\partial y^{2}}
$$

(ii) the static pressure $p(x, y, z)=p(z)$ has z dependence only.

Equation (1-1) becomes

$$
\begin{equation*}
\bar{v} \cdot \operatorname{grad} w=-\frac{1}{\rho} \frac{d p}{d z}+v \nabla^{2} w \tag{1-3}
\end{equation*}
$$

where the symbol $\nabla^{2}$ now represents the two dimensional Laplacian operator in $x$ and $y$. At this point of the analysis, we replace equation (1-3) with the Sparrow linearized version:

$$
\begin{equation*}
\varepsilon(z) \mathrm{W} \frac{\partial \mathrm{~W}}{\partial \mathrm{z}}=\Omega(\mathrm{z})+\nu \nabla^{2} \mathrm{~W} \tag{1-4}
\end{equation*}
$$

where $W$ is the average velocity over the cross-section A, $\varepsilon(z)$ is a to-be-determined function which weights the average velocity W , and $\Omega(z)$ is another to-be-determined function which includes the pressure term and the residual of the inertia terms. Integration of equation (1-4) over the cross-section A yields

$$
\Omega(z)=-\frac{\nu}{\mathrm{A}} \int_{\mathrm{A}} \nabla^{2} \mathrm{w} \mathrm{dA} .
$$

Thus by Green's theorem,

$$
\begin{equation*}
\Omega(z)=-\frac{\nu}{A} \oint_{C} \frac{\partial W}{\partial N} d S \tag{1-5}
\end{equation*}
$$

where C represents the boundary of the cross-section and $N$ represents the outward normal. Since $\varepsilon(z)$ is unknown, we temporarily suppress it by introducing a stretched axial coordinate $z^{*}$ defined by

$$
\begin{equation*}
\mathrm{d} z=\varepsilon(z) \mathrm{d}^{*} \tag{1-6}
\end{equation*}
$$

Using equations (1-5) and (1-6), equation (1-4) becomes

$$
\begin{equation*}
W \frac{\partial w}{\partial z^{*}}+\frac{v}{A} \phi_{C} \frac{\partial w}{\partial N} d S=\nu \nabla^{2} w . \tag{1-7}
\end{equation*}
$$

The problem undertaken in this dissertation is to solve equation (1-7) for $w(x, y, z)$, the $z$ component of the velocity vector, under the following conditions:
(i) the function $w(x, y, z)=0$ on the duct wall.
(ii) the initial profile at the duct entrance is uniform across the cross-section, i.e., $w(x, y, 0)=W$.

We now introduce the following dimensionless variables:
(i) $\phi=\frac{W}{W}$

$$
\text { (ii) } \xi=\frac{\mathrm{x}}{\sqrt{\mathrm{~A}}}
$$

$$
\text { (iii) } n=\frac{y}{\sqrt{A}}
$$

$$
\text { (iv) } B=\frac{z^{*}}{\sqrt{A} R}
$$

$$
\text { (v) } \quad \mathrm{n}=\frac{\mathrm{N}}{\sqrt{\mathrm{~A}}}
$$

$$
\text { (vi) } \quad s=\frac{\mathrm{S}}{\sqrt{\mathrm{~A}}}
$$

$$
\text { (vii) } R=\frac{W \sqrt{A}}{v}
$$

Substitution of these variables into equation (1-7) yields

$$
\begin{equation*}
\frac{\partial \phi}{\partial \beta}+\phi_{\overline{\mathrm{C}}} \frac{\partial \phi}{\partial \mathrm{n}} \mathrm{ds}=\nabla^{2}{ }_{\phi} \tag{1-8}
\end{equation*}
$$

where the symbol $\nabla^{2}$ represents the two dimensional

Laplacian operator in $\xi$ and $\eta$. We seek a solution to equation (1-8) of the form

$$
\begin{equation*}
\phi(\xi, \eta, \beta)=\phi_{e}(\xi, \eta, \beta)+\phi_{f}(\xi, \eta) \tag{1-9}
\end{equation*}
$$

where $\phi_{f}(\xi, \eta)$ is the fully developed velocity and $\phi_{e}(\xi, \eta, \beta)$ is a difference velocity. Substitution of equation (1-9) into equation (1-8) yields

$$
\begin{align*}
& {\left[\nabla^{2} \phi_{e}-\frac{\partial \phi_{e}}{\partial \beta}-\phi_{\bar{C}} \frac{\partial \phi_{e}}{\partial n} d s\right]} \\
& \quad+\left[\nabla^{2} \phi_{f}-\oint_{\bar{C}} \frac{\partial \phi_{f}}{\partial n} d s\right]=0 . \tag{1-10}
\end{align*}
$$

The fully developed velocity $\phi_{f}(\xi, n)$ is a solution of

$$
\begin{equation*}
\nabla^{2}{ }_{\phi_{f}}(\xi, n)=\frac{d p}{d z} \tag{1-11}
\end{equation*}
$$

where

$$
\begin{equation*}
\phi_{f}(\xi, \eta)=0 \quad \text { on } \quad \overline{\mathrm{C}} \tag{1-12}
\end{equation*}
$$

The pressure drop $\frac{d p}{d z}$ is assumed $z$-independent. Using Green's theorem and equation (1-11) we obtain

$$
\begin{aligned}
\oint_{\bar{C}} \frac{\partial \phi_{f}}{\partial \mathrm{n}} \mathrm{ds} & =\int_{\bar{A}} \nabla_{\mathrm{f}}^{2} \phi_{\mathrm{A}} \overline{\mathrm{~A}} \\
& =\int_{\bar{A}} \frac{\mathrm{dp}}{\mathrm{dz}} \mathrm{~d} \overline{\mathrm{~A}} \\
& =\frac{d p}{\mathrm{dz}}
\end{aligned}
$$

since the pressure $p(x, y, z)=p(z)$ is assumed to be a function of $z$ only. This shows that the second term in equation (1-10) is zero, i.e.,

$$
\begin{equation*}
\nabla^{2} \phi_{f}-\oint_{\bar{C}} \frac{\partial \phi_{f}}{\partial n} d s=0 \tag{1-13}
\end{equation*}
$$

Equation (1-10) becomes

$$
\begin{equation*}
\nabla^{2} \phi_{e}-\frac{\partial \phi_{e}}{\partial \beta}-\oint_{\bar{C}} \frac{\partial \phi_{e}}{\partial \mathrm{n}} \mathrm{ds}=0 \tag{1-14}
\end{equation*}
$$

Clearly, the difference velocity $\phi_{e}(\xi, \eta, \beta)$ is of significance only in the entrance region, i.e.,

$$
\begin{equation*}
\operatorname{Lim}_{\beta \rightarrow \infty} \phi_{e}(\xi, \eta, \beta)=0 \tag{1-15}
\end{equation*}
$$

With condition (1-15) in mind, we seek a solution to equation (1-14) of the form

$$
\begin{equation*}
\phi_{e}(\xi, \eta, \beta)=\sum_{j=1}^{\infty} c_{j} g_{j}(\xi, \eta) \exp \left(-\alpha_{j}^{2} \beta\right) \tag{1-16}
\end{equation*}
$$

where $c_{j}, g_{j}$, and $\alpha_{j}$ are to be determined. Substitution of equation (1-16) into equation (1-14) yields equations of the form

$$
\begin{equation*}
\nabla^{2} g_{j}+\alpha_{j}^{2} g_{j}=\oint_{\bar{C}} \frac{\partial g_{j}}{\partial n} d s \tag{1-17}
\end{equation*}
$$

for the determination of the $g_{j}$ 's. The $g_{j}$ 's are to satisfy the boundary condition

$$
\begin{equation*}
g_{j}(\xi, n)=0 \quad \text { on } \quad \bar{c} . \tag{1-18}
\end{equation*}
$$

The above analysis does not depend on a particular cross-sectional geometry. We now consider an eccentric annular cross-section (see Fig. 2). We have two boundary value problems to solve. One of them is the eigenvalue problem

$$
\begin{equation*}
-\nabla^{2} g(\xi, \eta)+\oint_{\bar{C}} \frac{\partial g}{\partial n} d s=\lambda g(\xi, \eta) \tag{1-19}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{g}(\xi, n)=0 \quad \text { on } \quad \overline{\mathrm{c}} \tag{1-20}
\end{equation*}
$$

(see Fig. 3). The other equation to solve is the fully developed velocity equation:

$$
\begin{equation*}
\nabla^{2} \phi_{f}(\xi, n)-\oint_{\bar{C}} \frac{\partial g}{\partial n} d s=0 \tag{1-21}
\end{equation*}
$$

where

$$
\begin{equation*}
\phi_{f}(\xi, \eta)=0 \quad \text { on } \quad \overline{\mathrm{C}} . \tag{1-22}
\end{equation*}
$$

We do not choose to solve equations (1-19) and (1-21) on the annulus. Instead, in Chapter 2 we transform the annulus conformally onto a rectangle and transform the eigenvalue and fully developed velocity equations accordingly.

In Chapter 3, we formulate our eigenvalue problem in a Hilbert space setting. First, we define the appropriate Hilbert spaces and an operator L. The operator L is then shown to be symmetric and positive definite. Finally, we define the energy space which is the space where Galerkin's method converges.

In Chapter 4, we find a linearly independent complete sequence for the energy space. We use linear combinations of elements of this sequence to approximate solutions of equation (2-26).

Chapter 5 contains an explanation of the method of Galerkin for eigenvalue problems and the proof of convergence for our particular operator L. The approximate eigenvalues of equation (2-26) are found by setting the Galerkin determinant equal to zero. In Chapter 6, we
obtain expressions for the inner products required in the Galerkin determinant.

In Chapter 7, we solve the transformed fully developed velocity equation. The method of solution is to expand $\phi_{f}$ in a Fourier series.

In Chapter 8, we give a procedure for finding approximations to the dimensionless velocity $\phi(\xi, \eta, \beta)$.

CHAPTER 2

THE TRANSFORMATION

In Chapter 1, we stated that we are seeking a solution $\phi_{e}$ of equation (1-14) of the form

$$
\phi_{e}(\xi, \eta, \beta)=\sum_{j=1}^{\infty} c_{j} g_{j}(\xi, \eta) \exp \left(-\alpha_{j}^{2} \beta\right)
$$

(see Fig. 2). Since our flow is laminar, our geometry is symmetric with respect to the $\xi \beta-\mathrm{plane}$, and the effect of gravity is neglected, $\phi_{e}$ is symmetric with respect to the $\xi \beta-p l a n e . ~ T h e s e ~ a s s u m p t i o n s ~ a l s o ~ i m p l y ~$ that $\frac{\partial \phi_{e}}{\partial \eta}=0$ on $\bar{B}$ (see Fig. 3). Thus we may assume that the $g_{j}(\xi, \eta)$ are symmetric with respect to the $\xi$-axis so that we need consider only the lower half of the crosssection of the duct when solving

$$
\begin{equation*}
\nabla^{2} g(\xi, n)-\oint_{\bar{C}} \frac{\partial g}{\partial n} d s=\lambda g(\xi, \eta) \tag{2-1}
\end{equation*}
$$

where

$$
\begin{equation*}
g(\xi, \eta)=0 \quad \text { on } \quad \bar{C} \tag{2-2}
\end{equation*}
$$

By Green's theorem and the symmetry of $g$, we have

$$
\begin{aligned}
\oint_{\overline{\mathrm{C}}} \frac{\partial \mathrm{~g}}{\partial \mathrm{n}} \mathrm{ds} & =\int_{\overline{\mathrm{A}}} \nabla^{2} \mathrm{~g} \mathrm{~d} \overline{\mathrm{~A}} \\
& =2 \int_{\overline{\mathrm{A}}} \overline{\mathrm{D}}^{2} \mathrm{~g} \mathrm{~d} \overline{\mathrm{~A}}^{\mathrm{L}} \\
& =2 \oint_{\overline{\mathrm{B}}, \overline{\mathrm{C}}^{\mathrm{L}}} \frac{\partial \mathrm{~g}}{\partial \mathrm{n}} \mathrm{ds}
\end{aligned}
$$

where $\overline{\mathrm{A}}^{\mathrm{L}}, \overrightarrow{\mathrm{B}}$, and $\overline{\mathrm{C}}^{\mathrm{L}}$ are as indicated in Figure 4. Thus we must solve the equation

$$
\begin{equation*}
-\nabla^{2} g(\xi, \eta)+2 \int_{-\bar{A}^{-}} \nabla^{2} g d A^{-L}=\lambda g(\xi, \eta) \tag{2-3}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathrm{g}(\xi, n)=0 \quad \text { on } \quad \overline{\mathrm{C}}^{\mathrm{L}} \tag{2-4}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial g(\xi, \eta)}{\partial \eta}=0 \quad \text { on } \quad \bar{B} . \tag{2-5}
\end{equation*}
$$

We choose not to solve equation (2-3) on the lower half of the eccentric annulus but to transform this region onto a rectangle and solve the transformed equation. We use a composition of transformations to achieve our purpose. The first transformation maps the lower half of the eccentric annulus onto the upper half of a concentric annulus with center at the origin.

The second transformation maps the upper half of the concentric annulus onto a rectangle. We now present the details of these transformations.

Considering Figure 3, we recall that nothing has been said about $c_{1}$ and $c_{2}$,i.e., about the location of the eccentric annulus on the $\xi$-axis. We locate the eccentric annulus in such a way that the two boundary circles form a coaxial system of circles with the $\xi$-axis as the line of centers and the $n$-axis as the radial axis. It is shown in [5] that such a location is possible and that the circles have the following equations:

$$
\begin{align*}
& \mathrm{C}_{1}: \xi^{2}+n^{2}+2 \frac{\mathrm{c}_{1}}{\sqrt{\mathrm{~A}}} \xi+c=0  \tag{2-6}\\
& \mathrm{c}_{2}: \xi^{2}+n^{2}+2 \frac{\mathrm{c}_{2}}{\sqrt{\mathrm{~A}}} \xi+c=0 \tag{2-7}
\end{align*}
$$

where $c_{1}, c_{2}$, and $c$ are to be determined. We have that

$$
\left(\xi+\frac{c_{1}}{\sqrt{A}}\right)^{2}+\eta^{2}=\frac{c_{1}^{2}}{A}-c=\frac{r_{1}^{2}}{A}
$$

and

$$
\left(\xi+\frac{c_{2}}{\sqrt{A}}\right)^{2}+n^{2}=\frac{c_{2}^{2}}{A}-c=\frac{r_{2}^{2}}{A} .
$$

Solving the equations

$$
\begin{aligned}
& c_{1}^{2}-c A=r_{1}^{2} \\
& c_{2}^{2}-c A=r_{2}^{2} \\
& c_{2}-c_{1}=\rho
\end{aligned}
$$

for $c_{1}, c_{2}$, and $c$ yields

$$
\begin{align*}
& c_{1}=\frac{r_{2}^{2}-r_{1}^{2}-\rho^{2}}{2 \rho}  \tag{2-8}\\
& c_{2}=\frac{r_{2}^{2}-r_{1}^{2}+\rho^{2}}{2 \rho}  \tag{2-9}\\
& c=\frac{c_{2}^{2}-r_{2}^{2}}{A}=\frac{c_{1}^{2}-r_{1}^{2}}{A} . \tag{2-10}
\end{align*}
$$

For convenience in this and later chapters, we set
(i) $K_{1}=\sqrt{c_{1}+r_{1}}$
(ii) $\mathrm{K}_{2}=\sqrt{\mathrm{c}_{1}-\mathrm{r}_{1}}$
(iii) $K_{3}=\sqrt{c_{2}+r_{2}}$
(iv) $K_{4}=\sqrt{c_{2}-r_{2}}$
(v) $K_{5}=\ln \left[\frac{K_{3}+K_{4}}{K_{3}-K_{4}}\right]$
(vi) $K_{6}=\ln \left[\frac{K_{1}+K_{2}}{K_{1}-K_{2}}\right]$
(vii) $K_{7}=K_{6}-K_{5}$
(viii) $K_{8}=\frac{K_{1}-K_{2}}{K_{1}+K_{2}}$
(ix) $K_{9}=\frac{\left(\mathrm{K}_{1}-\mathrm{K}_{2}\right)^{2}}{2 \mathrm{~K}_{1} \mathrm{~K}_{2}}$
(x) $K_{10}=\frac{\left(K_{1}^{2}-K_{2}^{2}\right)^{2}}{4 K_{1}^{2} K_{2}^{2}}$
(xi) $K_{11}=\frac{\mathrm{K}_{1}^{2}+\mathrm{K}_{2}^{2}}{2 \mathrm{~K}_{1} \mathrm{~K}_{2}}$
(xii) $K_{12}=\frac{K_{3}-K_{4}}{K_{3}+K_{4}}$
(xiii) $K_{13}=\frac{\left(K_{3}-K_{4}\right)^{2}}{2 K_{3} K_{4}}$
(xiv) $K_{14}=\frac{\left(\mathrm{K}_{3}^{2}-\mathrm{K}_{4}^{2}\right)^{2}}{4 \mathrm{~K}_{3}^{2} \mathrm{~K}_{4}^{2}}$
(xv) $K_{15}=\frac{K_{3}^{2}+K_{4}^{2}}{2 K_{3} K_{4}}$.

We now show that the bilinear transformation

$$
\begin{equation*}
W_{1}(z)=\xi_{1}+i \eta_{1}=\frac{z+p}{z-p} \tag{2-11}
\end{equation*}
$$

where $z=\xi+i n$ and $p=\sqrt{c}$ maps $C_{1}$ and $C_{2}(\xi \eta-p l a n e)$ onto circles ( $\xi_{1} \eta_{1}$-plane) whose centers are at the origin.

Solving equation (2-11) for $\xi$ and $\eta$ yields

$$
\xi=\frac{p\left[\left(\xi_{1}^{2}-1\right)+\eta_{1}^{2}\right]}{\left(\xi_{1}-1\right)^{2}+\eta_{1}^{2}}
$$

and

$$
\begin{equation*}
n=\frac{-2 p n_{1}}{\left(\xi_{1}-1\right)^{2}+n_{1}^{2}} \tag{2-13}
\end{equation*}
$$

Substitution of equations (2-12) and (2-13) into equations (2-6) and (2-7) yields

$$
\begin{align*}
& C_{1}^{\prime}\left(\xi_{1} \eta_{1}-p \text { lane }\right): \quad \xi_{1}^{2}+\eta_{1}^{2}=\frac{c_{1}-\sqrt{c_{1}-r_{1}^{2}}}{c_{1}+\sqrt{c_{2}^{2}+r_{1}^{2}}}  \tag{2-14}\\
& C_{2}^{\prime}\left(\xi_{1} \eta_{1}-\text { plane }\right): \quad \xi_{1}^{2}+\eta_{1}^{2}=\frac{c_{2}-\sqrt{c_{2}^{2}-r_{2}^{2}}}{c_{2}+\sqrt{c_{2}^{2}+r_{2}^{2}}} . \tag{2-15}
\end{align*}
$$

Figure 5 shows the location of $C_{1}$ and $C_{2}$ in the $\xi_{1} \eta_{1}-$ plane ( $C_{1}^{\prime}$ and $C_{2}^{\prime}$, respectively). Thus $W_{1}(z)$ maps the
circles $C_{1}$ and $C_{2}$ onto circles $C_{1}^{\prime}$ and $C_{2}^{\prime}$ with centers at the origin. That $W_{1}(z)$ maps the lower half of the eccentric annulus onto the upper half of the concentric annulus is seen from the equations

$$
\begin{align*}
& \xi_{1}=\frac{\xi^{2}+\eta^{2}-p^{2}}{(\xi-p)^{2}+\eta^{2}}  \tag{2-16}\\
& \eta_{1}=\frac{-2 p \eta}{(\xi-p)^{2}+\eta^{2}} \tag{2-17}
\end{align*}
$$

which are found by solving equation (2-11) for $\xi_{1}$ and $\eta_{1}$. We now apply the transformation
$W_{2}\left(z_{1}\right)=\xi_{2}+i \eta_{2}=\log W_{1}=\log _{e}\left|W_{1}\right|+i \arg \left(W_{1}\right)$.

Since the Log transformation maps circles onto straight lines, the upper half of Figure 5 is transformed onto the rectangle shown in Figure 6. Hence the transformation

$$
\begin{equation*}
W_{2}(z)=\log \frac{z+p}{z-p} \tag{2-19}
\end{equation*}
$$

maps the region shown in Figure 4 onto the region shown in Figure 6.

The Jacobian $J$ of the transformation $W_{2}(z)$ is given by

$$
J\left(\xi_{2}, \eta_{2}\right)=\left(\xi_{\xi_{2}}{ }^{\eta} \eta_{2}-\xi_{\eta_{2}} \eta_{2}\right)^{-1}
$$

where,e.g.,$\xi_{\xi_{2}}$ represents the derivative of $\xi$ with respect to $\xi_{2}$ [6]. A straight forward, but tedious calculation shows that

$$
\begin{equation*}
J\left(\xi_{2}, \eta_{2}\right)=\frac{A}{K_{1}^{2} K_{2}^{2}}\left(\cosh \xi_{2}-\cos \eta_{2}\right)^{2} \tag{2-20}
\end{equation*}
$$

Some properties of the Jacobian are presented in Appendix C. In terms of the coordinates $\xi_{2}$ and $\eta_{2}$, equation (2-3) is given by

$$
\begin{align*}
-J\left(\xi_{2}, \eta_{2}\right) & \nabla^{2} g\left(\xi_{2}, \eta_{2}\right) \\
& +2 \int_{\bar{A}_{2}} \nabla^{2} g \bar{d}_{2}=\lambda g\left(\xi_{2}, \eta_{2}\right) \tag{2-21}
\end{align*}
$$

Conditions (2-4) and (2-5) are given by

$$
\begin{equation*}
\mathrm{g}\left(\xi_{2}, \eta_{2}\right)=0 \quad \text { on } \quad \overline{\mathrm{C}}_{2} \tag{2-22}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial g\left(\xi_{2}, \eta_{2}\right)}{\partial n_{2}}=0 \quad \text { on } \quad \bar{B}_{2} \tag{2-23}
\end{equation*}
$$

respectively (see Fig. 6).

$$
\begin{align*}
& \xi_{3}=\xi_{2}-K_{5}  \tag{2-24}\\
& \eta_{3}=\eta_{2} . \tag{2-25}
\end{align*}
$$

The transformed region is shown in Figure 7. The transformed equation is

$$
\begin{align*}
& -J\left(\xi_{3}, \eta_{3}\right) \nabla^{2} g\left(\xi_{3}, \eta_{3}\right) \\
&  \tag{2-26}\\
& \quad+2 \int_{\bar{A}_{3}} \nabla^{2} g \bar{A}_{3}=\lambda g\left(\xi_{3}, \eta_{3}\right)
\end{align*}
$$

where

$$
\begin{equation*}
\mathrm{g}\left(\xi_{3}, \eta_{3}\right)=0 \quad \text { on } \quad \overline{\mathrm{C}}_{3} \tag{2-27}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial g\left(\xi_{3}, \eta_{3}\right)}{\partial \eta_{3}}=0 \quad \text { on } \quad \overline{\mathrm{B}}_{3} \tag{2-28}
\end{equation*}
$$

We note that

$$
\begin{equation*}
\int_{\bar{A}_{3}} \frac{1}{J\left(\xi_{3}, n_{3}\right)} d \bar{A}_{3}=\vec{A}^{L}=\frac{1}{2} \tag{2-29}
\end{equation*}
$$

In terms of the transformed variables $\xi_{3}$ and $\eta_{3}$, the
equation for the fully developed velocity is given by

$$
\begin{equation*}
-J\left(\xi_{3}, \eta_{3}\right) \nabla^{2} \phi_{f}\left(\xi_{3}, \eta_{3}\right)+2 \int_{\bar{A}_{3}} \nabla^{2} \phi_{f} d \bar{A}_{3}=0 \tag{2-30}
\end{equation*}
$$

where

$$
\begin{equation*}
\phi_{f}\left(\xi_{3}, \eta_{3}\right)=0 \quad \text { on } \quad \bar{c}_{3} \tag{2-31}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial \phi_{f}\left(\xi_{3}, \eta_{3}\right)}{\partial \eta_{3}}=0 \quad \text { on } \quad \bar{B}_{3} \tag{2-32}
\end{equation*}
$$

For convenience in later chapters, we make the following notation changes:
(i) $\xi_{3} \rightarrow x$

$$
\begin{aligned}
& \text { (ii) } \eta_{3} \rightarrow y \\
& \text { (iii) } \bar{A}_{3} \rightarrow \bar{A}
\end{aligned}
$$

$$
\text { (iv) } \bar{B}_{3} \rightarrow B
$$

$$
\text { (v) } \overline{\mathrm{c}}_{3} \rightarrow \mathrm{C}
$$

## CHAPTER 3

A HILBERT SPACE SETTING FOR $L(\mathrm{~g})=\lambda \mathrm{g}$

Let us formulate our eigenvalue problem in a Hilbert space setting. Consider the Hilbert space $H_{1}=\left\{L_{2}(\bar{A}) ;(),\right\}$ of all square integrable real-valued functions defined on $\overline{\mathrm{A}}$. The usual inner product for $\mathrm{L}_{2}(\overline{\mathrm{~A}})^{*}$ is defined by

$$
(g, h)=\int_{\bar{A}} g h d \bar{A}
$$

but for our purposes it is advantageous to consider the inner product

$$
(g, h)_{J}=\int_{\bar{A}} \frac{g h}{J} d \bar{A}
$$

where $J=J(x, y)=\frac{A}{K_{1}^{2} K_{2}^{2}}\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}$. We set

$$
\|g\|=\sqrt{(\mathrm{g}, \mathrm{~g})} \text { and }\|\mathrm{g}\|_{\mathrm{J}}=\sqrt{(\mathrm{g}, \mathrm{~g})} \mathrm{J} .
$$

LEMMA 3-1: Let $q(x, y)$ be an element of $L_{2}(\bar{A})$ and $p(x, y)$ be a continuous function defined on $\bar{A}$. Then $p q$ is an element of $L_{2}(\overline{\mathrm{~A}})$.

[^1]PROOF:
Since $p(x, y)$ is continuous on $\bar{A}$, there exists a real number $M$ such that $|p(x, y)|^{2}<M$. Thus $\int_{\bar{A}}(p q)^{2} d \bar{A}$, which is less than $M \int_{\bar{A}} q^{2} d \bar{A}$, is finite since $\int_{\bar{A}} q^{2} d \bar{A}$ is finite.
THEOREM 3-1: $\mathrm{H}_{2}=\left\{\mathrm{L}_{2}(\overline{\mathrm{~A}}) ;(,)_{\mathrm{J}}\right\}$ is a Hilbert space.
PROOF:
Let $g$ and $h$ be elements of $H_{2}$. By Lemma $3-1, \frac{g}{\sqrt{J}}$ and $\frac{h}{\sqrt{J}}$ are also in $\mathrm{H}_{2}$. This implies that

$$
(g, h)_{J}=\int_{-} \frac{g h}{J} d \bar{A}=\int_{\bar{A}} \frac{g}{\sqrt{J}} \frac{h}{\sqrt{J}} d \bar{A}
$$

is finite,i.e., $(\mathrm{g}, \mathrm{h})_{\mathrm{J}}$ is defined. It is easy to see that
(1) $\quad(g, h)_{J}=(h, g)_{J}$
(2) $(g+h, k)_{J}=(g, k)_{J}+(h, k)_{J}$
(3) $(\alpha g, h)_{J}=\alpha(g, h)_{J}$
where $g, h$, and $k$ are in $H_{2}$ and $\alpha$ is a real number. Also,

$$
(g, g)_{J}=\int_{\bar{A}} \frac{g^{2}}{J} d \bar{A}>0
$$

if $g \neq 0$ since $J(x, y)$ is positive on $\bar{A}$. Thus (, ) $J$ defines an inner product. We now prove completeness. Let $\left\{g_{k}\right\}_{k=1}^{\infty}$ be a Cauchy sequence in $H_{2}$. This implies that

$$
\left\|\frac{g_{m}}{\sqrt{J}}-\frac{g_{n}}{\sqrt{J}}\right\|^{2}=\int_{\bar{A}} \frac{\left(g_{m}-g_{n}\right)^{2}}{J} d \bar{A}=\left\|g_{m}-g_{n}\right\|_{J}
$$

converges to zero as $m$ and $n$ approach infinity. Since
$\left\{\frac{\mathrm{g}_{\mathrm{k}}}{\sqrt{\mathrm{J}}}\right\}_{\mathrm{k}=1}^{\infty}$ is Cauchy is $\mathrm{H}_{1}$, there exists a $\mathrm{g}^{\prime}$ in $\mathrm{H}_{1}$ so that

$$
\left\|g_{n}-g\right\|_{J}=\left\|\frac{g_{n}}{\sqrt{J}}-\frac{g}{\sqrt{J}}\right\|=\left\|\frac{g_{n}}{\sqrt{J}}-g^{\prime}\right\|
$$

converges to zero where $g=\sqrt{J} g^{\prime}$. Since $g$ is an element of $\mathrm{H}_{2}, \mathrm{H}_{2}$ is complete and hence $\mathrm{H}_{2}$ is a Hilbert space.

THEOREM 3-2: The mapping $T: H_{1} \rightarrow H_{2}$ defined by

$$
\mathrm{T}(\mathrm{~g})=\sqrt{\mathrm{J}} \mathrm{~g}
$$

is a Hilbert space isomorphism.

PROOF:
Clearly T defines an isomorphism between $\mathrm{H}_{1}$ and $\mathrm{H}_{2}$. Also,

$$
(T(g), T(h))_{J}=\int_{\bar{A}} \frac{(\sqrt{J} g)(\sqrt{J h})}{J} d \bar{A}=(g, h) .
$$

Thus $T$ defines a Hilbert space isomorphism.

Define a subspace $M$ of $H_{2}$ as follows: Let $g$ be an element of $M$ if and only if
(i) $g$, the first partial derivatives of $g$, and the
second partial derivatives of $g$ are continuous in the closure of $\bar{A}$.
(ii) $g=0$ on $C$ and $\frac{\partial g}{\partial y}=0$ on $B$.

Now define an operator $L: M \rightarrow H_{2}$ by

$$
L(g)=-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}
$$

Note that $\mathrm{L}(\mathrm{g})$ is continuous in the closed region. Our problem then is to solve the eigenvalue problem

$$
\begin{equation*}
L(g)=\lambda g \tag{3-1}
\end{equation*}
$$

where $g$ is an element of $M$ and $M$ is contained in $H_{2}$.
THEOREM 3-3: If $g$ is an element of $M$, then $(1, L(g))_{J}=0$.
PROOF:
For g in M , we have that

$$
\begin{aligned}
(1, L(g))_{J} & =\int_{\bar{A}} \frac{1}{J}\left(-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}\right) d \bar{A} \\
& =-\int_{A}^{A} \nabla^{2} g d \bar{A}+2 \int_{\bar{A}} \nabla^{2} g d \bar{A} \int_{\bar{A}} \frac{1}{J} d \bar{A} \\
& =0
\end{aligned}
$$

since $\int_{\bar{A}} \frac{1}{J} d \bar{A}=\frac{1}{2}$.

COROLLARY 1: If $\lambda, g$ is a solution of (3-1) for which
$\lambda \neq 0$, then $(1, g)_{J}=0$.
PROOF:
We have that $\lambda(1, g)_{J}=(1, \lambda g)_{J}=(1, L(g))_{J}=0$.
THEOREM 3-4: Let $\lambda, g$ and $\mu, h$ be solutions of (3-1) for which $\lambda \neq \mu$. Then

$$
(\mathrm{g}, \mathrm{~h})_{\mathrm{J}}=0 .
$$

PROOF:
We have that

$$
\begin{aligned}
\lambda(g, h)_{J} & =(\lambda g, h)_{J} \\
& =(L(g), h)_{J} \\
& =\int_{\bar{A}}^{\left(-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}\right) h} \begin{array}{l}
\mathrm{J} \\
\\
\end{array}=-\int_{\bar{A}} h \nabla^{2} g d \bar{A}+2 \int_{\bar{A}} \nabla^{2} g d \bar{A} \int_{\bar{A}} \frac{h}{J} d \bar{A} \\
& =-\int_{\bar{A}} h \nabla^{2} g d \bar{A}
\end{aligned}
$$

since $(1, \mathrm{~h})_{\mathrm{J}}=0$. Similarly,

$$
\mu(h, g)_{J}=-\int_{-}^{A} g \nabla^{2} h d \bar{A} .
$$

$$
\begin{aligned}
(\lambda-\mu)(g, h)_{J} & =\int_{\bar{A}}\left(g \nabla^{2} h-h \nabla^{2} g\right) d \bar{A} \\
& =\oint_{S}\left(g \frac{\partial h}{\partial n}-h \frac{\partial g}{\partial n}\right) d s \\
& =0
\end{aligned}
$$

since

$$
\mathrm{g}=0=\mathrm{h} \quad \text { on } \quad \mathrm{C}
$$

and

$$
\frac{\partial g}{\partial n}=\frac{\partial g}{\partial y}=0=\frac{\partial h}{\partial y}=\frac{\partial h}{\partial n} \text { on } B .
$$

This completes the proof.
Set $H_{3}=\left\{g\right.$ in $\left.H_{2} \mid(1, g)_{J}=0\right\}$ and $N=M \bigcap H_{3}$. Then Theorem 3-3 and Corollary 1 allow us to solve

$$
\begin{equation*}
L(g)=\lambda g \tag{3-2}
\end{equation*}
$$

where $g$ is an element of $N$ and $N$ is contained in $H_{3}$. We assume from now on that the domain of L is N , ie., $\mathrm{L}: \mathrm{N} \rightarrow \mathrm{H}_{3}$. Note that $\mathrm{H}_{3}$ is a Hilbert space since it is a closed subspace of a Hilbert space.

We now show that our operator $L: N \rightarrow H_{3}$ is symmetric, positive, and positive definite. The definition of these properties assumes that the domain of the operator is dense in the space. In the next chapter we will construct a set which is complete in $H_{3}$ and whose span is in $H$; this remark
implies that $N$ is dense in $H_{3}$.
THEOREM 3-5: The operator $\mathrm{L}: \mathrm{N} \rightarrow \mathrm{H}_{3}$ is symmetric, i.e.,

$$
(\mathrm{L}(\mathrm{~g}), \mathrm{h})_{J}=(\mathrm{g}, \mathrm{~L}(\mathrm{~h}))_{\mathrm{J}}
$$

for $a l l \mathrm{~g}$ and h in N .

PROOF:
Let $g$ and $h$ be elements of $M$. The proof of the Theorem 3-4 shows that

$$
(\mathrm{L}(\mathrm{~g}), \mathrm{h})_{J}=-\int_{\bar{A}} \mathrm{~h} \nabla^{2} \mathrm{~g} d \bar{A}
$$

and

$$
(g, L(h))_{J}=-\int_{\bar{A}} g \nabla^{2} h d \bar{A} .
$$

By Green's Theorem

$$
\begin{aligned}
(L(g), h)_{J}-(g, L(h))_{J} & =\int_{\bar{A}}\left(g \nabla^{2} h-h \nabla^{2} g\right) d \bar{A} \\
& =\oint_{S}\left(g \frac{\partial h}{\partial n}-h \frac{\partial g}{\partial n}\right) d s \\
& =0 .
\end{aligned}
$$

Hence L is symmetric.

THEOREM 3-6: The operator $\mathrm{L}: \mathrm{N} \rightarrow \mathrm{H}_{3}$ is positive, ice.,

$$
(\mathrm{L}(\mathrm{~g}), \mathrm{g})_{\mathrm{J}}>0
$$

for all $g$ in $N, g \neq 0$.

PROOF:
Let $g$ be an element of $N$. By Green's Theorem

$$
\begin{aligned}
(L(g), g)_{J} & =-\int_{\bar{A}} g \nabla^{2} g \mathrm{~d} \overline{\mathrm{~A}} \\
& =\int_{\bar{A}}\left\{\left(\frac{\partial g}{\partial x}\right)^{2}+\left(\frac{\partial g}{\partial y}\right)^{2}\right\} d \bar{A}-\oint_{S} g \frac{\partial g}{\partial \mathrm{n}} \mathrm{ds} \\
& =\int_{\bar{A}}\left\{\left(\frac{\partial g}{\partial x}\right)^{2}+\left(\frac{\partial g}{\partial y}\right)^{2}\right\} d \bar{A} \\
& \geq 0
\end{aligned}
$$

since $g=0$ on $C$ and $\frac{\partial g}{\partial n}=0$ on $\quad B . \quad$ If $(L(g), g)_{J}=0$, then $\frac{\partial g}{\partial x}=0=\frac{\partial g}{\partial y}$ on $\bar{A}$. This implies that $g$ is a constant
function. Since $g=0$ on $C, g$ must be the zero function.
Therefore $(\mathrm{L}(\mathrm{g}), \mathrm{g})_{J}>0$ if $\mathrm{g} \neq 0$. The proof is now complete.
THEOREM 3-7: The operator $L: N \rightarrow \mathrm{H}_{3}$ is positive definite, i.e., there exists a positive number $\gamma$ such that

$$
(L(g), g)_{J} \geq \gamma^{2}| | g \|_{J}^{2}
$$

for all $g$ in $N$.

PROOF:
Let $g$ be an element of $N$. For every $(x, y)$ in $\bar{A}$, we have

$$
\begin{aligned}
g(x, y) & =g(x, y)-g(0, y) \\
& =\int_{0}^{x} \frac{\partial g(\xi, y)}{\partial \xi} d \xi
\end{aligned}
$$

Using the Cauchy-Schwartz inequality, we obtain

$$
\begin{aligned}
|g(x, y)|^{2} & =\left|\int_{0}^{x} \frac{\partial g}{\partial \xi} d \xi\right|^{2} \\
& \leq \int_{0}^{x} 1^{2} d \xi \int_{0}^{x}\left(\frac{\partial g}{\partial \xi}\right)^{2} d \xi \\
& \leq K_{7} \int_{0}^{K} 7\left(\frac{\partial g}{\partial \xi}\right)^{2} d \xi
\end{aligned}
$$

Hence

$$
\begin{aligned}
& ||g||_{J}^{2}=\int_{0}^{\pi} \int_{0}^{K} 7 \frac{|g(x, y)|^{2}}{J(x, y)} d x d y \\
& \leq \int_{0}^{\pi} \int_{0}^{K} 7\left\{K_{7} \int_{0}^{K_{7}}\left(\frac{\partial g}{\partial \xi}\right)^{2} d \xi\right\} \frac{1}{\operatorname{Min} J} d x d y \\
& =\frac{\mathrm{K}_{7}^{2}}{\operatorname{Min} \mathrm{~J}} \int_{0}^{\pi} \int_{0}^{\mathrm{K}_{7}}\left(\frac{\partial g}{\partial \xi}\right)^{2} \mathrm{~d} \xi \mathrm{dy} \\
& =\frac{K_{7}^{2}}{\operatorname{Min} J} \int_{0}^{\pi} \int_{0}^{K_{7}}\left(\frac{\partial g}{\partial x}\right)^{2} d x d y \\
& \leq \frac{K_{7}^{2}}{\operatorname{Min} J} \int_{\bar{A}}\left\{\left(\frac{\partial g}{\partial x}\right)^{2}+\left(\frac{\partial g}{\partial y}\right)^{2}\right\} d \bar{A} \\
& =\frac{-\mathrm{K}_{7}^{2}}{\operatorname{Min} J} \int_{\bar{A}} g \nabla^{2} g d \bar{A} \\
& =\frac{K_{7}^{2}}{\operatorname{Min} J} \quad(L(g), g)_{J} \text {. }
\end{aligned}
$$

Setting

$$
\begin{equation*}
\gamma=\frac{\sqrt{\operatorname{Min~J}}}{\mathrm{K}_{7}} \tag{3-3}
\end{equation*}
$$

we obtain

$$
(\mathrm{L}(\mathrm{~g}), \mathrm{g})_{\mathrm{J}} \geq \gamma^{2}\|\mathrm{~g}\|_{\mathrm{J}}^{2} .
$$

Therefore L is positive definite.

We have previously defined an operator $L$ on a dense subspace $N$ (the denseness will be shown in Chapter 4) of a Hilbert space $H_{3}$ and have shown that it is symmetric and positive definite. A new inner product (, ) for $N$ is defined by setting

$$
{\underset{\mathrm{E}}{\mathrm{~J}}}_{\mathrm{g}, \mathrm{~h})_{\mathrm{J}}=(\mathrm{L}(\mathrm{~g}), \mathrm{h})_{\mathrm{J}} . . . . .}
$$

The completion of $N$ with respect to (, ) will be denoted by $\mathrm{H}_{4}$ and is called the energy space associated with the operator L . It can be shown that $\mathrm{H}_{4}$ is contained in $\mathrm{H}_{3}$
[7]. It is in the energy space where convergence of Galerkin's method will be established.

## CHAPTER 4

A COMPLETE SEQUENCE FOR THE ENERGY SPACE $H_{4}$

The purpose of this chapter is to construct a complete linearly independent sequence for $H_{4}$. With this in mind, we first construct a complete sequence for $H_{2}$. It is well known that

$$
\left\{\sqrt{\frac{2}{\mathrm{~K}_{7}}} \sin \frac{\mathrm{n} \pi \mathrm{x}}{\mathrm{~K}_{7}}\right\}_{\mathrm{n}=1}^{\infty}
$$

and

$$
\left\{\sqrt{\frac{1}{\pi}}, \sqrt{\frac{2}{\pi}} \cos m y\right\}_{m=1}^{\infty}
$$

are complete orthonormal sequences for $L_{2}\left(0, K_{7}\right)$ and - $L_{2}(0, \pi)$, respectively. Hence

$$
\begin{aligned}
S_{1}= & \left\{\sqrt{\frac{2}{\pi K_{7}}} \sin \frac{n \pi x}{K_{7}}\right\}_{n=1}^{\infty} \quad U \\
& \left\{\frac{2}{\sqrt{\pi K_{7}}} \sin \frac{n \pi x}{K_{7}} \cos m y\right\}_{m, n=1}^{\infty}
\end{aligned}
$$

is a complete orthonormal sequence for $H_{1}$ [8]. Since
$H_{1}$ and $H_{2}$ are isometrically isomorphic via multiplication by $\sqrt{J(x, y)}$,

$$
\begin{aligned}
S_{2}= & \left\{\sqrt{\frac{2}{\pi K_{7}}} \sqrt{J(x, y)} \sin \frac{n \pi x}{K_{7}}\right\}_{n=1}^{\infty} \quad U \\
& \left\{\frac{2}{\sqrt{\pi K_{7}}} \sqrt{J(x, y)} \sin \frac{n \pi x}{K_{7}} \cos m y\right\}_{m, n=1}^{\infty}
\end{aligned}
$$

is a complete orthonormal sequence for $H_{2}$.
THEOREM 4-1: The sequence

$$
\begin{aligned}
S_{3}= & \left\{\sqrt{\frac{2}{\pi K_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}}\right\}_{n=1}^{\infty} U \\
& \left\{\frac{2}{\sqrt{\pi K_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}} \cos m y\right\}_{m, n=1}^{\infty}
\end{aligned}
$$

is complete in $\mathrm{H}_{2}$.
PROOF:
Let $g$ be an element of $H_{2}$. Then $g^{\prime}=\frac{g}{\sqrt{J}}$ is an element of $H_{2}$. Let $\varepsilon$ be a positive number and set $\varepsilon^{\prime}=\varepsilon^{2} \operatorname{Min}\left(\frac{1}{J}\right)$.

Denote the elements of $S_{1}$ by $s_{k}$. Since $S_{2}$ is complete in $H_{2}$, there exists an element $\sum_{k=1}^{q} c_{k} \sqrt{J} s_{k}$ in the $\operatorname{Span}\left(S_{2}\right)$ such that

$$
\left\|g^{\prime}-\sum_{k=1}^{q} c_{k} \sqrt{J} s_{k}\right\|_{J}^{2}<\varepsilon^{\prime} .
$$

We have that
so that

$$
\left[\operatorname{Min}\left(\frac{1}{J}\right)\right]\left\|g-\sum_{k=1}^{q} c_{k} J s_{k}\right\|_{J}^{2}<\varepsilon^{\prime} .
$$

This implies that

$$
\left\|g-\sum_{k=1}^{q} c_{k} J s_{k}\right\|_{J}<\varepsilon
$$

where $\sum_{k=1}^{q} c_{k} J s_{k}$ is an element of the Span $\left(S_{3}\right)$. Therefore $S_{3}$ is complete in $H_{2}$.

The next theorem shows one way of using $S_{3}$ to form a complete sequence for $\mathrm{H}_{3}$.

THEOREM 4-2: Let $h$ be an element of a Hilbert space $\{\mathrm{H} ;()$,$\} and \mathrm{C}$ the orthogonal complement of h :

$$
C=\{g \text { in } H \mid(h, g)=0\} .
$$

Assume that $S \cup T=\left\{s_{k}\right\}_{k=1}^{\infty} U\left\{t_{j}\right\}_{j=1}^{\infty}$ is linearly
independent and complete in $H$ where $S$ is contained in C and T is contained in $\mathrm{H} \backslash \mathrm{C}$. If

$$
g=\sum_{k=1}^{\infty} a_{k} s_{k}+\sum_{j=1}^{\infty} b_{j} t_{j}
$$

is an element of $C$, then $g$ is an element of the $\overline{S p a n ~(S U T ')}$ where

$$
T^{\prime}=\left\{\frac{t_{j}}{\left(h, t_{j}\right)}-\frac{t_{j+1}}{\left(h, t_{j+1}\right)}\right\}_{j=1}^{\infty} .
$$

PROOF:
Let

$$
g=\sum_{k=1}^{\infty} a_{k} s_{k}+\sum_{j=1}^{\infty} b_{j} t_{j}
$$

be an element of $C$. We first consider the special case where only a finite number of the $b_{j}$ are different from zero. Then $g$ has the form

$$
g=\sum_{k=1}^{\infty} a_{k} s_{k}+\sum_{j=1}^{q} b_{j} t_{j}
$$

and so we must find $c_{j}$ 's such that

$$
\begin{aligned}
\sum_{j=1}^{q} b_{j} t_{j} & =\sum_{j=1}^{q-1} c_{j}\left\{\frac{t_{j}}{\left(h, t_{j}\right)}-\frac{t_{j+1}}{\left(h, t_{j+1}\right)}\right\} \\
& =c_{1} \frac{t_{1}}{\left(h, t_{1}\right)}+\left(c_{2}-c_{1}\right) \frac{t_{2}}{\left(h, t_{2}\right)} \\
& +\cdots+\left(c_{q-1}-c_{q-2}\right) \frac{t_{q-1}}{\left(h, t_{q-1}\right)} \\
& -c_{q-1} \frac{t_{q}}{\left(h, t_{q}\right)} .
\end{aligned}
$$

By equating coefficients, we obtain

$$
\begin{aligned}
& b_{1}=\frac{c_{1}}{\left(h, t_{1}\right)} \\
& b_{2}=\frac{c_{2}^{-c_{1}}}{\left(h, t_{2}\right)} \\
& \vdots \\
& b_{q-1}=\frac{c_{q-1}-c_{q-2}}{\left(h, t_{q-1}\right)} \\
& b_{q}=\frac{-c_{q-1}}{\left(h, t_{q}\right)} .
\end{aligned}
$$

Solving for the $c_{j}$ 's we obtain

$$
\begin{aligned}
& c_{1}=b_{1}\left(h, t_{1}\right) \\
& c_{2}=c_{1}+b_{2}\left(h, t_{2}\right)=\sum_{j=1}^{2} b_{j}\left(h, t_{j}\right) \\
& \vdots \\
& c_{q-1}=c_{q-2}+b_{q-1}\left(h, t_{q-1}\right)=\sum_{j=1}^{q-1} b_{j}\left(h, t_{j}\right) .
\end{aligned}
$$

Also $c_{q-1}=-b_{q}\left(h, t_{q}\right)$. It is clear that we can find $c_{1}, \cdots, c_{q-2}$ but we have two equations for finding $c_{q-1}$.

It is easy to see that both equations yield the same value for $c_{q-1}$. In fact, we have that

$$
\begin{aligned}
0 & =(h, g) \\
& =\sum_{k=1}^{\infty} a_{k}\left(h, s_{k}\right)+\sum_{j=1}^{q} b_{j}\left(h, t_{j}\right) \\
& =\sum_{j=1}^{q} b_{j}\left(h, t_{j}\right)
\end{aligned}
$$

which implies that

$$
\sum_{j=1}^{q-1} b_{j}\left(h, t_{j}\right)=-b_{q}\left(h, t_{q}\right) .
$$

q
Thus $\sum_{j=1} b{ }_{j} t_{j}$ is in the set $\overline{\operatorname{Span}\left(T^{\top}\right)}$.
We now prove the general case by showing that

$$
\sum_{j=1}^{\infty} b_{j} t_{j}=\sum_{j=1}^{\infty} c_{j}\left\{\frac{t_{j}}{\left(h, t_{j}\right)}-\frac{t_{j+1}}{\left(h, t_{j+1}\right)}\right\}
$$

j
where $c_{j}=\sum_{k=1} b_{k}\left(h, t_{k}\right)$. We have that

$$
\begin{aligned}
& \left\|\sum_{j=1}^{\infty} b_{j} t_{j}-\sum_{j=1}^{\infty} c_{j}\left\{\frac{t_{j}}{\left(h, t_{j}\right)}-\frac{t_{j+1}}{\left(h, t_{i+1}\right)}\right\}\right\| \\
& \left.=\left\|\sum_{j=1}^{\infty} b_{j} t_{j}-\sum_{j=1}^{\infty}\left(\sum_{k=1}^{j} b_{k}\left(h, t_{k}\right)\right)\left\{\frac{t_{j}}{\left(h, t_{j}\right)}-\frac{t_{j+1}}{\left(h, t_{j+1}\right)}\right\}\right\| \right\rvert\, \\
& =\| \sum_{j=1}^{\infty} b_{j} t_{j}-\sum_{j=1}^{\infty}\left(\sum_{k=1}^{j} b_{k}\left(h, t_{k}\right)\right) \frac{t_{j}}{\left(h, t_{j}\right)} \\
& \quad+\sum_{j=2}^{\infty}\left(\sum_{k=1}^{j-1} b_{k}\left(h, t_{k}\right)\right) \frac{t_{j}}{\left(h, t_{j}\right)} \| \\
& =\left\|\sum_{j=2}^{\infty} b_{j} t_{j}-\sum_{j=2}^{\infty}\left(b_{j}\left(h, t_{j}\right)\right) \frac{t_{j}}{\left(h, t_{j}\right)}\right\| \\
& =0
\end{aligned}
$$

This completes the proof.

COROLLARY 1: If S UT is a complete orthonormal sequence in $H$, then $S U T$ is complete in $C$.

COROLLARY 2: Let

$$
\begin{aligned}
& S_{4}=\left\{\sqrt{\frac{2}{\pi K_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}}\right\}_{\text {even } n} U \\
& \left\{\frac{2}{\sqrt{\pi K_{7}}} J(x, y) \sin \frac{\mathrm{n} \pi x}{\mathrm{~K}_{7}} \cos m y\right\}_{m, n=1}^{\infty}
\end{aligned}
$$

and

$$
S_{5}=\left\{\frac{1}{2 K_{7}} J(x, y)\left\{n \sin \frac{n \pi x}{K_{7}}-(n+2) \sin \frac{(n+2) \pi x}{K_{7}}\right\}\right\} \text { odd } n
$$

Then $\mathrm{S}_{4} \mathrm{U} \mathrm{S}_{5}$ is complete in $\mathrm{H}_{3}$.
PROOF:
Let $g$ be an element of $H_{3}$. Since $S_{2}$ is a complete orthonormal sequence for $\mathrm{H}_{2}$,

$$
\begin{aligned}
& \frac{\mathrm{g}}{\sqrt{J}}=\sqrt{\frac{2}{\pi K_{7}}} \sum_{\mathrm{n}=1}^{\infty} a_{\mathrm{n}} \sqrt{J(x, y)} \sin \frac{\mathrm{n} \pi x}{\mathrm{~K}_{7}}+ \\
& \frac{2}{\sqrt{\pi K_{7}}} \sum_{\mathrm{m}=1}^{\infty} \sum_{\mathrm{n}=1}^{\infty} b_{m n} \sqrt{J(x, y)} \sin \frac{\mathrm{n} \pi x}{\mathrm{~K}_{7}} \cos m y
\end{aligned}
$$

Therefore $g=\sqrt{\frac{2}{\pi K_{7}}} \sum_{n=1}^{\infty} a_{n} J(x, y) \sin \frac{n \pi x}{K_{7}}+$

$$
\frac{2}{\sqrt{\pi K_{7}}} \sum_{m=1}^{\infty} \sum_{n=1}^{\infty} b_{m n} J(x, y) \sin \frac{n \pi x}{K_{7}} \cos m y
$$

Thus each element of $\mathrm{H}_{3}$ can be represented as an infinite linear combination of elements of $S_{3}$. The set $S_{3}$ is linearly independent and complete in $\mathrm{H}_{3}$. It is easy to see that each element of $\mathrm{S}_{4}$ is orthogonal to the constant function one and that

$$
\left(1, \sqrt{\frac{2}{\pi K_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}}\right)_{J}=\frac{2}{n} \sqrt{\frac{2 K_{7}}{\pi}}
$$

if n is odd. Also if n is odd,

$$
\begin{aligned}
& \frac{\sqrt{\frac{2}{\pi K_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}}}{\frac{2}{n} \sqrt{\frac{2 K_{7}}{\pi}}}-\frac{\sqrt{\frac{2}{\pi K_{7}}} J(x, y) \sin \frac{(n+2) \pi x}{K_{7}}}{\left.\frac{2}{(n+2}\right) \sqrt{\frac{2 K_{7}}{\pi}}} \\
& =\frac{1}{2 K_{7}} J(x, y)\left\{n \sin \frac{n \pi x}{K_{7}}-(n+2) \sin \frac{(n+2) \pi x}{K_{7}}\right\} .
\end{aligned}
$$

The conclusion now follows from the theorem by setting

$$
H=H_{2},(,)=(,)_{J}, h=1, C=H_{3}, S=S_{4},
$$

$$
T=S_{3} \backslash S_{4}, \text { and } T^{\prime}=S_{5}
$$

COROLLARY 3: The domain of the operator

$$
L(g)=-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}
$$

is dense in $\mathrm{H}_{3}$, ie, $\overline{\mathrm{N}}=\mathrm{H}_{3}$.

PROOF:
It is sufficient to show that $S_{4} \mathrm{US}_{5}$ is contained in $N$ since $N$ is a subspace. Clearly each element of $S_{4} U S_{5}$ satisfies condition (i) and the first part of
(ii) in the definition of $M$ (see Chapter 3). Let

$$
f(x, y)=\sqrt{\frac{2}{\pi K_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}} .
$$

Then

$$
\frac{\partial f}{\partial y}=\sqrt{\frac{2}{\pi K_{7}}} \frac{\partial J}{\partial y} \sin \frac{n \pi x}{K_{7}} .
$$

Since $\frac{\partial J}{\partial y}=0$ on $B, \frac{\partial f}{\partial y}=0$ on $B$.

Let

$$
g(x, y)=\frac{2}{\sqrt{\pi_{7}}} J(x, y) \sin \frac{n \pi x}{K_{7}} \cos m y
$$

Then

$$
\frac{\partial g}{\partial y}=\frac{2}{\sqrt{\pi K_{7}}} \sin \frac{n \pi x}{K_{7}}\left\{-m J(x, y) \sin m y+\frac{\partial J}{\partial y} \cos m y\right\}
$$

Thus $\frac{\partial g}{\partial y}=0$ on $B$. Hence each element in $S_{4} \cup S_{5}$ satisfies the last part of (ii). Also, each element of $S_{4} \mathrm{US}_{5}$ is orthogonal to the constant function equal to 1 . Therefore N is dense in $\mathrm{H}_{3}$.

We are now in a position to construct a complete set for $H_{4}$. The construction is based on the following:

THEOREM 4-3: Let $\mathrm{L}: \mathrm{N} \rightarrow \mathrm{H}_{3}$ be a symmetric positive definite operator and let $\left\{g_{n}\right\}_{n=1}^{\infty}$ be a subset of $N$.
If $\left\{\operatorname{Lg}_{n}\right\}_{n=1}^{\infty}$ is complete in $H_{3}$, then $\left\{g_{n}\right\}_{n=1}^{\infty}$ is complete
in $H_{4}$, the energy space associated with $H_{3}$ [9].
For our operator

$$
L(g)=-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}
$$

it is easy to see that

$$
L\left(S_{6}\right)=S_{4}
$$

and

$$
L\left(S_{7}\right)=S_{5}
$$

where

$$
\begin{aligned}
S_{6}= & \left\{\sqrt{\frac{2}{\pi K_{7}}}\left(\frac{K_{7}}{n_{\pi}}\right)^{2} \sin \frac{n \pi x}{K_{7}}\right\} \text { even } n \\
& \left\{\frac{2}{\sqrt{\pi K_{7}}} \frac{1}{\left(\frac{n_{\pi}}{K_{7}}\right)^{2}+m^{2}} \sin \frac{n \pi x}{K_{7}} \cos m y\right\}_{m, n=1}^{\infty}
\end{aligned}
$$

and

$$
S_{7}=\left\{\frac { 1 } { 2 K _ { 7 } } \left[n\left(\frac{K_{7}}{n \pi}\right)^{2} \sin \frac{n \pi x}{K_{7}}-\right.\right.
$$

$$
\left.\left.(n+2)\left(\frac{K_{7}}{(n+2) \pi}\right)^{2} \sin \frac{(n+2) \pi x}{K_{7}}\right]\right\} \text { odd } n
$$

Since $\mathrm{S}_{4} \mathrm{US}_{5}$ is complete in $\mathrm{H}_{3}$, Theorem 4-3 implies that $S_{6} \mathrm{US}_{7}$ is complete in $\mathrm{H}_{4}$. We state this result as THEOREM 4-4: $S_{6} \mathrm{US}_{7}$ is complete in $\mathrm{H}_{4}$.

The next theorem follows from the definition of $S_{6} \cup S_{7}$ and the linear independence of $S_{1}$.

THEOREM 4-5: $S_{6} \cup S_{7}$ is linearly independent.

## CHAPTER 5

## THE GALERKIN METHOD AND ITS CONVERGENCE

Galerkin's method is a technique for finding approximate solutions of certain types of operator equations defined on Hilbert spaces. We now present a particular formulation of this method and state conditions sufficient for convergence. Then we will show that these conditions hold for our operator equation.

Let L be a symmetric positive definite linear operator defined on a dense subspace $N$ of the Hilbert space $\left\{\mathrm{H}_{3} ;(,)_{J}\right\}$ of real-valued square integrable functions defined on a closed domain $\overline{\mathrm{A}}$ contained in $\mathrm{E}^{2}$ which are orthogonal to the constant function one. Consider the eigenvalue operator equation

$$
\begin{equation*}
L(g)=\lambda g \tag{5-1}
\end{equation*}
$$

where $g$ must satisfy certain homogeneous boundary conditions. Select a linearly independent sequence $F=\left\{f_{i}\right\}_{i=1}^{\infty}$ of functions with the following properties:
(i) each element of $F$ is in the domain of $L$.
(ii) each element of $F$ satisfies the same homogeneous
boundary conditions as does g .
(iii) $F$ is complete with respect to the energy space $\left\{\mathrm{H}_{4} ;(,)_{E}\right\}$ associated with $\mathrm{H}_{3}$.

The set F is called a set of coordinate functions for the eigenvalue problem. We call a function $h_{n j}$ of the form

$$
h_{n j}=\sum_{i=1}^{n} a_{i j}^{n} f_{i}
$$

the $n$ 'th Galerkin approximation of a solution $g_{j}$ of equation (5-1) with respect to the $n$ 'th Galerkin approximation $\lambda_{n j}$ of the $j^{\prime}$ th eigenvalue $\lambda_{j}$ if the $a_{i j}^{n}$ and $\lambda_{n j}$ are defined as follows: first, we substitute $h_{n j}$ for $g$ in equation (5-1) and require that

$$
\begin{equation*}
\left(L\left(h_{n j}\right)-\lambda h_{n j}, f_{m}\right)_{J}=0 \tag{5-2}
\end{equation*}
$$

for $m=1,2, \cdots, n$. We have that

$$
\begin{aligned}
& \left(L\left(h_{n j}\right)-\lambda h_{n j}, f_{m}\right) \\
& =\left(L\left(\sum_{i=1}^{n} a_{i j}^{n} f_{i}\right)-\lambda \sum_{i=1}^{n} a_{i j}^{n} f_{i}, f_{m}\right) \\
& =\sum_{i=1}^{n} a_{i j}^{n}\left\{\left(L\left(f_{i}\right), f_{m}\right)-\lambda\left(f_{i}, f_{m}\right)_{J}\right\} \\
& =\sum_{i=1}^{n} a_{i j}^{n}\left\{\left(f_{i}, f_{m}\right)-\lambda\left(f_{i}, f_{m}\right)\right\}
\end{aligned}
$$

Thus equation (5-2) yields the homogeneous system

$$
\begin{aligned}
& \sum_{i=1}^{n} a_{i j}^{n}\left\{\left(f_{i}, f_{m}\right){ }_{E}-\lambda\left(f_{i}, f_{m}\right)_{J}\right\}=0 \\
& (m=1,2, \cdots, n)
\end{aligned}
$$

of $n$ equations in $n$ unknowns. For it to have a unique solution it is necessary and sufficient that the determinant $D_{n}$ of the system be equal to zero, where $D_{n}$ is as shown on the next page. Setting $D_{n}=0$ yields a polynomial in $\lambda$, the roots of which we denote by $\lambda_{n j}(j=1, \cdots, n)$. Without loss of generality we assume $\lambda_{\mathrm{n} 1} \leq \lambda_{\mathrm{n} 2} \leq \cdots \leq \lambda_{\mathrm{nn}}$. Now to find the $a_{i j}^{n}$ corresponding to $\lambda_{n j}$, we simply solve system (5-3) with $\lambda$ replaced by $\lambda_{n j}$.

We now state a sequence of theorems which lead to the convergence of Galerkin's method. S.G. Mikhlin has proved the following theorems:

THEOREM 5-1: Assume that every bounded set in $\mathrm{H}_{4}$ is compact in $\mathrm{H}_{3}$. Then
(i) there exists a countable number of eigenvalues of $L(g)=\lambda g: 0<\lambda_{I} \leq \lambda_{2} \leq \cdots \lambda_{j} \leq \cdots$; $\operatorname{Lim}_{j \rightarrow \infty} \lambda_{j}=\infty$.
(ii) the sequence $\left\{g_{j}\right\}_{j=1}^{\infty}$ of eigenvectors is complete in both $\mathrm{H}_{3}$ and $\mathrm{H}_{4}$ [7].

COROLLARY 1: $L^{-1}$ is completely continuous in $H_{4}$. THEOREM 5-2: Assume that $\mathrm{L}^{-1}$ is completely continuous in $H_{4}$. Then for each $j, \lambda_{n j}$ converges to an eigenvalue $\lambda_{j}$ of equation (5-1) as $n \rightarrow \infty$. Conversely, every eigenvalue of equation (5-1) is the limit point of one of the Galerkin approximate sequences, i.e., of $\left\{\lambda_{n j}\right\}_{n=1}^{\infty}$ for some j [10].

## N.I. Pol'skii has proved the following:

THEOREM 5-3: Assume that $\mathrm{L}^{-1}$ is completely continuous in $H_{4}$. Let $\left\{\lambda_{n j}\right\}_{n=1}^{\infty}$ be a Galerkin approximate sequence of eigenvalues and $\left\{h_{n j}\right\}_{n=1}^{\infty}$ the corresponding Galerkin approximate sequence of normalized eigenvectors. Then $\left\{h_{n j}\right\}_{n=1}^{\infty}$ is compact in $H_{4}$ and every limit point of this
set is an eigenvector of equation (5-1). Conversely, every eigenvector of equation (5-1) is a limit point of one of the Galerkin approximate sequences of normalized eigenvectors [11].

COROLLARY 1: Assume that $g_{j}$ is a limit point of $\left\{h_{n j}\right\}_{n=1}^{\infty}$ so that $h_{n_{k}}$ converges to $g_{j}$ in $H_{4}$. Then $h_{n_{k} j}$
converges to $\mathrm{g}_{\mathrm{j}}$ in $\mathrm{H}_{3}$,i.e., in the mean.
We now consider the case where $L$ is the operator

$$
L(g)=-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}
$$

and where $\mathrm{N},(,)_{\mathrm{J}}, \mathrm{H}_{3}$, etc. are defined as in Chapter 3 . Since we have already shown that this operator L is symmetric and positive definite on $\mathrm{H}_{3}$ and that its domain N is dense in $H_{3}$, the previous theorems show that Galerkin's method will converge if any set bounded in $\mathrm{H}_{4}$ is compact in $H_{3}$. As an aid in proving this we introduce the Hilbert space $\mathrm{W}_{2}^{(1)}(\overline{\mathrm{A}})=\mathrm{W}$ of functions in $\mathrm{H}_{1}=\mathrm{L}_{2}(\overline{\mathrm{~A}})$ having generalized first derivatives (see Appendix $G$ for details). The usual norm for $W,\|g\|_{W}$, is not convenient for our purposes. We use an equivalent norm for $\mathrm{W},\|\mathrm{g}\|_{\mathrm{W}^{*}}$, defined by

$$
\|g\|_{W^{*}}^{2}=|B(g)|^{2}+\int_{\bar{A}}\left\{D_{x}^{* 2} g+D_{y}^{* 2} g\right\} d \vec{A}
$$

where $B(g)$ can be any bounded linear functional with respect to $\|g\|_{W}$ which does not vanish for the constant function $h(x, y)=1[2]$. We define $B(g)=\int_{C} g$ ds; note that we are using our particular $\bar{A}$ to define $B(g)$. We now show that this $B(g)$ satisfies the required conditions. We have that

$$
|B(g)|^{2} \leq\left(\int_{C}|g| \mathrm{ds}\right)^{2} .
$$

Using the Cauchy-Schwartz inequality, we obtain

$$
\begin{aligned}
|B(g)|^{2} & \leq 2 \pi \int_{C} g^{2} d s \\
& \leq 2 \pi \int_{S} g^{2} d s \\
& \leq 2 \pi k\left\{\int_{\bar{A}}\left\{g^{2}+D_{x}^{* 2} g+D_{y}^{* 2} g\right\} d \bar{A}\right\} \\
& =2 \pi k| | g \|_{W}^{2}
\end{aligned}
$$

where $k$ is some positive real number [7]. Also,

$$
B(h)=\int_{C} h d s=2 \pi \neq 0 .
$$

Thus our $\mathrm{B}(\mathrm{g})$ satisfies the required conditions.
THEOREM 5-4: For $L(g)=-J \nabla^{2} g+2 \int_{-A} \nabla^{2} g d \bar{A}, H_{4}$ is
isometrically isomorphic to a subspace of $\mathrm{W}_{2}^{(1)}(\overline{\mathrm{A}})$.
PROOF:
For $g$ in $N$, we have that

$$
\begin{align*}
\|g\|_{W^{*}}^{2} & =\|\left.\int_{C} g d s\right|^{2}+\int_{A}\left\{\left(\frac{\partial g}{\partial x}\right)^{2}+\left(\frac{\partial g}{\partial y}\right)^{2}\right\} d \bar{A} \\
& =\int_{\bar{A}}\left\{\left(\frac{\partial g}{\partial x}\right)^{2}+\left(\frac{\partial g}{\partial y}\right)^{2}\right\} d \bar{A} \\
& =(L(g), g)_{J} \\
& =\|g\|_{E}^{2} . \tag{5-4}
\end{align*}
$$

Let $g$ be an element of $\mathrm{H}_{4} \backslash \mathrm{~N}$. Then there exists a sequence $\left\{\mathrm{g}_{\mathrm{k}}\right\}_{\mathrm{k}=1}^{\infty}$ contained in N such that $\left\|\mathrm{g}_{\mathrm{k}}-\mathrm{g}\right\|_{\mathrm{E}} \rightarrow 0$
as $k \rightarrow \infty$. Since $\left\{g_{k}\right\}_{k=1}^{\infty}$ is Cauchy in $H_{4}$, inequality (5-4)
shows that it is Cauchy in $W$. Thus there exists a $\hat{g}$ in W such that $\left\|g_{k}-\hat{g}\right\|_{W^{*}} \rightarrow 0$ as $k \rightarrow \infty$. We define

P: $\mathrm{H}_{4} \rightarrow \mathrm{~W}$ by

$$
P(g)= \begin{cases}g & \text { if } g \text { is in } N \\ \hat{g} & \text { if } g \text { is in } H_{4} \backslash N\end{cases}
$$

It is not difficult to prove that $P$ is well defined and linear. To see that $P$ is one-to-one, assume that $P(g)=\hat{g}$ and $P(h)=\hat{g}$. Then there exist $\left\{g_{k}\right\}_{k=1}^{\infty}$ and $\left\{h_{k}\right\}_{k=1}^{\infty}$ contained in $N$ such that $\left\|g_{k}-g\right\|_{E} \rightarrow 0$ and $\left\|h_{k}-h\right\|_{E} \rightarrow 0$. Since $P(g-h)=0,\left\|g_{k}-h_{k}\right\|_{W} \|^{*} \rightarrow 0$.

Thus $g_{k}-h_{k}$ converges to 0 weakly in $W$ and hence in
$H_{4}$. Thus for $p$ in $N,\left(g_{k}-h_{k}, p\right)_{E} \rightarrow 0$, ie, $(g-h, p)_{E}=0$.
Since $N$ is dense in $H_{4}$ we have that $g=h$; therefore $P$ is one-to-one. Thus $P$ is an isomorphism into $W$. The proof that $P$ is an isometry is as follows. Let $P(g)=\hat{g}$ and $\left\{g_{k}\right\}_{k=1}^{\infty}$ be a sequence in $N$ which converges to g in $\mathrm{H}_{4}$. Then

The proof is now complete.

COROLLARY 1: For all g in $\mathrm{H}_{4}$

$$
\|g\|_{E}=\|g\|_{W^{*}}
$$

The importance of being able to consider $\mathrm{H}_{4}$ as a subspace of W is seen in the following theorem:

THEOREM 5-5 (Imbedding Theorem): The injection mapping from $W$ into $H_{1}=L_{2}(\overline{\mathrm{~A}})$ is completely continuous [12].

From a previous remark, the next theorem establishes convergence of Galerkin's method for our eigenvalue problem.

THEOREM 5-6: For $L(g)=-J \nabla^{2} g+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}$, any set bounded in $\mathrm{H}_{4}$ is compact in $\mathrm{H}_{3}$.

PROOF:
Let $G$ be a bounded set in $H_{4}$. By Corollary 1 of Theorem 5-4, $G$ is bounded in $W$. Hence by the Imbedding Theorem, $G$ is compact in $H_{1}$. Since $H_{1}$ and $H_{2}$ are
isometrically isomorphic, $G$ is compact in $H_{2}$ and hence in $\mathrm{H}_{3}$. This completes the proof.

## CHAPTER 6

THE GALERKIN DETERMINE NT

We now construct the Galerkin determinant for our eigenvalue problem. The set $S_{6} U S_{7}$ (see Chapter 4) is linearly independent and complete in $\mathrm{H}_{4}$. Each of its elements is in the domain of our operator

$$
L(g)=-J(x, y) \nabla^{2} g(x, y)+2 \int_{\bar{A}} \nabla^{2} g d \bar{A}
$$

and satisfies the required homogeneous boundary conditions. We choose this set to be our set of coordinate functions. For convenience, we set

$$
F=\left\{f_{i}\right\}_{i=1}^{\infty}=S_{6} \cup^{\prime} S_{7} .
$$

Since the elements of the Galerkin determinant $D_{n}$ (see Chapter 5) are of the form $\left(f_{i}, f_{j}\right)_{E}-\lambda\left(f_{i}, f_{j}\right){ }_{J}$ it is necessary to evaluate the inner products $\left(f_{i}, f_{j}\right)$ and $\left(f_{i}, f_{j}\right)$ for various combinations of elements of $\mathrm{S}_{6} \mathrm{US}_{7}$.

The possibilities for the inner product $\left(f_{i}, f_{j}\right)_{J}$ are given by
(i) $\frac{2}{\pi K_{7}}\left(\frac{K_{7}}{n \pi}\right)^{2}\left(\frac{K_{7}}{p \pi}\right)^{2} \frac{K_{1}^{2} K_{2}^{2}}{A} \int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \sin \frac{p \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)} 2^{d y d x}$
( $\mathrm{n}, \mathrm{p}$ even)
(ii) $\frac{2 \sqrt{2}}{\pi K_{7}}\left(\frac{K_{7}}{n \pi}\right)^{2} \frac{1}{\left(\frac{\mathrm{p} \pi}{K_{7}}\right)^{2}+q^{2}} \frac{\mathrm{~K}_{1}^{2} \mathrm{~K}_{2}^{2}}{A} \int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{\mathrm{n} \pi \mathrm{x}}{\mathrm{K}_{7}} \sin \frac{\mathrm{p} \pi \mathrm{x}}{\mathrm{K}_{7}} \cos \mathrm{qy}}{\left(\cosh \left(\mathrm{x}+\mathrm{K}_{5}\right)-\cos \mathrm{y}\right)^{2}} d y d x$
(n even)
(iii) $\sqrt{\frac{2}{\pi K_{7}}}\left(\frac{K_{7}}{n_{\pi}}\right)^{2} \frac{K_{1}^{2} K_{2}^{2}}{A}\left\{\frac{p}{2 K_{7}}\left(\frac{K_{7}}{p \pi}\right)^{2}\right.$

$$
\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \sin \frac{p \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x
$$

$-\frac{\mathrm{p}+2}{2 \mathrm{~K}_{7}}\left(\frac{\mathrm{~K}_{7}}{(\mathrm{p}+2) \pi}\right)^{2}$

$$
\left.\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \sin \frac{(p+2) \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x\right\}
$$

( n even, p odd)
(iv) $\frac{4}{\pi K_{7}} \frac{1}{\left(\frac{n \pi}{K_{7}}\right)^{2}+m^{2}} \frac{1}{\left(\frac{\mathrm{R}_{7}}{K_{7}}\right)^{2}+\mathrm{q}^{2}} \frac{\mathrm{~K}_{1}^{2} \mathrm{~K}_{2}^{2}}{\mathrm{~A}}$

$$
\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \cos m y \sin \frac{\mathrm{p} \pi \mathrm{x}}{\mathrm{~K}_{7}} \cos q y}{\left(\cosh \left(x+\mathrm{K}_{5}\right)-\cos y\right)^{2}} d y d x
$$

(v) $\frac{2}{\sqrt{\pi K_{7}}} \frac{1}{\left(\frac{\mathrm{n}_{2}}{\mathrm{~K}_{7}}\right)^{2}+\mathrm{m}^{2}} \frac{\mathrm{~K}_{1}^{2} \mathrm{~K}_{2}^{2}}{\mathrm{~A}}$

$$
\left\{\frac{p}{2 K_{7}}\left(\frac{K_{7}}{p \pi}\right)^{2} \int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \cos m y \sin \frac{p \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x\right.
$$

$$
\left.-\frac{(p+2)}{2 K_{7}}\left(\frac{K_{7}}{(p+2) \pi}\right)^{2} \int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \cos m y \sin \frac{(p+2) \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x\right\}
$$

(p odd)
(iv) $\frac{1}{\left(2 K_{7}\right)^{2}} \frac{\mathrm{~K}_{1} \mathrm{~K}_{2}^{2}}{\mathrm{~A}} \ln \left(\frac{\mathrm{~K}_{7}}{\mathrm{n} \pi}\right)^{2} \mathrm{p}\left(\frac{\mathrm{K}_{7}}{\mathrm{p} \pi}\right)^{2}$

$$
\begin{gathered}
\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \sin \frac{p \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x \\
-n\left(\frac{K_{7}}{n \pi}\right)^{2}(p+2)\left(\frac{K_{7}}{(p+2)}\right)^{2} \\
\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \sin \frac{(p+2) \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} 2 d y d x
\end{gathered}
$$

$$
\begin{aligned}
& -(n+2)\left(\frac{K_{7}}{(n+2) \pi}\right)^{2} p\left(\frac{K_{7}}{p \pi}\right)^{2} \\
& \int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{(n+2) \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} \sin \frac{p \pi x}{K_{7}} \\
& +(n+2) \\
& \left(\frac{K_{7}}{(n+2) \pi}\right)^{2}(p+2)\left(\frac{K_{7}}{(p+2) \pi}\right)^{2} \\
& \left.\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{(n+2) \pi x}{K_{7}} \sin \frac{(p+2) \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x\right\} \\
& (n, p \text { odd). }
\end{aligned}
$$

Thus each inner product $\left(f_{i}, f_{j}\right)_{J}$ requires the evaluation of an integral of the form

$$
I=\int_{0}^{K} 7 \int_{0}^{\pi} \frac{\sin \frac{t \pi x}{K_{7}} \cos u y \sin \frac{v \pi x}{K_{7}} \cos w y}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x
$$

where $t$ and $v$ are positive integers and $u$ and $w$ are nonnegative integers. A closed form expression for $I$ is derived in Appendix $E$.

The possibilities for the inner products $\left(f_{i}, f_{j}\right)$
are given by
(vii) $\frac{2}{\pi K_{7}}\left(\frac{\mathrm{~K}_{7}}{\mathrm{p} \pi}\right)^{2} \int_{0}^{\mathrm{K}_{7}} \int_{0}^{\pi} \sin \frac{\mathrm{n} \pi \mathrm{x}}{\mathrm{K}_{7}} \sin \frac{\mathrm{p} \pi \mathrm{x}}{\mathrm{K}_{7}} d y d x$
( $n, \mathrm{p}$ even)
(viii) $\frac{2 \sqrt{2}}{K_{7}} \frac{1}{\left(\frac{\mathrm{p} \pi}{\mathrm{K}_{7}}\right)^{2}+\mathrm{q}^{2}} \int_{0}^{\mathrm{K}_{7}} \int_{0}^{\pi} \sin \frac{\mathrm{n} \pi \mathrm{x}}{\mathrm{K}_{7}} \sin \frac{\mathrm{p} \pi \mathrm{x}}{\mathrm{K}_{7}} \cos q y d y d x$
(n even)
(ix) $\sqrt{\frac{2}{\pi K_{7}}}\left\{\frac{p}{2 K_{7}}\left(\frac{K_{7}}{p \pi}\right)^{2} \int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{n \pi x}{K_{7}} \sin \frac{p \pi x}{K_{7}} d y d x\right.$

$$
\left.-\frac{p+2}{2 K_{7}}\left(\frac{K_{7}}{(p+2) \pi}\right)^{2} \int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{n \pi x}{K_{7}} \sin \frac{(p+2) \pi x}{K_{7}} d y d x\right\}
$$

( n even, p odd)
(x) $\frac{4}{\pi K_{7}} \frac{1}{\left(\frac{p^{\pi}}{K_{7}}\right)^{2}+q^{2}} \int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{n \pi x}{K_{7}} \cos m y \sin \frac{p \pi x}{K_{7}} \cos q y d y d x$
(xi) $\frac{2}{\sqrt{\pi K_{7}}}\left\{\frac{\mathrm{p}}{2 \mathrm{~K}_{7}}\left(\frac{\mathrm{~K}_{7}}{\mathrm{p} \pi}\right)^{2} \int_{0}^{\mathrm{K}_{7}} \int_{0}^{\pi} \sin \frac{\mathrm{n} \pi \mathrm{x}}{\mathrm{K}_{7}} \cos m y \sin \frac{\mathrm{p} \pi \mathrm{x}}{\mathrm{K}_{7}} d y d x\right.$

$$
\begin{aligned}
& -\frac{(p+2)}{2 K_{7}}\left(\frac{K_{7}}{(p+2) \pi}\right)^{2} \\
& \left.\int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{n \pi x}{K_{7}} \cos m y \sin \frac{(p+2) \pi x}{K_{7}} d y d x\right\} \\
& (p \text { odd) }
\end{aligned}
$$

(xii) $\frac{1}{4 K_{7}^{2}}\left\{\mathrm{n} \mathrm{p}\left(\frac{\mathrm{K}_{7}}{\mathrm{p} \pi}\right)^{2}\right.$

$$
\begin{aligned}
& \int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{n \pi x}{K_{7}} \sin \frac{p \pi x}{K_{7}} d y d x \\
& -\mathrm{n}(\mathrm{p}+2)\left(\frac{\mathrm{K}_{7}}{(\mathrm{p}+2) \pi}\right)^{2} \\
& \int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{n \pi x}{K_{7}} \sin \frac{(p+2) \pi x}{K_{7}} d y d x \\
& -(n+2) p\left(\frac{K_{7}^{2}}{p \pi}\right)^{2} \\
& \int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{(n+2) \pi x}{K_{7}} \sin \frac{p \pi x}{K_{7}} d y d x \\
& +(n+2)(p+2)\left(\frac{K_{7}}{(p+2) \pi}\right)^{2} \\
& \left.\int_{0}^{K_{7}} \int_{0}^{\pi} \sin \frac{(n+2) \pi x}{K_{7}} \sin \frac{(p+2) \pi x}{K_{7}} d y d x \quad\right\} \\
& \text { ( } \mathrm{n}, \mathrm{p} \text { odd) . }
\end{aligned}
$$

The integrals in (viii), (ix), and (xi) are always zero. The expression in (vii) is zero unless $n=p$ in which case it is equal to $\frac{\mathrm{K}_{7}^{2}}{\pi n^{2}}$. The expression in
( x ) is zero unless $\mathrm{n}=\mathrm{p}$ and $\mathrm{m}=\mathrm{q}$ in which case it is equal to $\frac{1}{\left(\frac{n \pi}{K_{7}}\right)^{2}+m^{2}}$. The possible values for the expression in (xii) are given by the following THEOREM 6-1: If $f_{i}$ and $f_{j}$ are elements of $S_{7}$, then

$$
-\frac{\mathrm{K}_{7}}{8} \text { if } \mathrm{n}=\mathrm{p}-2
$$

$$
\left(f_{i}, f_{j}\right)_{E}=\left\{\begin{array}{l}
\frac{k_{7}}{4} \text { if } n=p \\
-\frac{k_{7}}{8} \text { if } n=p+2
\end{array}\right.
$$

0 otherwise .

We have now determined closed form expressions for all of the inner products required in the Galerkin determinant.

Once we have chosen $n$ elements from $F$ and have evaluated the appropriate inner products, we must determine the $\lambda$ 's which satisfy $D_{n}=0$. Solving $D_{n}=0$ for $\lambda$ is equivalent to solving the matrix eigenvalue problem

$$
\begin{equation*}
A_{n} x=\lambda B_{n} x \tag{6-1}
\end{equation*}
$$

where

$$
A_{n}=\left[\begin{array}{cccc}
\left(f_{1}, f_{1}\right)_{E} & \left(f_{1}, f_{2}\right)_{E} & \cdots & \left(f_{1}, f_{n}\right)_{E} \\
\left(f_{2}, f_{1}\right)_{E} & \left(f_{2}, f_{2}\right)_{E} & & \\
\cdot & & \cdot \\
\cdot & & \\
\cdot & & & \\
\left(f_{n}, f_{1}\right)_{E} & \cdot & . & \left(f_{n}, f_{n}\right)_{E}
\end{array}\right]
$$

and


If $\lambda_{n j}$ is an eigenvalue of equation (6-1) and $x_{n j}$ is the corresponding eigenvector, then

$$
x_{n j}=\left(a_{i j}^{n}, a_{2 j}^{n}, \cdots, a_{n j}^{n}\right)
$$

gives the coefficients for the approximation

$$
h_{n j}=\sum_{i=1}^{n} a_{i j}^{n} f_{i} .
$$

Since $F$ is linearly independent and the determinants of $A_{n}$ and $B_{n}$ are Gram determinants, $A_{n}$ and $B_{n}$ are nonsingular
[9]. In fact, since the operator

$$
L(g)=-J \nabla^{2} g+2 \int_{-}^{A} \nabla^{2} g d \bar{A}
$$

is positive definite, the matrix $A_{n}$ is positive definite. This implies that $A_{n}$ can be written as

$$
A_{n}=C_{n} C_{n}^{T}
$$

for some nonsingular matrix $C_{n}$. Defining $\mu=\frac{1}{\lambda}$, we
now write equation (6-1) as

$$
B_{n} x=\mu C_{n} C_{n}^{T}
$$

which we in turn write as

$$
\begin{equation*}
G_{n} y=\mu y \tag{6-2}
\end{equation*}
$$

where

$$
G_{n}=C_{n}^{-1} B_{n} C_{n}^{-T}
$$

and

$$
y=C_{n}^{T} x
$$

If $\mu_{n j}$ is an eigenvalue of equation (6-2) and $y_{n j}$ is the corresponding eigenvector, then

$$
\lambda_{n j}=\frac{1}{\mu_{n j}}
$$

is an eigenvalue of equation (6-1) and

$$
x_{n j}=C_{n}^{-T} y_{n j}
$$

is the corresponding eigenvector. Thus we have reduced the problem of finding the $\lambda_{n j}$ 's which satisfy $D_{n}=0$ and the corresponding $a_{i j}^{n}$ 's for $h_{n j}$ to the problem of finding the eigenvalues and eigenvectors of the real symmetric matrix $G_{n}$; techniques for solving this latter problem are widely known.

The number n of coordinate functions to use in a particular approximation as well as which $f_{i}$ to choose from $S_{6}$ and which $f_{i}$ to choose from $S_{7}$ is arbitrary. The next theorem shows, however, that approximations with all the $f_{i}$ in $S_{6}$ do not contribute to the solution for $\phi_{e}$, the difference velocity.

LEMMA 6-1: The coefficients $c_{j}$ for the difference velocity $\phi_{e}$ are given by

$$
\begin{equation*}
c_{j}=\frac{-\left(\phi_{f}, g_{j}\right)_{J}}{\left\|g_{j}\right\|_{J}^{2}} \tag{6-3}
\end{equation*}
$$

PROOF:
From Chapter 1 and Chapter 2, we have that

$$
\begin{aligned}
& \text { (i) } \phi_{e}(x, y, 0)=\sum_{j=1}^{\infty} c_{j} g_{j}(x, y) \\
& \text { (ii) } \phi_{0}=\phi(x, y, 0)=1 \\
& \text { (iii) } \phi_{e}(x, y, 0)=\phi_{o}-\phi_{f}(x, y) .
\end{aligned}
$$

This implies that

$$
\begin{aligned}
\int_{\bar{A}} \frac{\left(\phi_{o}-\phi_{f}\right) g_{j}}{J} & =\int_{\bar{A}} \frac{\phi_{e} g_{j}}{J} d \bar{A} \\
& =\sum_{k=1}^{\infty} c_{k} \int_{\bar{A}} \frac{g_{k} g_{j}}{J} d \bar{A} \\
& =\sum_{k=1}^{\infty} c_{k}\left(g_{k}, g_{j}\right) \\
& =c_{j}\left\|g_{j}\right\|_{J}^{2}
\end{aligned}
$$

by Theorem 3-4. We also have that

$$
\begin{aligned}
\int_{\bar{A}} \frac{\left(\phi_{o}-\phi_{f}\right) g_{j}}{J} & =\int_{\bar{A}} \frac{\phi_{0} g_{j}}{J} d \bar{A}-\int_{\bar{A}} \frac{\phi_{f} g_{j}}{J} d \bar{A} \\
& =\left(1, g_{j}\right)-\left(\phi_{f}, g_{j}\right)_{J} \\
& =-\left(\phi_{f}, g_{j}\right)_{J}
\end{aligned}
$$

by Corollary 1 of Theorem 3-3. Therefore

$$
c_{j}=\frac{-\left(\phi_{f}, g_{j}\right)_{J}}{\left\|g_{j}\right\|_{J}^{2}}
$$

THEOREM 6-2: Let

$$
g_{j}=\sum_{k=1}^{\infty} a_{k} s_{k}
$$

be a solution of $L(g)=\lambda g$ where $s_{k}$ is in $S_{6}$ for every $k$. Then the coefficient $c_{j}$ corresponding to $g_{j}$ in the series expansion of $\phi_{e}$ is zero.

PROOF:
We set $\nabla^{2} s_{k}=b_{k} s_{k}$. We have that

$$
\begin{aligned}
\lambda\left(\phi_{f}, g_{j}\right) & =\left(\phi_{f},-J \nabla^{2} g_{j}+2 \int_{A} \nabla^{2} g_{j} d \bar{A}\right)_{J} \\
& =\left(\phi_{f},-J \nabla^{2} g_{j}+2 \sum_{k=1}^{\infty} a_{k} b_{k} \int_{\bar{A}} s_{k} d \bar{A}\right)_{J} \\
& =\left(\phi_{f},-J \nabla^{2} g_{j}\right)
\end{aligned}
$$

since $\int_{\bar{A}} s_{k} d \bar{A}=0$ for every $k$. By Green's theorem,

$$
\begin{aligned}
&\left(\phi_{f},-J \nabla^{2} g_{j}\right)=-\int_{\bar{A}} \phi_{f} \nabla^{2} g_{j} d \bar{A} \\
&=-\int_{\bar{A}} g_{j} \nabla^{2} \phi_{f} d \bar{A}+\phi_{S}\left(g_{j} \frac{\partial \phi_{f}}{\partial n}-\phi_{f} \frac{\partial g_{j}}{\partial n}\right) d s \\
&=-\int_{\bar{A}} g_{j} \nabla^{2} \phi_{f} d \bar{A} \\
& \text { since } g_{j}=0=\phi_{f} \text { on } C \text { and } \frac{\partial g_{j}}{\partial n}=0=\frac{\partial \phi_{f}}{\partial n} \text { on } B .
\end{aligned}
$$

We have that

$$
\begin{aligned}
-\int_{\bar{A}} g_{j} \nabla^{2} \phi_{f} d \bar{A} & =-\int_{\bar{A}} J \nabla{ }^{2} \phi_{f} \frac{g_{j}}{J} d \bar{A} \\
& =-\int_{\bar{A}} \frac{d p}{d z} \frac{g_{j}}{J} d \bar{A} \\
& =0
\end{aligned}
$$

since $\frac{d p}{d z}$ is a constant and $\int_{\bar{A}} \frac{g_{j}}{J} d \bar{A}=\left(1, g_{j}\right) J_{J}=0$.

Therefore $\left(\phi_{f}, g_{j}\right)_{J}=0$. This completes the proof.

## CHAPTER 7

THE FULLY DEVELOPED VELOCITY

The fully developed velocity $\phi_{f}$ (in dimensionless variables) is a solution of

$$
\begin{equation*}
\nabla^{2} \phi_{f}(\xi, \eta)=\frac{d p}{d z} \tag{7-1}
\end{equation*}
$$

where

$$
\begin{equation*}
\phi_{f}(\xi, \eta)=0 \quad \text { on } \quad \bar{C} \tag{7-2}
\end{equation*}
$$

(see Fig. 2 of Appendix A). Heyda [14] has solved the dimensional form of equation (7-1). His method of solution was to write $\phi_{f}$ as the integral of the Green's function for the Laplacian operator defined on an eccentric annulus. The evaluation of this integral leads to a series solution for $\phi_{f}$. Heyda's solution is defined on the original annulus. In order to have a solution defined on the rectangle, we now solve the transformed fully developed velocity equation. We use the method of Fourier to obtain an expression for $\phi_{f}$ in terms of the coordinate functions which were used in Galerkin's method.

In Chapter 2, we saw that the transformed fully developed velocity $\phi_{f}(x, y)$ is a solution of the equation

$$
\begin{equation*}
-J(x, y) \nabla^{2} g(x, y)+2 \int_{\bar{A}} \nabla^{2} g(x, y) d \bar{A}=0 \tag{7-3}
\end{equation*}
$$

where

$$
\begin{equation*}
g(x, y)=0 \quad \text { on } \quad C \tag{7-4}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial g(x, y)}{\partial y}=0 \quad \text { on } \quad B \tag{7-5}
\end{equation*}
$$

Since $2 \int_{\bar{A}} \nabla^{2} \mathrm{~g}(\mathrm{x}, \mathrm{y}) \mathrm{d} \overline{\mathrm{A}}$ represents the pressure drop,
we can assume that $\int_{\bar{A}} \nabla^{2} \mathrm{~g} d \overline{\mathrm{~A}} \neq 0$. Under this assumption equation (7-3) can be written as

$$
\begin{equation*}
\nabla^{2} \bar{g}=\frac{2}{J(x, y)} \tag{7-6}
\end{equation*}
$$

where

$$
\bar{g}=\frac{g}{\int_{\bar{A}} \nabla^{2} g d \bar{A}}
$$

The boundary conditions become

$$
\begin{equation*}
\bar{g}(x, y)=0 \quad \text { on } \quad C \tag{7-7}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\partial \bar{g}(x, y)}{\partial y}=0 \quad \text { on } \quad B \tag{7-8}
\end{equation*}
$$

We consider the Hilbert space $H_{1}=\left\{L_{2}(\overline{\mathrm{~A}}) ;(),\right\}$
which was defined in Chapter 3. Since a complete orthonormal sequence for $H_{1}$ is given by

$$
S_{1}=\left\{\sqrt{\frac{2}{\pi K_{7}}} \sin \frac{\mathrm{n} \pi x^{K_{7}}}{\}_{7}} \bigcup_{\mathrm{n}=1}^{\infty}\left\{\frac{2}{\sqrt{\pi \mathrm{~K}_{7}}} \sin \frac{\mathrm{n} \pi \mathrm{x}}{\mathrm{~K}_{7}} \cos m \mathrm{y}\right\}_{\mathrm{m}, \mathrm{n}=1}^{\infty},\right.
$$

the function $\bar{g}(x, y)$ can be represented in the form

$$
\begin{align*}
\bar{g}(x, y)= & \sum_{n=1}^{\infty} a_{n} \sin \frac{n \pi x}{K_{7}} \\
& +\sum_{n=1}^{\infty} \sum_{m=1}^{\infty} b_{m n} \sin \frac{n \pi x}{k_{7}} \cos m y \tag{7-9}
\end{align*}
$$

Since each element of $S_{1}$ satisfies conditions (7-7) and (7-8), so does $\bar{g}(x, y)$. Substituting equation (7-9) into equation (7-6) yields
$-\sum_{n=1}^{\infty}\left(\frac{n \pi}{K_{7}}\right)^{2} a_{n} \sin \frac{n \pi x}{K_{7}}$

$$
-\sum_{n=1}^{\infty} \sum_{m=1}^{\infty}\left\{\left(\frac{n \pi}{K_{7}}\right)^{2}+m^{2}\right\} b_{m n} \sin \frac{n \pi x}{K_{7}} \cos m y=\frac{2}{J(x, y)} .
$$

Since the elements of $S_{1}$ are orthonormal, it follows that
(i) $a_{n}=\frac{-4 K_{7} K_{1}^{2} K_{2}^{2}}{n^{2} \pi^{3} A} \int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x$
(ii) $\quad b_{m n}=\frac{-8 K_{7} \mathrm{~K}_{1}^{2} \mathrm{~K}_{2}^{2}}{\pi\left\{(\mathrm{n} \pi)^{2}+\left(\mathrm{mK}_{7}\right)^{2}\right\} \mathrm{A}}$

$$
\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K_{7}} \cos m y}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x
$$

Closed form expressions for the integrals contained in (i) and (ii) are derived in Appendix F. Using these results, the Fourier series solution of equation (7-6) can be written in closed form.

## APPROXIMATIONS OF THE VELOCITY $\phi(\xi, \eta, \beta)$

The purpose of this chapter is to outline a procedure for finding approximations to $\phi(\xi, \eta, \beta)$, the $z$-component of the velocity in dimensionless form. Recall that in Chapter 2, we made the following notation changes:
(i) $\xi_{3} \rightarrow x$
(ii) $\eta_{3} \rightarrow y$
(iii) $\bar{A}_{3} \rightarrow \bar{A}$
(iv) $\bar{B}_{3} \rightarrow B$
(v) $\overline{\mathrm{C}}_{3} \rightarrow \mathrm{C}$.

We now revert to the original notation so that $\left(\xi_{3}, \eta_{3}\right)$ again refers to a point on the rectangle (see Fig. 7).

The first step in our procedure is the selection of the eccentric annulus parameters: $r_{1}, r_{2}$, and $\rho$. The symbol $r_{1}$ represents the radius of the inner circle, $r_{2}$
*The figures for Chapter 8 are contained in Appendix A.
represents the radius of the outer circle and $\rho$ represents the eccentricity, i.e., the distance between the centers. Once $r_{1}, r_{2}$, and $\rho$ are chosen, we calculate the following additional parameters;
(i) The area $A$ of the original annulus:

$$
A=\pi\left(r_{2}^{2}-r_{1}^{2}\right)
$$

(ii) The location of the annulus on the $\xi$-axis: $c_{1}, c_{2}$ (see Equations (2-8) and (2-9) and Fig. 2).
(iii) The parameters related to the Jacobian:
$K_{1}, \cdots, K_{4}$ (see Chapter 2).
(iv) The parameters related to the rectangle:
$K_{5}, \cdots, K_{7}$ (see Chapter 2) .
(v) The parameters related to the $\left(f_{i}, f_{j}\right)_{J}$ inner products and the coefficients of the fully developed velocity: $\mathrm{K}_{8}, \cdots, \mathrm{~K}_{15}$. (see Chapter 2, Appendix E, and Appendix F).
(vi) The parameter p defined in the transformation $\mathrm{W}_{1}(z)$ (see Equation (2-11)).
(vii) The minimum and maximum values of the Jacobian defined on the rectangle: $\operatorname{Min}(J), \operatorname{Max}(J)$ (see Chapter 2 and Fig. 7).
(viiị) A positive definite bound $\gamma$ for our operator $L(g)$ (see Equation (3-3)).

In terms of the coordinates $\xi_{3}, \eta_{3}$, and $\beta$, the velocity $\phi$ is given by

$$
\phi\left(\xi_{3}, \eta_{3}, \beta\right)=\phi_{e}\left(\xi_{3}, \eta_{3}, \beta\right)+\phi_{f}\left(\xi_{3}, \eta_{3}\right)
$$

(see Equation (1-19)). The fully developed velocity $\phi_{f}\left(\xi_{3}, \eta_{3}\right)$ is given by

$$
\phi_{f}\left(\xi_{3}, \eta_{3}\right)=\sum_{p=1}^{\infty} a_{p} \sin \frac{p \pi x}{k_{7}}+\sum_{p=1}^{\infty} \sum_{q=1}^{\infty} b_{p q} \sin \frac{p \pi x}{k_{7}} \cos q y
$$

(see Equation (7-9)). We approximate it by truncating this series when the coefficients $a_{p}, b_{p q}$ become smaller than some preassigned value. We now approximate the difference velocity $\phi_{e}$ :

$$
\phi_{e}\left(\xi_{3}, \eta_{3}, \beta\right)=\sum_{j=1}^{\infty} c_{j} g_{j}\left(\xi_{3}, \eta_{3}\right) \exp \left(-\lambda_{j} \beta\right)
$$

(see Equation (1-16)). First, we select the number, say n , of $\mathrm{g}_{\mathrm{j}}$ 's to use in our approximation of $\phi_{e}$. Then we select $n$ coordinate functions from $F=S_{6} \mathrm{US}_{7}$ to use to approximate the $\mathrm{g}_{\mathrm{j}}$ 's (see Chapter 4). The approximation $h_{n j}$ of $g_{j}$ is given by

$$
h_{n j}=\sum_{i=1}^{n} a_{i j}^{n} f_{i} \quad(j=1, \cdots, n)
$$

(see Chapter 5). We must determine the approximation $\lambda_{n j}$ to $\lambda_{j}$ and the $a_{i j}^{n}$ 's for $j=1, \cdots, n$. To do this we first calculate the inner products $\left(f_{i}, f_{j}\right)$ and $\left(f_{i}, f_{j}\right)$ for $i, j=1, \cdots, n$. (see Chapter 6 and

Appendix E). We then form the matrices
and

Next, we find a Cholesky decomposition of $A_{n}$

$$
A_{n}=C_{n} C_{n}^{T}
$$

and form the matrix

$$
G_{n}=C_{n}^{-1} B C^{-T}
$$

(see Chapter 6). We then find the eigenvalues $\mu_{n j}$ and the eigenvectors $y_{n j}$ of $G_{n} \quad(j=1, \cdots, n)$.

It follows that

$$
\lambda_{n j}=\frac{1}{\mu_{n j}}
$$

and

$$
C_{n}^{-T} y_{n j}=\left(a_{i j}^{n}, \cdots, a_{n j}^{n}\right) \quad(j=1, \cdots, n)
$$

To complete our approximation of $\phi_{e}$, we need approximations for the $c_{j}$ 's. We find these by using the formula

$$
c_{j}=\frac{-\left(\phi_{f}, g_{j}\right)_{J}}{\left\|g_{j}\right\|_{J}^{2}}
$$

wherein we replace $\phi_{f}$ with its truncated series approximation and $g_{j}$ with $h_{\mathrm{nj}}$ (see Equation (6-3)).

At this point, we have our approximation of $\phi$ in terms of the coordinates $\xi_{3}, \eta_{3}$, and $\beta$, i.e., we have $\phi\left(\xi_{3}, \eta_{3}, \beta\right)$. Since we want $\phi(\xi, \eta, \beta)$, we make the following substitutions:

$$
\begin{aligned}
& \text { (i) } \xi_{3}=\frac{1}{2} \ln \left[\frac{(\xi+p)^{2}+\eta^{2}}{(\xi-p)^{2}+n^{2}}\right]-\mathrm{K}_{5} \\
& \text { (ii) } \eta_{3}=\arctan \left[\frac{-2 p \eta}{\xi^{2}+\eta^{2}-p^{2}}\right]
\end{aligned}
$$

(see Chapter 2). Recall from Chapter 1 that

$$
\beta\left(z^{*}\right)=\frac{z^{*}}{\sqrt{\mathrm{~A}} \mathrm{R}}
$$

We want the dimensionless variable $\beta$ to be a function of $z$,i.e., we want $\beta(z)=\frac{z}{\sqrt{A} R}$. It is shown in [2] that

$$
\beta(z)=\frac{z}{\sqrt{\mathrm{~A} R}}=\frac{\int_{0}^{\beta} \varepsilon(\beta) \mathrm{d} \beta}{\sqrt{\mathrm{~A} R}}
$$

where

$$
\varepsilon(\beta)=\frac{\frac{3}{2} \int_{A}{ }^{2}(\xi, \eta, \beta) \frac{\partial \phi_{e}}{\partial \beta} d A}{\int_{A} \phi(\xi, n, \beta) \frac{\partial \phi_{e}}{\partial \beta} d A}-2 .
$$

Using $\beta=\beta(z)$ in $\phi(\xi, \eta, \beta)$, we have the desired approximation to the velocity.

APPENDIX A
FIGURES


FIGURE 1

$$
\begin{aligned}
& \xi=\frac{x}{\sqrt{A}} \\
& \eta=\frac{y}{\sqrt{A}} \\
& \bar{A}=\frac{\pi}{A}\left(r_{2}^{2}-r_{1}^{2}\right)=1
\end{aligned}
$$

FIGURE 2


FIGURE 3


FIGURE 4


FIGURE 5



## APPENDIX B

SIMPLIFICATION OF THE GOVERNING EQUATIONS

The equations governing laminar flow along the axis of a duct (see Fig. 1) are the continuity equation

$$
\begin{equation*}
\frac{\partial \rho}{\partial t}+\operatorname{div}(\overline{\mathrm{V}})=0 \tag{B-1}
\end{equation*}
$$

and the $z$-direction momentum equation

$$
\begin{align*}
\rho \frac{D w}{D t}=f_{z}- & \frac{\partial p}{\partial z}+\frac{\partial}{\partial z}\left[\mu\left(2 \frac{\partial w}{\partial t}-\frac{2}{3} \operatorname{div} \bar{V}\right)\right] \\
& +\frac{\partial}{\partial x}\left[\mu\left(\frac{\partial w}{\partial x}+\frac{\partial u}{\partial z}\right)\right] \\
& +\frac{\partial}{\partial y}\left[\mu\left(\frac{\partial v}{\partial z}+\frac{\partial w}{\partial y}\right)\right] \tag{B-2}
\end{align*}
$$

where
(i) $\overline{\mathrm{V}}=(\mathrm{u}, \mathrm{v}, \mathrm{w})$ represents the velocity
(ii) $t$ represents time
(iii) $\rho(x, y, z, t)$ represents the density
(iv) $\mu(x, y, z, t)$ represents the viscosity
(v) $p(x, y, z, t)$ represents the pressure
(vi) $f=\left(f_{x}, f_{y}, f_{z}\right)$ represents the external body force
(vii) $\frac{D}{D t}$ represents the substantial derivative [4].

We are considering Newtonian fluid under incompressible steady state laminar conditions. We assume that there are no external body forces present. The effect of these assumptions on equations ( $B-1$ ) and ( $B-2$ ) will now be determined. The assumptions of incompressible flow and Newtonian fluid mean that $\rho$ and $\mu$ are constants, respectively. The continuity equation now simplifies to

$$
\begin{equation*}
\operatorname{div} \overline{\mathrm{V}}=0 \tag{B-3}
\end{equation*}
$$

and the $z$-direction momentum equation to

$$
\begin{equation*}
\rho \frac{D w}{D t}=f_{z}-\frac{\partial p}{\partial z}+\mu\left(\frac{\partial^{2} w}{\partial x^{2}}+\frac{\partial^{2} w}{\partial y^{2}}+\frac{\partial^{2} w}{\partial z^{2}}\right) \tag{B-4}
\end{equation*}
$$

The substantial derivative of $w, \frac{D w}{D t}$, is given by

$$
\frac{D w}{D t}=\frac{\partial w}{\partial t}+u \frac{\partial w}{\partial x}+v \frac{\partial w}{\partial y}+w \frac{\partial w}{\partial z}
$$

Since body forces are neglected,$\overline{\mathrm{f}}=\overline{0}$. The assumption
of steady state flow means that the flow in independent of time. The z-direction momentum equation now simplifies to

$$
\begin{equation*}
\overline{\mathrm{V}} \cdot \operatorname{grad} \mathrm{w}=-\frac{1}{\rho} \frac{\partial p}{\partial z}+\nu \nabla^{2} w \tag{B-5}
\end{equation*}
$$

where

$$
\nu=\frac{\mu}{\rho}
$$

is the kinematic viscosity.

## APPENDIX C

## PROPERTIES OF THE JACOBIAN

We now establish several properties of the Jacobian

$$
J(x, y)=\frac{A}{K_{1}^{2} K_{2}^{2}}\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}
$$

THEOREM C-1: For a fixed $x$ in $\left[0, K_{7}\right], J(x, y)$ is increasing on $[0, \pi]$ and for a fixed $y$ in $[0, \pi], J(x, y)$ is increasing on $\left[0, \mathrm{~K}_{7}\right.$ ].

PROOF:
Let $x$ be an element of $\left[0, K_{7}\right]$. Then

$$
\frac{\partial J}{\partial y}=\frac{2 A}{K_{1}^{2} K_{2}^{2}}\left(\cosh \left(x+K_{5}\right)-\cos y\right) \sin y .
$$

Since $x$ is nonnegative and $K_{5}$ is positive, $\cosh \left(x+K_{5}\right)$ is positive so that $\left(\cosh \left(x+K_{5}\right)-\cos y\right)$ is positive. On $(0, \pi), \sin y$ is positive and for $y=0, \pi, \sin y=0$. Therefore $\frac{\partial J}{\partial y}$ is positive on $(0, \pi)$ which implies that $J(x, y)$ is increasing on [ $0, \pi$ ] for a fixed $x$. Let $y$ be
an element of $[0, \pi]$. Then

$$
\frac{\partial J}{\partial x}=\frac{2 A}{K_{1}^{2} K_{2}^{2}}\left(\cosh \left(x+K_{5}\right)-\cos y\right) \sinh \left(x+K_{5}\right)
$$

Since $\sinh \left(x+K_{5}\right)$ is positive on $\left[0, K_{7}\right]$, we have that $\frac{\partial J}{\partial x}$ is positive on $\left[0, K_{7}\right]$. Hence $J(x, y)$ is increasing on $\left[0, K_{7}\right]$ for a fixed $y$. The proof is now complete.

COROLLARY 1: We have the following:
(i) $\frac{\partial J(x, 0)}{\partial y}=0$ for $x$ in $\left[0, K_{7}\right]$.
(ii) $\frac{\partial J(x, \pi)}{\partial y}=0 \quad$ for $x$ in $\left[0, K_{7}\right]$.

COROLLARY 2: We have the following:
(i) The minimum of $J(x, y)$ on the rectangle is

$$
J(0,0)=\frac{\pi\left(r_{2}^{2}-r_{1}^{2}\right)\left(\left[r_{2}-\rho\right]^{2}-r_{1}^{2}\right)}{r_{2}^{2}\left(\left[r_{2}+\rho\right]^{2}-r_{1}^{2}\right)}
$$

(ii) The maximum of $J(x, y)$ on the rectangle is

$$
J\left(K_{7}, \pi\right)=\frac{\pi\left(r_{2}^{2}-r_{1}^{2}\right)\left(r_{2}^{2}-\left[r_{1}-\rho\right]^{2}\right)}{r_{1}^{2}\left(r_{2}^{2}-\left[r_{1}+\rho\right]^{2}\right)}
$$

## APPENDIX D

CALCULATIONS OF IMPORTANT INTEGRALS

The purpose of this appendix is to evaluate the following integrals:

$$
\begin{aligned}
& \text { (i) } I_{1}=\int_{-\pi}^{\pi} \frac{\cos (r y)}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y \\
& \text { (ii) } I_{2}=\int_{0}^{K_{7}} \frac{\exp (-r x) \sin \frac{s \pi x}{K_{7}}\left[\cosh \left(x+K_{5}\right)+r \sinh \left(x+K_{5}\right)\right]}{\sinh ^{3}\left(x+K_{5}\right)} d x \\
& \text { (iii) } I_{3}=\int_{0}^{K_{7}} \frac{\exp (-r x) \cos \frac{s \pi x}{K_{7}}\left[\cosh \left(x+K_{5}\right)+r \sinh \left(x+K_{5}\right)\right]}{\sinh ^{3}\left(x+K_{5}\right)} d x
\end{aligned}
$$

where $r$ and $s$ are integers. We have that

$$
I_{1}=\operatorname{Re} \int_{-\pi}^{\pi} \frac{\exp (\operatorname{iry})}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y
$$

For convenience we set

$$
z_{0}=\exp \left(x+K_{5}\right)
$$

and

$$
z_{1}=\exp \left(-\left(x+K_{5}\right)\right)
$$

If we make the substitution $z=\exp (i y)$, then $I_{1}$ becomes

$$
\begin{equation*}
I_{1}=4 \operatorname{Re}\left\{-i \oint_{C} f(z) d z\right\} \tag{D-1}
\end{equation*}
$$

where

$$
f(z)=\frac{z^{r+1}}{\left(z-z_{0}\right)^{2}\left(z-z_{1}\right)^{2}}
$$

and

$$
C:|z|=1
$$

is the unit circle. Since $x$ is nonnegative and $K_{5}$ is positive, the Residue Theorem implies that

$$
\begin{equation*}
\oint_{C} f(z) d z=2 \pi i \operatorname{Res}\left(z_{1}\right) \tag{D-2}
\end{equation*}
$$

where

$$
\operatorname{Res}\left(z_{1}\right)=\operatorname{Lim}_{z \rightarrow z_{1}} \frac{d}{d z}\left\{\left(z-z_{1}\right)^{2} f(z)\right\}
$$

It is not difficult to show that

$$
\begin{equation*}
\operatorname{Res}\left(z_{1}\right)=\frac{(r-1) z_{o}^{(r+1)}-(r+1) z_{o}^{(r-1)}}{\left(z_{0}-z_{1}\right)^{3}} \tag{D-3}
\end{equation*}
$$

Substitution of equation (D-3) into equation (D-2) and (D-2) into equation (D-1) yields, after some manipulation,

$$
I_{1}=\frac{2 \pi \exp \left(-r\left(x+K_{5}\right)\right)\left\{\cosh \left(x+K_{5}\right)+r \sinh \left(x+K_{5}\right)\right\}}{\sinh ^{3}\left(x+K_{5}\right)} .(D-4)
$$

We now evaluate $I_{2}$. Consider the following functions:
(i) $f(x)=-\frac{1}{2} \frac{\exp (-r x) \sin \frac{s \pi x}{K_{7}}}{\sinh ^{2}\left(x+K_{5}\right)}$
(ii) $g(x)=-\frac{r}{2} \exp (-r x) \sin \frac{s \pi x}{K_{7}} \operatorname{coth}\left(x+K_{5}\right)$
(iii) $h(x)=-\frac{s \pi}{2 K_{7}} \exp (-r x) \cos \frac{s \pi x}{K_{7}} \operatorname{coth}\left(x+K_{5}\right)$.

A straight forward calculation shows that

$$
\begin{align*}
I_{2}=[f(x) & +g(x)+h(x)]_{0}^{K_{7}} \\
& -\frac{1}{2}\left\{r^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}\right\} I \tag{D-5}
\end{align*}
$$

where

$$
I=\int_{0}^{K_{7}} \exp (-r x) \sin \frac{s \pi x}{K_{7}} \operatorname{coth}\left(x+K_{5}\right) d x
$$

Using the well-known fact that for $\mathrm{x}+\mathrm{K}_{5}>0$,

$$
\operatorname{coth}\left(x+K_{5}\right)=1+2 \sum_{j=1}^{\infty} \exp \left(-2 j\left(x+K_{5}\right)\right.
$$

we obtain

$$
\begin{align*}
I= & \int_{0}^{K_{7}} \exp (-r x) \sin \frac{s \pi x}{K_{7}} d x \\
& \quad+2 \sum_{j=1}^{\infty} \exp \left(-2 j K_{5}\right) \int_{0}^{K_{7}} \exp (-(r+2 j) x) \sin \frac{s \pi x}{K_{7}} d x . \tag{D-6}
\end{align*}
$$

Evaluation of $[f(x)+g(x)+h(x)]{ }_{0}^{K_{7}}$ and equation (D-6)
leads to, after considerable simplification, the following expression for $\mathrm{I}_{2}$ :

$$
\left.\begin{array}{rl}
I_{2}= & -\frac{s \pi}{2 K_{7}}\left\{(-1)^{s}\left(\frac{K_{8}}{K_{10}}\right)^{r} K_{11}-K_{15}\right\} \\
& -\frac{s \pi}{2 K_{7}}\left\{1+(-1)^{s+1} \sum_{\left.\left(\frac{K_{8}}{K_{12}}\right)^{r}\right\}}\right. \\
& -\frac{s \pi}{K_{7}}\left\{r^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}\right\} \sum_{j=1}^{\infty} K_{12}^{2 j}\left\{\frac{1+(-1)^{s+1}\left(\frac{K_{8}}{K_{10}}\right)^{r+2 j}}{(r+2 j)^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}}\right\}
\end{array}\right\}
$$

We now evaluate $I_{3}$. Consider the following functions:
(i) $f(x)=\frac{\frac{3}{2}}{2} \frac{\exp (-r x) \cos \frac{s \pi x}{K_{7}}}{\sinh ^{2}\left(x+K_{5}\right)}$
(ii) $g(x)=-\frac{r}{2} \exp (-r x) \cos \frac{s \pi x}{K_{7}} \operatorname{coth}\left(x+K_{5}\right)$
(iii) $h(x)=\frac{s \pi}{K_{7}} \exp (-r x) \sin \frac{s \pi x}{K_{7}} \operatorname{coth}\left(x+K_{5}\right)$.

It is not difficult to show that

$$
\begin{aligned}
I_{3} & =[f(x)+g(x)+h(x)]_{0}^{K_{7}} \\
& -\frac{3}{2}\left\{r^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}\right\} \int_{0}^{K_{7}} \exp (-r x) \cos \frac{s \pi x}{K_{7}} \operatorname{coth}\left(x+K_{5}\right) d x
\end{aligned}
$$

Calculations similar to those made in the last paragraph then show that

$$
\begin{aligned}
I_{3} & =-\frac{3}{2}\left\{(-1)^{s} K_{10}\left(\frac{K_{8}}{K_{12}}\right)^{r}-K_{14}\right\} \\
& -\frac{r}{2}\left\{(-1)^{s}{\left.\left(\frac{K_{8}}{K_{12}}\right)^{r} K_{11}-K_{15}\right\}}\right. \\
& -\left\{r^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}\right\} \sum_{j=1}^{\infty} K_{12}^{2 j}\left\{\frac{(r+2 j)\left\{(-1)^{s+1}\left(\frac{K_{8}}{K_{10}}\right)^{r+2 j}+1\right\}}{(r+2 j)^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}}\right)
\end{aligned}
$$

## APPENDIX E

THE $\left(\mathrm{f}_{\mathrm{i}}, \mathrm{f}_{\mathrm{j}}\right)_{\mathrm{J}}$ INNER PRODUCTS FOR THE GALERKIN DETERMINANT

The $\left(f_{i}, f_{j}\right)_{J}$ inner products needed for the Galerkin determinant require evaluation of an integral of the form

$$
I=\int_{0}^{K_{7}} \int_{0}^{\pi} \frac{\sin \frac{t \pi x}{K_{7}} \cos u y \sin \frac{v \pi x}{K_{7}} \cos w y}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x
$$

where $t$ and $v$ are positive integers and $u$ and $w$ are nonnegative integers. Using equation ( $D-4$ ) and the trigonometric identities
(i) $\cos u y \cos w y=\frac{1}{2}(\cos (u-w) y+\cos (u+w) y)$
(ii) $\sin \frac{t \pi x}{K_{7}} \sin \frac{v \pi x}{K_{7}}=\frac{1}{2}\left(\cos \frac{(t-v) \pi x}{K_{7}}-\cos \frac{(t+v) \pi x}{K_{7}}\right)$,
we can show that

$$
I=\frac{\pi}{4}\left(I_{1}-I_{2}+I_{3}-I_{4}\right)
$$

where $I_{1}, I_{2}, I_{3}$, and $I_{4}$ are of the form
$I^{\prime}=\exp \left(-r K_{5}\right) \int_{0}^{K_{7}} \frac{\exp (-r x) \cos \frac{s \pi x}{K_{7}}\left\{\cosh \left(x+K_{5}\right)+r \sinh \left(x+K_{5}\right)\right\}}{\sinh ^{3}\left(x+K_{5}\right)} d x$.

Using the equation (D-8) , I' simplifies to

$$
\begin{aligned}
& I^{\prime}=-\frac{1}{2}\left\{(-1)^{s} K_{10} K_{8}^{\mathrm{K}}-\mathrm{K}_{14} \mathrm{~K}_{12}^{\mathrm{r}}\right\} \\
& -\frac{\mathrm{r}}{2}\left\{(-1)^{\mathrm{s}} \mathrm{~K}_{8}^{\mathrm{r}} \mathrm{~K}_{9}-\mathrm{K}_{12}^{\mathrm{r}} \mathrm{~K}_{13}\right\} \\
& -\left\{r^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}\right\} \sum_{j=1}^{\infty}(r+2 j)\left\{\frac{K_{12}^{r+2 j}-(-1)^{s} K_{8}^{r+2 j}}{(r+2 j)^{2}+\left(\frac{s \pi}{K_{7}}\right)^{2}}\right\} \text {. }
\end{aligned}
$$

The values of $r$ and $s$ for $I_{1}, I_{2}, I_{3}$, and $I_{4}$ are given by

$$
\begin{aligned}
& \text { (i) } I_{1}: r=u-w ; s=t-v \\
& \text { (ii) } I_{2}: r=u-w ; s=t+v \\
& \text { (iii) } I_{3}: r=u+w ; s=t-v \\
& \text { (iv) } I_{4}: r=u+w ; s=t+v .
\end{aligned}
$$

## APPENDIX F

## THE COEFFICIENTS OF THE FULLY DEVELOPED VELOCITY

The coefficients $a_{n}$ and $b_{m n}$ in the Fourier series expansion of the fully developed velocity require evaluation of integrals of the form

$$
I=\int_{0}^{K} 7 \int_{0}^{\pi} \frac{\sin \frac{n \pi x}{K 7} \cos m y}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y d x
$$

where $n$ is a positive integer and $m$ is a nonnegative integer. We established in Appendix $D$ that

$$
\begin{aligned}
& \int_{0}^{\pi} \frac{\cos m y}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d y \\
& \quad=\frac{3}{2} \int_{-\pi\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}}^{\pi} d y \\
& =\frac{\cos m y}{\sinh ^{3}\left(x+K_{5}\right)}
\end{aligned}
$$

Using another result from Appendix D, we obtain

$$
\begin{aligned}
I & =\pi \exp \left(-m K_{5}\right) \int_{0}^{K_{7}} \frac{\exp (-m x) \sin \frac{n \pi x}{K_{7}}\left\{\cosh \left(x+K_{5}\right)+m \sinh \left(x+K_{5}\right)\right\}}{\left(\cosh \left(x+K_{5}\right)-\cos y\right)^{2}} d x \\
& =-\frac{n_{\pi}^{2}}{2 K_{7}}\left\{(-1)^{n} K_{8}^{m} K_{9}-K_{12}^{m} K_{13}\right\} \\
& -\frac{n \pi}{K_{7}}\left\{m^{2}+\left(\frac{n \pi}{K_{7}}\right)^{2}\right\} \sum_{j=1}^{\infty}\left\{\frac{K_{12}-(-1)^{n} K_{8}^{m+2 j}}{(m+2 j)^{2}+\left(\frac{n \pi}{K_{7}}\right)^{2}}\right\}
\end{aligned}
$$

## APPENDIX G

THE HILBERT SPACE $\mathrm{W}_{2}^{(1)}(\overline{\mathrm{A}})$
Consider $L_{2}(\bar{A})$ where $\bar{A}$ is a compact domain in $E^{2}$. We say that an element $g(x, y)$ of $L_{2}(\bar{A})$ vanishes in a boundary layer if there exists a positive number $\delta$ such that for all ( $x, y$ ) in the set

$$
\overline{\mathrm{A}}_{\delta}=\{(x, y) \text { in } \overline{\mathrm{A}} \mid \operatorname{dist}[(x, y), \operatorname{Bdy}(\overline{\mathrm{A}})]<\delta\}
$$

$g(x, y)=0$. Define $M_{1}(\bar{A})$ to be the set of all continuously differentiable functions in $L_{2}(\bar{A})$ which vanish in a boundary layer. For $g$ in $M_{1}(\bar{A})$,

$$
D_{x} g=\frac{\partial g}{\partial x}
$$

exists. Since $M_{1}(\bar{A})$ is dense in $L_{2}(\bar{A})$, the adjoint $D_{x}^{*}$ of $D_{x}$ exists $[7,13]$. We call the operator $-D_{x}^{*}$ the operation of generalized differentiation with respect to $x$. If $h$ is an element of the domain of $D_{x}^{*}$, we call $-D{ }_{x}^{*} h$ the generalized first derivative of $h$ with respect to $x$. If $h$ has continuous derivatives in $\vec{A}$, then

$$
-D{ }_{x}^{*} h=\frac{\partial h}{\partial x}
$$

i.e.,the generalized derivative coincides with the ordinary derivative. Similar remarks can be made with $x$ replaced with $y$.

Consider the set $W_{2}^{(1)}(\bar{A})=W$ of all elements of
$L_{2}(\bar{A})$ which have generalized first derivatives. An inner product on $W$ is defined by

$$
\begin{array}{rl}
(g, h)_{W}=\int_{\bar{A}} g h & d \bar{A}+\int_{\bar{A}} D_{x}^{*} g D_{x}^{*} h d \bar{A} \\
& +\int_{\bar{A}} D_{y}^{*} g D_{y}^{*} h d \bar{A}
\end{array}
$$

and the corresponding norm by

$$
\begin{gathered}
\|\left. g\right|_{W} ^{2}=\int_{-} g^{2} d \bar{A}+\int_{\bar{A}} D_{x}^{* 2} g d \bar{A} \\
+\int_{\bar{A}} D_{y}^{* 2} g d \bar{A}
\end{gathered}
$$

With these definitions, W becomes a Hilbert space [7].

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[^0]:    *Numbers in brackets refer to references in the bibliography.

[^1]:    *We, of course, consider functions equivalent if they differ only on a set of measure zero. In the sequel we will usually require continuous representations from the resulting equivalence classes.

