# APPLICATION OF THE RITZ METHOD TO NON-STANDARD EIGENVALUE PROBLEMS 

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## 1. Introduction

There is an extensive literature on application of the Ritz method to eigenvalue problems of the type

$$
\begin{equation*}
L_{1} w=\lambda L_{2} w \tag{1}
\end{equation*}
$$

where $L_{1}, L_{2}$ are positive definite linear operators in a Hilbert space (see for example [1]). The classical theory concerns the case in which there exists a minimum (or maximum) eigenvalue, and subsequent eigenvalues can be located by a well-known mini-max principle [2; p. 405]. This paper considers the possibility of application of the Ritz method to eigenvalue problems of the type (1) where the linear operators $L_{1}, L_{2}$ are not necessarily positive definite and a minimum (or maximum) eigenvalue may not exist. The special cases considered may be written with the eigenvalue occurring in a non-linear manner.

This investigation was suggested by results [3] recently obtained when the Ritz method was applied to a variational formulation by Chandrasekhar [4] of the boundary value problem describing non-radial oscillations of stars. This problem, which is described in section 2 , does not have a minimum (or maximum) eigenvalue and hence is not dealt with by the classical theory. It is the aim of section 2 to indicate the use that has been made by astrophysicists of the methods whose validity is examined in the succeeding sections, to discuss the results they have obtained, and to examine the mathematical properties of the eigenvalue problem they considered, so that later the similarities between that problem and the one considered in section 4 may become apparent.

Some elementary results for quite general operators are established in section 3. Section 4 considers a simpler problem for which an exact solution may be obtained but which is in many ways similar to that studied in [3]. Like the problem considered in [3], this simpler problem has no minimum or maximum eigenvalue, though its eigenvalues are bounded from below. Some theorems, which it is hoped are also interesting in their own right,
are proved concerning the results of formal application of the Ritz method to this simpler problem. Section 5 reconsiders the astrophysical problem in the light of the results of section 4 which are also extended slightly.

## 2. Non-radial oscillations of stars

The fourth order system of differential equations governing small adiabatic non-radial oscillations of stars is derived and its properties discussed in [5; p. 509]. This system, and the simpler second order system obtained when, as in [3] and many other papers, the Eulerian perturbation $\delta \phi$ of gravitational potential is neglected, will be denoted here, as in [3], by $I A$ and $I B$ respectively.

When $\delta \phi$ is neglected, Chandrasekhar's variational principle for oscillations corresponding to the spherical harmonic $\ell$ may be written

$$
\begin{align*}
& \sigma^{2} \int_{0}^{R} \rho\left(\frac{\psi^{2}}{r^{2}}+\frac{\left(\chi^{\prime}\right)^{2}}{\ell(\ell+1)}\right) d r  \tag{2}\\
& \quad=\int_{0}^{R}\left\{\Gamma_{1} p\left(\frac{d p}{d r} \frac{\psi}{\Gamma_{1} p}+\psi^{\prime}-\chi^{\prime}\right)^{2}+A \frac{d p}{d r} \psi^{2}\right\} \frac{d r}{r^{2}}
\end{align*}
$$

where $2 \pi / \sigma$ is the period, $R$ is the radius,

$$
A=\frac{1}{\rho} \frac{d \rho}{d r}-\frac{1}{\Gamma_{1} p} \frac{d p}{d r}
$$

and the components $\psi, \chi^{\prime}$ of the vector-valued eigenfunction, which depend on the radial and transverse components of displacement, are defined in [4]. The density $\rho$, the total pressure $p$, and the adiabatic exponent $\Gamma_{1}$ are functions of the distance $r$ from the centre. These functions depend on the stellar model chosen. It is required that $\psi / r^{2}$ and $\chi^{\prime} / r$ be bounded at $r=0$ and $r=R$. If the stellar model satisfies certain conditions that are satisfied by all the stellar models considered in this paper, Lebovitz has shown [6] that if $\sigma^{2} \neq 0$, the eigenfunctions of $I A$ must satisfy stronger conditions (which are satisfied by (6)) at the centre. An identical proof establishes the same result for $I B$.

The Euler-Lagrange equations obtained from (2) by considering variations in $\psi$ and $\chi^{\prime}$ respectively are

$$
\begin{align*}
& \sigma^{2} \rho \psi / r^{2}=-F^{\prime}(r)+\left(p^{\prime} \mid \Gamma_{1} p\right) F(r)+A p^{\prime} \psi / \gamma^{2}  \tag{3a}\\
& \sigma^{2} \rho \chi^{\prime} / \ell(\ell+1)=-F(r) \tag{3b}
\end{align*}
$$

where $F(r)=\left(p^{\prime} \psi+\Gamma_{1} p \psi^{\prime}-\Gamma_{1} p \chi^{\prime}\right) / r^{2}$. Equations (12), (31), (33) and (34) of [4] show that equations (3a) and (3b) above are scalar multiples of the
radial and transverse components of the basic equation (9) of [4] when $\delta \phi$ is neglected.

Equation (3) is of the form (1). Denote the linear operators $L_{1}, L_{2}$ in this case by $A_{1}, A_{2}$ respectively. In this case $w$ corresponds to the vector valued function $\left\{\psi, \chi^{\prime}\right\}$, the inner product is defined by

$$
\left(\left\{\psi_{1}, \chi_{1}^{\prime}\right\},\left\{\psi_{2}, \chi_{2}^{\prime}\right\}\right)=\int_{0}^{R}\left(\psi_{1} \psi_{2}+\chi_{1}^{\prime} \chi_{2}^{\prime}\right) d r
$$

$\lambda=\sigma^{2}$ and $A_{1}$ and $A_{2}$ are symmetric, but not positive definite. (That the operators are symmetric, as is necessary for the variational formulation, was proved in [4].) As $A_{2}$ is given by $A_{2}\left\{\psi, \chi^{\prime}\right\}=\left\{\rho \psi / r^{2}, \rho \chi^{\prime} \mid \ell(\ell+1)\right\}$, the orthogonality relation satisfied by eigenfunctions $\left\{\psi_{1}, \chi_{1}^{\prime}\right\},\left\{\psi_{2}, \chi_{2}^{\prime}\right\}$ corresponding to distinct eigenvalues is

$$
\int_{0}^{R}\left(\rho \psi_{1} \psi_{2} / r^{2}+\rho \chi_{1}^{\prime} \chi_{2}^{\prime} l f(\ell+1)\right) d r=0 .
$$

Since $\ell>0$, and $\rho>0$ for $r<R$ in all models considered here (and all physical models), the operator $A_{2}$ is positive (that is for all non-zero $w$ in the domain of $\left.A_{2},\left(A_{2} w, w\right)>0\right)$. Hence (3) satisfies the conditions of both Theorems 1 and 2.

Substituting for $F(r)$ from (3b) in (3a) shows that when $\sigma^{2} \neq 0$ equation (3) may be written in standard matiix form as
(4) $\quad\binom{\psi}{\chi^{\prime}}^{\prime}=\left(\begin{array}{ll}-\frac{1}{\Gamma_{1} p} \frac{d p}{d r} & 1-\frac{\sigma^{2} \rho r^{2}}{\Gamma_{1} \not p \ell(\ell+1)} \\ \frac{\ell(\ell+1)}{r^{2}}\left(1-\frac{A}{\rho \sigma^{2}} \frac{d p}{d r}\right) & -A\end{array}\right)\binom{\psi}{\chi^{\prime}}$.

This differs only in notation from the form given for $I B$ in [5]. When Schwarzschild's stability criterion $A<0$ is satisfied throughout the star (except perhaps on a set of measure zero) it is obvious from (2) that, as $d p \mid d r<0$ and $\Gamma_{1} p>0$, all eigenvalues are strictly positive so that the above formulation is valid. Although (3) is written with $\sigma^{2}$ occurring in a linear manner, when that equation is written as the second order differential equation (4) (instead of as a third order one), $\sigma^{2}$ occurs in a non-linear manner. This property, which is also exhibited by the equation considered in section 4, appears to be intimately related to the properties of the spectrum.

Numerical solutions of $I A$ and $I B$ have been obtained for a number of stellar models. Generally these follow the pattern predicted by Cowling [7]. The spectrum is discrete, there being two families of modes - the $p$-modes with eigenvalues tending to $+\infty$ and the $g$-modes with eigenvalues tending to zero. In most cases the $n$th $p$-mode and the $n$th $g$-mode both have
exactly $n$ zeros in $\psi$ and $n$ zeros in $\chi^{\prime}$. In between these two families is a single mode, the $f$-mode, with no zeros in $\psi$ or $\chi^{\prime}$. Variations in this pattern have been noted for some models. (See for example [8], [9] and [5; p. 515].)

An exception relevant to what follows is the convective model, in which $A \equiv 0$. In this case it is obvious that $\sigma^{2}=0$ is an eigenvalue of (3) of infinite multiplicity, with the corresponding space of eigenfunctions given by $\chi^{\prime}=\psi^{\prime}+p^{\prime} \psi / \Gamma_{1} p$, where $\psi / r^{2}$ is any bounded function sufficiently differentiable on ( $0, R$ ). In this case (3) may be written

$$
\begin{equation*}
\sigma^{2}\left[\left(\frac{r^{2} p_{1} \chi^{\prime \prime}}{\ell(\ell+1)}\right)^{\prime}+\left(\frac{p_{1} \rho r^{2} \sigma^{2}}{\Gamma_{1} p \ell(\ell+1)}-p_{1}\right) \chi^{\prime}\right]=0 \tag{5}
\end{equation*}
$$

where $p_{1}(r)=\exp \left(\int_{0}^{r}\left(p^{\prime} \mid \Gamma_{1} p\right) d r\right)$. The non-zero eigenvalues are given by the equation obtained by dividing (5) by $\sigma^{2}$. With the boundary conditions this is a singular Sturm-Liouville problem in $\chi^{\prime}[2 ; p .324]$, since $\rho, \Gamma_{1} p, p_{1}$ and $\ell$ are all strictly positive in $[0, R$ ). The singularities are regular (compare [ 5 ; p. 461]) when, as is usual, $\rho / \Gamma_{1} p \sim(R-r)^{-1}$ as $r \rightarrow R$. It seems likely, as is generally accepted [5; p. 518], that with the convective model $I B$ has an $f$-mode and an infinite family of $p$-modes but that instead of the infinite family of $g$-modes there is an infinite family of eigenfunctions with eigenvalue zero.

In [3] the Ritz method was applied to (2) by assuming $\psi / r^{\ell+1}$ and $\chi^{\prime} \mid r^{\ell}$ to be continuous and piecewise linear. Good results were obtained for the $f$-mode and the $p$-modes but no trace of the $g$-modes was obtained. Of the eigenvalues and eigenvectors of the matrix equation yielded by the Ritz method, those which did not correspond to $p$-modes or the $f$-mode seemed completely spurious.

The only model considered in [3] was one constructed by Van der Borght [10]. It has a convective core and a radiative envelope with $A=0$ for $r<r_{f}<R$ and $A<0$ for $r_{f} \leqq r \leqq R$. (The possible effect on the solution of $I B$, of the vanishing of $A$ throughout a region, is considered in section 5.) For this model, Wan and Van der Borght [11] have solved both $I A$ and $I B$ by numerical integration and obtained both $p$-modes and $g$-modes in both cases.

The Ritz method, using similar continuous piecewise linear coordinate functions, has been applied successfully to the problem of purely radial oscillations of this model [12]. In that case however the associated differential operator is positive definite and the variational principle is of the classical minimum type. The existence of the $g$-modes with eigenvalues tending to zero ensures that the operator $A_{1}$ associated with non-radial oscillations is not positive definite. (A symmetric linear operator $L$ is said to be positive definite if there is a strictly positive constant $k$ such that $(L w, w) \geqq k(w, w)$ for all $w$ in the domain of $L$.) The variational principle
(2) is not a classical minimum principle as there is no minimum eigenvalue.

A second application of the Ritz method to (2) in [3] assumed $\psi$ and $\chi^{\prime}$ to be given by

$$
\begin{equation*}
\psi=\sum_{i=0}^{n} C_{i} r^{\ell+2 i+1}, \quad \chi^{\prime}=(\ell+1) C_{0} r^{\ell}+\sum_{i=1}^{n} C_{i}^{*} r^{\ell+2 i} \tag{6}
\end{equation*}
$$

where the $C_{i}, C_{i}^{*}$ are parameters. Coordinate functions of type (6) were initially used in [13] where only the case $n=1$ was considered and only the $f$-mode determined. However those numerical results in [3] which used (6) have subsequently been found to be incorrect. Amended results using values of $n$ up to $n=6$ are given in a later paper [20]. These amended results show that when (6) is used with this model crude approximations of $g$-modes may be obtained for both $I B$ and $I A$ but that they are far less accurate than the approximations obtained for $f$ - and $p$-modes. These calculations, like those in [3] and [11], consider only the spherical harmonic $\ell=2$.

Robe and Brandt [14] have applied the Ritz method to both $I A$ and $I B$ for several polytropic models. For $\ell=2$ they obtained both $p$ - and $g$-modes in both cases although, except for models of unrealistically low central condensation, the values they obtained, especially for $g$-modes, differed considerably from those obtained by numerical integration of the differential equation. They assumed $\psi$ and $\chi^{\prime}$ to be given by (6), using values of $n$ up to $n=2$. Their relatively good results for models of low central condensation may be explained by the fact that (6) is the form of the exact eigenfunctions for the homogeneous (constant density) model. For a model with $\rho$ proportional to $1-(r / R)^{2}$, which also has very low central condensation, Tassoul [15] obtained encouraging results for both $p$ - and $g$-modes of $I A$ using the Ritz method, again with $\psi$ and $\chi^{\prime}$ given by (6). However in [15] comparison of the results of the Ritz method with accurate results is limited to the case of purely radial oscillations when, as already noted, the Ritz method presents no difficulty. It cannot be assumed that the Ritz method will be as successful for the more realistic model used in [3] as it was for Tassoul's model or the low index polytropes. As well as having higher central condensation, Van der Borght's model has $A \equiv 0$ throughout an extensive core and Cowling's argument [7] is weakest when $A$ vanishes throughout some region.

No rigorous mathematical analysis has yet been undertaken for either $I A$ or $I B$ except with some highly unrealistic models [5, p. 514] and the numerical solutions obtained have not been rigorously justified. Accordingly before considering in more detail the reason for the partial failure of the Ritz method in this incompletely understood case, it would seem useful to consider a related problem which can be solved explicitly and shown
to have cluster points of eigenvalues at $+\infty$ and zero. This is done in section 4.

## 3. The general equation

First it would seem of interest to consider the general eigenvalue problem (1) where $L_{1}, L_{2}$ are symmetric linear operators in a Hilbert space. This equation may be regarded as the Euler-Lagrange equation of the variational problem $\left(L_{1} w, w\right)=\lambda\left(L_{2} w, w\right)$. Application of the Ritz method to this yields the matrix equation

$$
\begin{equation*}
P x=\Lambda Q x \tag{7}
\end{equation*}
$$

where the matrices $P$ and $Q$ are symmetric, and if the operators $L_{1}$ and $L_{2}$ are positive then $P$ and $Q$ are positive definite. The Ritz method has been rigorously justified [1] when the variational principle is a true mini-max principle, but this is not the case with the problems considered in sections 2 and 4 and is not assumed in this section.

Throughout the rest of this paper, eigenvalues and eigenfunctions yielded by formal application of the Ritz method to some problem will be termed Ritz approximate eigenvalues and eigenfunctions of that problem.

The following Theorem, which emphasises the importance of the choice of coordinate functions, will be required for later work. The proof is straightforward and is omitted.

Theorem 1. Let $L_{1}$ be any symmetric linear operator and $L_{2}$ any symmetric positive linear operator in a real Hilbert space of functions. Let $\lambda_{1}, \cdots, \lambda_{n}$ be any distinct eigenvalues of (1) and $w_{1}, \cdots, w_{n}$ the corresponding eigenfunctions (which are linearly independent since the $\lambda_{i}$ are distinct.) Let $\xi_{1}, \cdots, \xi_{n}$ be any basis of the space spanned by $w_{1}, \cdots, w_{n}$. Then when the Ritz method is applied to (1) using the $\xi_{i}$ as coordinate functions, the Ritz approximate eigenvalues obtained are the exact $\lambda_{i}$ and the corresponding Ritz approximate eigenfunctions are the corresponding exact $w_{i}$.

A particular case of Theorem 1 is the result obtained by Robe [16] when he applied the Ritz method to $I A$ for the homogeneous stellar model using as coordinate functions a basis of the space spanned by certain exact eigenfunctions, whose form in that very special (and non-physical) case is known [5; p. 514]. However this gives little indication of the reliability of the Ritz method when the solution is not known in advance.

In the theory of convergence of the Ritz method in the classical case, an important fact is that the Ritz approximate eigenfunctions corresponding to distinct eigenvalues satisfy the orthogonality relation of the true eigenfunctions. The method of proof normally used in that case is also valid more generally.

Theorem 2. Let $L_{1}, L_{2}$ be symmetric linear operators in a real Hilbert space of functions. Then if $\Lambda_{1}, \Lambda_{2}$ are distinct Ritz approximate eigenvalues of (1) and $W_{1}, W_{2}$ are the corresponding Ritz approximate eigenfunctions, $\left(L_{1} W_{1}, W_{2}\right)=\left(L_{2} W_{1}, W_{2}\right)=0$.

Proof. For $k=1$ and $2, W_{k}=\sum_{j=1}^{n} x_{k j} \xi_{j}$ where the coordinate functions $\xi_{j}$ are any linearly independent functions in the domains of $L_{1}$ and $L_{2}$ and the scalars $x_{k j}$ satisfy (7) with $\Lambda=\Lambda_{k}, x=x_{k}=\left\{x_{k 1}, \cdots, x_{k n}\right\}$, $p_{i j}=\left(L_{1} \xi_{i}, \xi_{j}\right)$, and $q_{i j}=\left(L_{2} \xi_{i}, \xi_{j}\right)$. Now

$$
\begin{aligned}
\left(L_{2} W_{1}, W_{2}\right) & =\left(\sum_{i=1}^{n} x_{1 i} L_{2} \xi_{i}, \sum_{j=1}^{n} x_{2 j} \xi_{j}\right) \\
& =\sum_{i=1}^{n} \sum_{j=1}^{n} x_{1 i} x_{2 j} q_{i j} \\
& =x_{2}^{\tau} Q x_{1} \text { (where } x^{\tau} \text { denotes the transpose of } x \text { ) } \\
& =0 \text { since } A_{1} \neq \Lambda_{2} \text { and } P \text { and } Q \text { are symmetric. }
\end{aligned}
$$

Similarly $\left(L_{1} W_{1}, W_{2}\right)=0$.

## 4. A simpler equation

Let $H$ be the Hilbert space of real square summable 2 -vector valued functions $\{u, v\}$ defined on the interval $[0,1]$, in which the inner product is defined by

$$
\begin{equation*}
\left(\left\{u_{1}, v_{1}\right\},\left\{u_{2}, v_{2}\right\}\right)=\int_{0}^{1}\left(u_{1} u_{2}+v_{1} v_{2}\right) \tag{10}
\end{equation*}
$$

where in (10) and throughout this section the abbreviated notation $\int f$ is used for $\int f(x) d x$.

Define the operators $B_{1}, B_{2}$ on those sufficiently differentiable functions in $H$ satisfying

$$
\begin{equation*}
v(0)=v(1)=0 \tag{11}
\end{equation*}
$$

by

$$
\begin{aligned}
B_{1}\binom{u}{v} & =\binom{a_{1} u+a_{2} v-a_{1} a_{4} v^{\prime}}{\left[a_{4}\left(a_{1} u+a_{2} v-a_{1} a_{4} v^{\prime}\right)\right]^{\prime}+a_{2} u+a_{3} v-a_{2} a_{4} v^{\prime}} \\
B_{2}\{u, v\} & =\left\{a_{5} u, a_{6} v\right\}
\end{aligned}
$$

where the $a_{i}$ are continuously differentiable, non-vanishing, real-valued functions defined on $[0,1]$. It follows from (10) and (11) that

$$
\begin{aligned}
\left(B_{1}\binom{u_{1}}{v_{1}},\binom{u_{2}}{v_{2}}\right)= & \int_{0}^{1}\left[a_{1}\left(u_{1}+\frac{a_{2}}{a_{1}} v_{1}-a_{4} v_{1}^{\prime}\right)\left(u_{2}+\frac{a_{2}}{a_{1}} v_{2}-a_{4} v_{2}^{\prime}\right)\right. \\
& \left.+\left(a_{3}-\frac{a_{2}^{2}}{a_{1}}\right) v_{1} v_{2}\right]
\end{aligned}
$$

and

$$
\left(B_{2}\left\{u_{1}, v_{1}\right\},\left\{u_{2}, v_{2}\right\}\right)=\int_{0}^{1}\left(a_{5} u_{1} u_{2}+a_{6} v_{1} v_{2}\right)
$$

These equations show that $B_{1}$ and $B_{2}$ are symmetric and that if, almost everywhere on ( 0,1 ),

$$
\begin{equation*}
a_{3}>a_{2}^{2} / a_{1}>0 \tag{12a}
\end{equation*}
$$

then $B_{1}$ is positive, while $B_{2}$ is positive if, almost everywhere on $(0,1)$,

$$
\begin{equation*}
a_{5}>0, a_{6}>0 \tag{12b}
\end{equation*}
$$

In the rest of this section it is assumed that (12) is satisfied since in the simplest case the operators $A_{1}, A_{2}$ of section 2 are positive and one purpose of the present investigation is to shed further light on the problem considered there. Other cases are mentioned briefly in section 5 .

In exactly the same manner in which (4) was derived from (3), the eigenvalue problem

$$
\begin{equation*}
B_{1}\{u, v\}=\lambda B_{2}\{u, v\} \tag{13}
\end{equation*}
$$

may be rewritten as

$$
\binom{u}{v}^{\prime}=\left(\begin{array}{cl}
-b_{1}-b_{6} & b_{5}-b_{3} / \lambda  \tag{14}\\
b_{2}-\lambda b_{4} & b_{1}
\end{array}\right)\binom{u}{v}
$$

where

$$
\begin{array}{lll}
b_{1}=a_{2} / a_{1} a_{4}, & b_{2}=1 / a_{4}, & b_{3}=\left(a_{1} a_{3}-a_{2}^{2}\right) / a_{1} a_{4} a_{5} \\
b_{4}=a_{5} / a_{1} a_{4}, & b_{5}=a_{6} / a_{4} a_{5}, & b_{6}=\left(a_{4} a_{5}\right)^{\prime} / a_{4} a_{5}
\end{array}
$$

The non-linear occurrence of $\lambda$ in (14) and its linear occurrence in (13) are reminiscent of (4) and (3).

If $a_{1} / a_{5}\left(=b_{2} / b_{4}\right)$ is a constant, (14) shows that $a_{1} / a_{5}$ is an eigenvalue of (13) with corresponding eigenfunction $\left\{u_{0}, v_{0}\right\}$ where

$$
u_{0}=c\left(a_{4} a_{5}\right)^{-1} \exp \left(-\int b_{1}\right), \quad v_{0} \equiv 0
$$

and $c$ is a constant. Denote this eigenvalue by $\lambda_{0}$. Substitution shows that $\left\{u_{0}, v_{0}\right\}$ is the only eigenfunction with $v \equiv 0$.

Application of the Ritz method to the solution of (13) involves extremising $\lambda=\left(B_{1}\{u, v\},\{u, v\}\right) /\left(B_{2}\{u, v\},\{u, v\}\right)$ where $\{u, v\}$ is restricted to a finite dimensional subspace of $H$. The coordinate functions (basis of the subspace) usually chosen are of the form $\left\{\xi_{i}, 0\right\},\left\{0, \eta_{j}\right\}$ so that it is assumed that

$$
\begin{equation*}
u=\sum_{i=1}^{n} c_{i} \xi_{i}, \quad v=\sum_{j=1}^{m} c_{j}^{*} \eta_{j} \tag{15}
\end{equation*}
$$

where the $c_{i}, c_{j}^{*}$ are constants and the $\left\{\xi_{i}, \eta_{j}\right\}$ are in the domain of $B_{1}$ and
hence in the domain of $B_{2}$. The coordinate functions used in [3], [14], [15] and [16] are of this type.

ThEOREM 3. If $a_{1} / a_{5}$ is constant, application of the Ritz method to (13) using (15) with $n>m$ will yield at least $n-m$ linearly independent Ritz approximate eigenfunctions in which $v \equiv 0$. The Ritz approximate eigenvalue corresponding to each of these will be $a_{1} / a_{5}$, the exact value of $\lambda_{0}$.

Proof. The matrix equation yielded by the Ritz method will be of the form (7) where in this case it is easily shown that

$$
P=\left(\begin{array}{ll}
a_{1}(0) P_{1} & P_{2}  \tag{16}\\
P_{2}^{\tau} & P_{3}
\end{array}\right), \quad Q=\left(\begin{array}{ll}
a_{5}(0) P_{1} & 0_{n m} \\
0_{m n} & P_{4}
\end{array}\right)
$$

where $P_{1}$ is an $n \times n$ matrix, $P_{3}$ and $P_{4}$ are $m \times m, P_{2}$ is $n \times m$ and $0_{m n}, 0_{n m}$ are zero matrices. If $\Lambda=a_{1} / a_{5}$, then $x^{\tau}=\{y, 0\}$, where $y$ is a $\mathbf{l} \times n$ vector and 0 the $1 \times m$ zero vector, will be a solution of (7) if and only if $P_{2}^{\tau} y^{\tau}=0$. Since $P_{2}^{\tau}$ is $m \times n$ and $n>m$ this will have a solution space of dimension at least $n-m$. The result follows. Alternatively the multiplicity of $\lambda_{0}$ could be established by noting that $\operatorname{det}(P-\Lambda Q)$ has a factor $\left(a_{1}-a_{5} \Lambda\right)^{n-m}$.

The remainder of this paper will be concerned with the case when the $a_{i}$ are all constant. In this case (14) is very simple as $b_{6}=0$ and the $b_{i}$ are all constant. To distinguish this case, (13) with all the $a_{i}$ constant will be labelled (13a).

The exact solution of (13a) may easily be found. Let

$$
d(\lambda)=b_{1}^{2}+b_{2} b_{5}+b_{3} b_{4}-\lambda b_{4} b_{5}-b_{2} b_{3} / \lambda
$$

Since $a_{1} / a_{5}$ is constant it has already been shown that $\lambda_{0}$ is an eigenvalue with unique corresponding eigenfunction $\left\{u_{0}, v_{0}\right\}$ and that no other eigenfunctions have $v \equiv \mathbf{0}$. For any given $\lambda$, the solutions of (13a), if any, may be found by the usual methods for homogeneous linear differential equations with constant coefficients. It is readily seen that there is no eigenvalue $\lambda$ (except perhaps $\lambda_{0}$ ) for which $d(\lambda)=0$ as (11) shows that this would require $v \equiv 0$. Hence all remaining eigenfunctions must be of the form

$$
\{u(x), v(x)\}=\left\{c_{1}, c_{2}\right\} \exp (\sqrt{ }(d(\lambda)) x)+\left\{c_{3}, c_{4}\right\} \exp (-\sqrt{ }(d(\lambda)) x)
$$

where the $c_{i}$ are constants. Boundary conditions (11) then show that the only solutions with $v$ not identically zero have $v(x)=c \sin (\sqrt{ }(-d(\lambda)) x)$, where $c$ is a constant and that this requires $d(\lambda)=-n^{2} \pi^{2}$ where $n$ is a positive integer, so that all eigenvalues except $\lambda_{0}$ must satisfy

$$
\begin{equation*}
\lambda^{2} b_{4} b_{5}-\lambda\left(b_{1}^{2}+b_{2} b_{5}+b_{3} b_{4}+n^{2} \pi^{2}\right)+b_{2} b_{3}=0 \tag{17}
\end{equation*}
$$

so that

$$
\begin{aligned}
\lambda=\left(2 b_{4} b_{5}\right)^{-1}\left[b_{1}^{2}+b_{2} b_{5}+b_{3} b_{4}+n^{2} \pi^{2}\right. & \pm\left(\left(b_{1}^{2}+b_{2} b_{5}+b_{3} b_{4}+n^{2} \pi^{2}\right)^{2}\right. \\
& \left.\left.-4 b_{2} b_{3} b_{4} b_{5}\right)^{\frac{1}{2}}\right] .
\end{aligned}
$$

For each $n$, denote the values of $\lambda$ corresponding to the $+\operatorname{sign}$ and the $-\operatorname{sign}$ above by $\lambda_{n_{+}}$and $\lambda_{n-}$ respectively. It follows from (12) that for all $n$

$$
0<\lambda_{(n+1)-}<\lambda_{n-}<\lambda_{n+}<\lambda_{(n+1)+}
$$

and

$$
\begin{equation*}
\lambda_{n+} \sim n^{2} \pi^{2}\left(b_{4} b_{5}\right)^{-1} \rightarrow \infty, \quad \lambda_{n-} \sim b_{2} b_{3}\left(n^{2} \pi^{2}\right)^{-1} \rightarrow+0 \quad \text { as } n \rightarrow \infty \tag{18}
\end{equation*}
$$

Substitution in (13a) shows that all $\lambda_{n \pm}$ thus defined are eigenvalues of (13a) and that the eigenfunctions corresponding to $\lambda_{n+}$ and $\lambda_{n-}$ are $\left\{u_{n+}, v_{n+}\right\}$ and $\left\{u_{n_{-}}, v_{n-}\right\}$ respectively where $v_{n_{ \pm}}(x)=c \sin n \pi x$ and

$$
\begin{align*}
\left(b_{2}-b_{\mathbf{4}} \lambda_{n \pm}\right) u_{n \pm}(x) & =v_{n \pm}^{\prime}(x)-b_{1} v_{n \pm}(x) \\
& =c\left(n \pi \cos n \pi x-b_{1} \sin n \pi x\right) \tag{19}
\end{align*}
$$

The above expressions and (17) show that

$$
\begin{align*}
& \lambda_{n+} \lambda_{n-}=b_{2} b_{3} / b_{4} b_{5}=\left(a_{1} a_{3}-a_{2}^{2}\right) / a_{5} a_{6}  \tag{20}\\
& \left(a_{1}-a_{5} \lambda_{n+}\right) u_{n+}=\left(a_{1}-a_{5} \lambda_{n-}\right) u_{n-} ; v_{n+}=v_{n-} \tag{21}
\end{align*}
$$

Since (12) implies

$$
\left(b_{1}^{2}+b_{2} b_{5}+b_{3} b_{4}+n^{2} \pi^{2}\right)^{2}-4 b_{2} b_{3} b_{4} b_{5}>\left(b_{1}^{2}-b_{2} b_{5}+b_{3} b_{4}+n^{2} \pi^{2}\right)^{2}
$$

and $b_{4} b_{5}>0$, it follows that

$$
\lambda_{n-}<\lambda_{0}<\lambda_{n+}
$$

Although (12a) ensures that $B_{1}$ is positive, putting $a_{5}=a_{6}=1$ in (18) shows that it is not positive definite. The + and -families of eigenvalues and eigenfunctions and the single eigenvalue and eigenfunction in between them are in many ways similar to Cowling's $p$-modes, $g$-modes and $f$-mode mentioned in Section 2. While $v_{0} \equiv 0, u_{0}$ has no zero in [0, 1]. Both $u_{n+}$ and $u_{n-}$ have exactly $n$ zeros in $(0,1)$ and do not vanish at either 0 or 1 while $v_{n+}$ and $v_{n-}$ have exactly $n-1$ zeros in $(0,1)$ but also vanish at both 0 and 1 . The positions of the zeros of $u_{n \pm}$ depend on $b_{1}$ but they are separated by those of $v_{n \pm}$ (which is independent of the $b_{i}$ ).

Direct calculation shows that if the Ritz method is applied to (13a) with the space of allowable functions given by

$$
\begin{equation*}
u(x)=c_{1} \sin n \pi x+c_{2} \cos n \pi x, v(x)=c_{3} \sin n \pi x \tag{22}
\end{equation*}
$$

where the $c_{i}$ are the variational parameters, then the Ritz approximate eigenvalues will be exactly $\lambda_{n-}, \lambda_{n+}$ and $\lambda_{0}$, and that the Ritz approximate eigenfunction corresponding to $\lambda_{n_{ \pm}}$will be exactly $\left\{u_{n \pm}, v_{n \pm}\right\}$. The boundary
conditions and the fact that the $a_{i}$ are constant make (22) quite a logical choice. However it can hardly be expected that the results yielded by this choice will be representative of the accuracy obtainable by the Ritz method in general. (See Theorems $\mathbf{l}$ and 3.)

Polynomials are a favourite choice of coordinate functions with the Ritz method and were one of the types of coordinate functions used in [3]. Since the eigenfunctions of (13a) are not polynomials, solutions of (13a) obtained using polynomial coordinate functions should give a truer picture of the reliability of the Ritz method for this problem than (22). It would then seem useful to consider the application of the Ritz method to (13a) using polynomial coordinate functions. The simplest choice of polynomial coordinate functions satisfying (11) is given by (15) where in this case

$$
\begin{equation*}
\xi_{i}(x)=x^{i-1}, \quad \eta_{i}(x)=x^{i}(1-x) . \tag{23}
\end{equation*}
$$

Theorem 4. If $n-m \geqq 2$ when the Ritz method is applied to (13a) using (15) and (23) then, apart from the multiplicity of $\lambda_{0}$, the Ritz approximate eigenvalues and eigenfunctions will be independent of $n$, depending on $m$ only. In this case the Ritz approximate eigenvalues other than $\lambda_{0}$ will occur in pairs satisfying (20) and the corresponding Ritz approximate eigenfunctions will satisfy (21).

Proof. With this choice of coordinate functions, direct calculation shows that $P$ and $Q$ in the matrix equation (7) obtained by the Ritz method are of the form (16) where on this occasion $P_{k}=\left(p_{k i j}\right), k=1, \cdots, 4$ and

$$
\begin{aligned}
& p_{1 i j}=\frac{1}{i+j-1}, \\
& p_{2 i j}=\frac{-a_{1} a_{4} j}{i+j-1}+\frac{a_{1} a_{4}(j+1)+a_{2}}{i+j}-\frac{a_{2}}{i+j+1}, \\
& p_{3 i j}=\frac{2 a_{1} a_{4}^{2} i j}{(i+j-1)(i+j)(i+j+1)}+\frac{2 a_{3}}{(i+j+1)(i+j+2)(i+j+3)}, \\
& p_{4 i j}=\frac{2 a_{6}}{(i+j+1)(i+j+2)(i+j+3)} .
\end{aligned}
$$

Clearly

$$
p_{2 i j}=-a_{1} a_{4} j p_{1 i j}+\left(a_{1} a_{4}(j+1)+a_{2}\right) p_{1, i, j+1}-a_{2} p_{1, i, j+2}
$$

so that, if $\Lambda \neq a_{1} / a_{5}$, column operations on ( $P-\Lambda Q$ ) can reduce the first $n$ elements of the last $m$ columns to zero. These column operations are equivalent to post multiplication by the matrix $R(\Lambda)$ given by

$$
R(\Lambda)=\left(\begin{array}{ll}
I_{n} & R_{1}(\Lambda) \\
0_{m n} & I_{m}
\end{array}\right)
$$

where $\left(a_{1}-a_{5} \Lambda\right) R_{1}(\Lambda)=R_{2}=\left(r_{i j}\right)$ and, for $1 \leqq i \leqq m, \quad r_{i i}=i a_{1} a_{4}$, $r_{i+1, i}=-(i+1) a_{1} a_{4}-a_{2}, r_{i+2, i}=a_{2}$, and $r_{i j}=0$ for $|i-j-1|>1$. Now

$$
(P-\Lambda Q) R(\Lambda)=\left(\begin{array}{ll}
\left(a_{1}-a_{5} \Lambda\right) P_{1} & 0_{n m} \\
P_{2}^{\tau} & P_{5}(\Lambda)
\end{array}\right)
$$

where $\left(a_{1}-a_{5} \Lambda\right) P_{5}(\Lambda)=P_{6}(\Lambda)=\left(p_{6 i j}(\Lambda)\right)$ and

$$
\begin{align*}
p_{6 i j}(\Lambda)= & \frac{-2 \Lambda i j a_{1} a_{4}^{2} a_{5}}{(i+j-1)(i+j)(i+j+1)} \\
& +\frac{2\left[a_{1} a_{3}-a_{2}^{2}-\left(a_{1} a_{6}+a_{3} a_{5}\right) \Lambda+a_{5} a_{6} \Lambda^{2}\right]}{(i+j+1)(i+j+2)(i+j+3)} . \tag{24}
\end{align*}
$$

Since $P_{1}$ is the Gram matrix of a linearly independent set of functions it is non-singular (in fact positive definite) and hence the characteristic equation $\operatorname{det}(P-\Lambda Q)=0$, which determines the Ritz approximate eigenvalues, may be written

$$
\begin{equation*}
\left(a_{1}-a_{5} \Lambda\right)^{n-m} \operatorname{det} P_{6}(\Lambda)=0 \tag{25}
\end{equation*}
$$

The left hand side of (25) is identically equal to $(\operatorname{det}(P-\Lambda Q)) / \operatorname{det} P_{1}$ for $\Lambda \neq a_{1} / a_{5}$ and hence, by continuity, for all $\Lambda$. The factor $\left(a_{1}-a_{5} \Lambda\right)^{n-m}$ gives the roots predicted by Theorem 3 and det $P_{6}(\Lambda)$, a polynomial of degree $2 m$ and independent of $n$, gives the others. Clearly $P_{6}\left(\left(a_{1} a_{3}-a_{2}^{2}\right) / \Lambda a_{5} a_{6}\right)$ is a scalar multiple of $P_{6}(\Lambda)$, so that $\left(a_{1} a_{3}-a_{2}^{2}\right) / \Lambda a_{5} a_{6}$ is a zero of det $P_{6}$ if and only if $\Lambda$ is also. Hence the Ritz approximate eigenvalues occur in pairs satisfying (20).

The equation $(P-\Lambda Q) x=0$ may be thought of as

$$
\begin{equation*}
(P-\Lambda Q) R(\Lambda)\left(R^{-1}(\Lambda) x\right)=0 \tag{26}
\end{equation*}
$$

Let $x_{1}, x_{2}$ be column vectors representing the first $n$ and the last $m$ elements respectively of $x$ and let $x_{3}, x_{4}$ be the corresponding parts of $R^{-1}(\Lambda) x$. Then, for $\Lambda \neq a_{1} / a_{5}$, (26) may be written

$$
\begin{equation*}
P_{1} x_{3}=0, \quad P_{2}^{\tau} x_{3}+P_{5}(\Lambda) x_{4}=0 \tag{27}
\end{equation*}
$$

For each eigenvalue $\Lambda$ of (25) this has a non-zero solution and since $P_{1}$ is nonsingular, this solution must have $x_{3}=0$. From the form of $R$ it is obvious that

$$
R^{-1}(\Lambda)=\left(\begin{array}{cc}
I_{n} & -R_{1}(\Lambda) \\
0_{m n} & I_{m}
\end{array}\right)
$$

so that $x_{2}=x_{4}$ and

$$
\begin{equation*}
x_{1}-R_{1}(\Lambda) x_{2}=x_{3}=0 \tag{28}
\end{equation*}
$$

If $n-m>2$ the $n-m-2$ bottom rows of $R_{1}$, and hence the last $n-m-2$ elements of $x_{1}$ are zero. Also $x_{4}$ and the first $m+2$ rows of $R_{1}$, and hence the first $m+2$ elements of $x_{1}$, are independent of $n$. Hence the Ritz approximate eigenfunctions depend only on $m$. Since $R_{2}=\left(a_{1}-a_{5} \Lambda\right) R_{1}(\Lambda)$ is constant, and since (24) and (27) show that the same $x_{4}$ is obtained by the eigenvalues $\Lambda$ and $\left(a_{1} a_{3}-a_{2}^{2}\right) / \Lambda a_{5} a_{6}$ it follows from (15), (23), (28) and the definition of $x_{1}$ and $x_{2}$ that the corresponding Ritz approximate eigenfunctions satisfy (21). This completes the proof.

Theorem 4 shows that for each $m$, the Ritz approximate eigenvalues in this case will consist of $n-m$ eigenvalues equal to $\lambda_{0}$ together with $m$ eigenvalues greater than or equal to $\Lambda_{0}=\left(b_{2} b_{3} / b_{4} b_{5}\right)^{\frac{1}{2}}$ and $m$ less than or equal to $\Lambda_{0}$. That is the Ritz approximate eigenvalues obtained in this case may be labelled $\lambda_{0}(n-m$ times $), \Lambda_{1 \pm}(m), \Lambda_{2 \pm}(m), \cdots, \Lambda_{m \pm}(m)$ where

$$
A_{m-}(m) \leqq \cdots \leqq \Lambda_{2-}(m) \leqq \Lambda_{1-}(m) \leqq \Lambda_{0} \leqq \Lambda_{1+-}(m) \leqq \cdots \leqq \Lambda_{m+}(m)
$$

Denote the Ritz approximate eigenfunction corresponding to $\Lambda_{k \pm}(m)$ by $\left\{U_{k \pm}(m), V_{k \pm}(m)\right\}$.

Theorem 4 shows that $\Lambda_{k-}(m)$ and $\left\{U_{k-}(m), V_{k-}(m)\right\}$ will be good approximations of $\lambda_{k-}$ and $\left\{u_{k-}, v_{k-}\right\}$ if and only if $\Lambda_{k+}(m)$ and $\left\{U_{k+}(m), V_{k+}(m)\right\}$ are good approximations of $\lambda_{k+}$ and $\left\{u_{k+}, v_{k+}\right\}$. This case contrasts with those results in [3] where $p$-modes were detected but not $g$-modes. However Theorem 4 gives no indication of whether $\Lambda_{k \pm}(m)$ and $\left\{U_{k \pm}(m), V_{k \pm}(m)\right\}$ will in fact be approximations of $\lambda_{k \pm}$ and $\left\{u_{k \pm}, v_{k \pm}\right\}$ in any meaningful sense, or whether they will be spurious like some of the modes discovered in [3]. To settle this question it is necessary to consider whether $\lim _{m \rightarrow \infty} \Lambda_{k \pm}(m)$ and $\lim _{m \rightarrow \infty}\left\{U_{k_{ \pm}}(m), V_{k_{ \pm}}(m)\right\}$ exist, and if so whether these limits are equal, or nearly equal, to $\lambda_{k \pm}$ and $\left\{u_{k \pm}, v_{k \pm}\right\}$, and how rapidly the limits are approached.

The methods normally used to discuss the convergence of the Ritz method (as in [1]) are not immediately applicable here as they use the fact that in standard problems the $k$ th eigenvalue of ( 1 ) is the minimum of $\left(L_{1} w, w\right) /\left(L_{2} w, w\right)$ when $w$ is subject to the restriction that it is orthogonal to the eigenfunctions corresponding to the $k-1$ lowest eigenvalues. In the present problem not only is there no minimum (or maximum) eigenvalue but also there are functions (not eigenfunctions) which are orthogonal to $\left\{u_{0}, v_{0}\right\}$ and yet give the same value of the ratio

$$
\left(B_{1}\{u, v\},\{u, v\}\right) /\left(B_{2}\{u, v\},\{u, v\}\right)
$$

as $\left\{u_{0}, v_{0}\right\}$. An example is

$$
\{u(x), v(x)\}=\left\{\cos (2 \pi n x+c) \exp b_{1} x, 0\right\}
$$

where $c$ is a constant and $n$ an integer.

It is planned to consider the question of convergence in this and related problems in a later paper. In the meantime it may be noted that $\left\{U_{k \pm}(m), V_{k \pm}(m)\right\}$ resembles $\left\{u_{k \pm}, v_{k \pm}\right\}$ in several important ways (Theorem 5 and corollaries) and that for small $m, \Lambda_{k \pm}(m)$ and $\left\{U_{k \pm}(m), V_{k \pm}(m)\right\}$ are in fact good approximations of $\lambda_{k \pm}$ and $\left\{u_{k \pm}, v_{k \pm}\right\}$ (Theorem 6 and the succeeding paragraphs).

Define $U_{k m+}, U_{k m-}$ and $V_{k m}$ to be the extensions of $U_{k+}(m), U_{k-}(m)$ and $V_{k+}(m)\left(=V_{k--}(m)\right)$ to the whole real line.

Theorem 5. (i) $\left(b_{2}-b_{4} \Lambda_{k \pm}(m)\right) U_{k m \pm}=V_{k m}^{\prime}-b_{1} V_{k m}$.
(ii) If $\Lambda_{i+}(m)>\Lambda_{j+}(m)$, then

$$
\int_{0}^{1} V_{i m} V_{j m}=\int_{0}^{1} U_{i m \pm} U_{i m \pm}=0
$$

Proof. (i) If

$$
V_{k m}(x)=x(1-x) \sum_{i=0}^{m-1} c_{i} x^{i}
$$

(28) shows that

$$
\begin{aligned}
\left(b_{2}-b_{4} \Lambda_{k \pm}(m)\right) U_{k m \pm}(x) & =\sum_{i=1}^{m}\left[i c_{i-1} x^{i-1}-\left(b_{1}+i+1\right) c_{i-1} x^{i}+b_{1} c_{i-1} x^{i+1}\right] \\
& =V_{k m}^{\prime}(x)-b_{1} V_{k m}(x)
\end{aligned}
$$

(ii) Since $b_{4} \neq 0$ and $\Lambda_{i+}(m)>\Lambda_{j+}(m) \geqq \Lambda_{j_{-}}(m)>\Lambda_{i-}(m)$ it follows that $b_{2}-b_{4} \Lambda_{i+}(m) \neq b_{2}-b_{4} \Lambda_{i-}(m)$. By Theorem 2,

$$
\int_{0}^{1}\left(a_{5} U_{i m+} U_{j m+}+a_{6} V_{i m} V_{j m}\right)=0=\int_{0}^{1}\left(a_{5} U_{i m-} U_{j m-}+a_{6} V_{i m} V_{j m}\right)
$$

Combining these relations, using

$$
\left(b_{2}-b_{4} \Lambda_{i+}(m)\right) U_{i m+}=\left(b_{2}-b_{4} \Lambda_{i-}(m)\right) U_{i m-}
$$

gives the result.
COROLLARY 1. If the zeros of $V_{k m}$ are all real, then between each pair of consecutive distinct zeros of $V_{k m}$ there is exactly one real zero of $U_{k m \pm}$.

Proof. Since $V_{k m}$ is a polynomial with only real zeros there is exactly one real zero of $V_{k m}^{\prime}$ between each pair of consecutive distinct zeros of $V_{k m}$. Hence by (i) the sign of $U_{k m \pm}$ immediately after a zero of $V_{k m}$ is opposite to its sign immediately before the next zero of $V_{k m}$. Hence there is an odd number of real zeros of $U_{k m \pm}$ between any two consecutive distinct zeros of $V_{k m}$. Moreover at a zero of multiplicity $r$ of $V_{k m}, U_{k m \pm}$ has a zero of multiplicity exactly $r-1$. Since $U_{k m \pm}$ has the same number of zeros as $V_{k m}$ and the zeros of $V_{k m}$ are all real the result follows.

Corollary 2. If $\Lambda_{i+}(m)>\Lambda_{j+}(m)$ then $\int_{0}^{1} V_{i m}^{\prime} V_{j m}^{\prime}=0$.

Proof. Since $\int_{0}^{1}\left(V_{i m}^{\prime} V_{j m}+V_{i m} V_{j m}^{\prime}\right)=0$, the result follows by combining results (i) and (ii) of the Theorem.

Theorem 6. For all non-zero constants $a_{i}$ satisfying (12),

$$
\begin{aligned}
& \lambda_{1+}<\Lambda_{1+}(3)<\Lambda_{1+}(2)=\Lambda_{1+}(1)<1.014 \lambda_{1+}, \\
& \lambda_{2+}<\Lambda_{2+}(3)=\Lambda_{2+}(2)<1.064 \lambda_{2+}, \\
& \lambda_{3+}<\Lambda_{3+}(3)<1.15 \lambda_{3+} \text { and } \Lambda_{1+}(3)<1.000015 \lambda_{1+} .
\end{aligned}
$$

Proof. Define $f(\lambda)$ by

$$
f(\lambda)=\left(a_{1} a_{3}-a_{2}^{2}-\left(a_{1} a_{6}+a_{3} a_{5}\right) \lambda+a_{5} a_{6} \lambda^{2}\right) / a_{1} a_{4}^{2} a_{5} .
$$

Equation (17), which determines the eigenvalues $\lambda_{k \pm}$ may be written $f\left(\lambda_{k \pm}\right)=k^{2} \pi^{2} \lambda_{k \pm}$. Direct calculation shows that the equation $\operatorname{det} P_{6}(\lambda)=0$, which determines the $\Lambda_{k \pm}(m)$, may be written as $f(\lambda)-10 \lambda=0$ when $m=1$, as $(f(\lambda)-10 \lambda)(f(\lambda)-42 \lambda)=0$ when $m=2$, and as

$$
(f(\lambda)-(56-4 \sqrt{ }(133)) \lambda)(f(\lambda)-42 \lambda)(f(\lambda)-(56+4 \sqrt{ }(133)) \lambda)=0
$$

when $m=3$. Since
and

$$
\begin{aligned}
& 0<\pi^{2}<56-4 \sqrt{ }(133)<10<1.014 \pi^{2}, \\
& 4 \pi^{2}<42<1.064\left(4 \pi^{2}\right), \\
& 9 \pi^{2}<56+4 \sqrt{ }(133)<1.15\left(9 \pi^{2}\right) \\
& 56-4 \sqrt{ }(133)<1.000015 \pi^{2},
\end{aligned}
$$

the result follows from the following lemma.
Lemma. If the $a_{i}$ are non-zero and satisfy (12), and $K>0$ the roots of $f(\lambda)=K \lambda$ are real, distinct and positive. If $K_{2}>K_{1}>0$, then

$$
K_{1} \lambda^{*}\left(K_{1}\right)<K_{1} \lambda^{*}\left(K_{2}\right)<K_{2} \lambda^{*}\left(K_{1}\right)
$$

where $\lambda^{*}(K)$ is the greater root of $f(\lambda)=K \lambda$.
Proof. If the constants $c_{1}, c_{2}, c_{3}, c_{4}, c_{2}^{2}-4 c_{1} c_{4}$ are all strictly positive, the roots of $c_{1} y^{2}-\left(c_{2}+c_{3}\right) y+c_{4}=0$ are real distinct and positive. Let $y\left(c_{3}\right)$ be the greater of these roots. Direct solution shows that if $\delta>0$, then $y\left(c_{3}+\delta\right)>y\left(c_{3}\right)$. Also

$$
\begin{aligned}
& y\left(c_{3}+\delta\right)-y\left(c_{3}\right) \\
& \quad=\frac{1}{2 c_{1}}\left[\delta+\left(c_{3}+\delta\right)\left(1+\frac{2 c_{2}}{c_{3}+\delta}+\frac{c_{2}^{2}-4 c_{1} c_{4}}{\left(c_{3}+\delta\right)^{2}}\right)^{\frac{1}{2}}-c_{3}\left(1+\frac{2 c_{2}}{c_{3}}+\frac{c_{2}^{2}-4 c_{1} c_{4}}{c_{3}^{2}}\right)^{\frac{1}{2}}\right] \\
& \quad<\frac{\delta}{2 c_{1}}\left[1+\left(1+\frac{2 c_{2}}{c_{3}}+\frac{c_{2}^{2}-4 c_{1} c_{4}}{c_{3}^{2}}\right)^{\frac{1}{2}}\right]<\frac{\delta}{c_{3}} y\left(c_{3}\right) .
\end{aligned}
$$

Since

$$
a_{5} a_{6}, a_{1} a_{6}+a_{3} a_{5}, a_{1} a_{4}^{2} a_{5}, a_{1} a_{3}-a_{2}^{2} \text { and }\left(a_{1} a_{6}+a_{3} a_{5}\right)^{2}-4 a_{5} a_{6}\left(a_{1} a_{3}-a_{2}^{2}\right)
$$

are all strictly positive, the result follows.
The bounds given by Theorem 6 are uniform bounds. Finer bounds can be found for specific values of the $a_{i}$. Since, by Theorem 4,

$$
\Lambda_{k+}(m) \Lambda_{k-}(m)=\lambda_{k+} \lambda_{k-}
$$

bounds for the $\Lambda_{k-}(m)$ follow immediately from Theorem 6. These initial results exhibit a convergence similar to that encountered in standard problems. Note that for $m \leqq 3, \Lambda_{k-}(m)<\lambda_{k-}<\lambda_{k+}<\Lambda_{k+}(m)$. This is similar to the result in [14].

Since the value of $f(\lambda) / \lambda$ at the zeros of $\operatorname{det} P_{6}(\lambda)$ in Theorem 4 is independent of the $a_{i}$, it follows that the eigenvectors $x_{4}$ of $P_{6}(\lambda)$, and hence the $V_{k \pm}(m)$, are also independent of the $a_{i}$, as are also the $v_{k \pm}$. Direct calculation shows that $x_{4}$ corresponding to $\Lambda_{2 \pm}(3)$ is $\{1,-2,0\}$. Hence $V_{2 \pm}(3)$, like $v_{2 \pm}$, has a zero at $\frac{1}{2}$ and no other zeros in $(0,1)$. Moreover $V_{2 \pm}(2)=V_{2 \pm}(3)$. Also with $m=3, x_{4}$ corresponding to the eigenvalue given by $f(\lambda)=(56 \pm 4 \sqrt{ }(133)) \lambda$ is $\{5 \pm 3 \sqrt{ }(133) / 7,-23 \mp 2 \sqrt{ }(133), 23 \pm 2 \sqrt{ }(133)\}$. This shows that $V_{3 \pm}(3)$ and $V_{1 \pm}(3)$ have the same number of zeros in ( 0,1 ) as $v_{3 \pm}$ and $v_{1 \pm}$ respectively. The zeros of $V_{3 \pm}(3)$ in $(0,1)$ are the roots of $(23+2 \sqrt{ }(133))\left(y^{2}-y\right)+5+3 \sqrt{ }(133) / 7=0$, which have the approximate values 0.32 and 0.68 . The zeros of $v_{3 \pm}$ in $(0,1)$ are $\frac{1}{3}, \frac{2}{3}$. Similarly direct calculation shows that $V_{1 \pm}(1)=V_{1 \pm}(2)$ and that these functions, like $v_{1 \pm}$ and $V_{1 \pm}(3)$, have no zeros in $(0,1)$.

It follows from Theorem 5 and Corollary 1 that for $m \leqq 3$, the zeros of $U_{k \pm}(m)$ follow a pattern similar to those of $u_{k \pm}$. Another point of resemblance between the two sets of vectors easily shown by direct calculation for $m \leqq 3$ is that $\left|U_{k m \pm}(0)\right|=\left|U_{k m \pm}(1)\right|$.

## 5. Concluding remarks

The results of the last section indicate that there exist cases not covered by the classical theory in which useful results may be obtained by the Ritz method. It remains to find general criteria of the reliability of the method in particular circumstances. Although, as noted in section 3, the results of [16] will not be representative of the effectiveness of application of the Ritz method to (2) for general models, and as noted in section 2, even [14] and [15] leave the efficacy of the method in doubt, it does seem that the results of [3] may not be typical.

Rounding errors, which affect all computer calculations such as those in [3], [14] and [15], were not considered in sections 3 and 4. The matrix
$P_{1}$ arising in Theorem 4 is the ill-conditioned Hilbert matrix [17] and some moderately ill-conditioned matrices were obtained in the calculations described in [3], as will always happen if the coordinate functions are not strongly minimal [18]. It is possible that the effect of rounding errors on the solution of (2) could be different when different stellar models are used, especially as (6) will approximate the true eigenfunctions of (2) more closely for some stellar models than for others. It is planned to consider the question of rounding errors in a later paper.

In section 2 it was suggested that the fact that the model used in [3], unlike those used in [14], [15] and [16], contains a core throughout which $A \equiv 0$, might be partly responsible for the difference in results. Comparison of equations (2) and (13a) shows that the condition $a_{3}>a_{2}^{2} / a_{1}$ in (12) may be regarded as analogous to the condition $A d p \mid d r>0$ with (2), while the conditions $a_{1}>0, a_{5}>0, a_{6}>0$ correspond to the conditions $\Gamma_{1} p>0$ and $\rho>0$. Since $\Gamma_{1} p, \rho$ and $-d p \mid d r$ (unlike $-A$ ) are strictly positive, except perhaps at the boundaries, in all the stellar models considered here (and indeed in any realistic stellar model), it would seem of interest to see how many of the results of the last section continue to hold when the requirement $a_{3}>a_{2}^{2} / a_{1}$ is relaxed, while the requirement that $a_{1}, a_{5}$ and $a_{6}$ be positive and all the $a_{i}$ non-zero is still imposed.

If $a_{3}<a_{2}^{2} / a_{1}$, (14) and (17) still hold and in this case the $\lambda_{n-}$ eigenvalues of (13a) are all negative. (Very similar results have been obtained for $I A$ and $I B$, considered in section 2, for stellar models with $A>0$ throughout [5, p. 514]. More complicated results have been obtained [8] when $A$ changes sign.) Theorems 3, 4 and 5 also still hold, and provided $a_{1} a_{6}+a_{3} a_{5}>0$, Theorem 6 will hold also.

When $a_{1} a_{3}=a_{2}^{2}$, all functions $\{u, v\}$ in the domain of $B_{1}$, satisfying $u=a_{4} v^{\prime}-\left(a_{2} / a_{1}\right) v$ are eigenfunctions of (13a). The corresponding eigenvalue, which has infinite multiplicity, is zero. For the non-zero eigenvalues, many of the results of section 4, including equations (14), (17), (19), (25) and (28) and Theorems 3, 5(i) and 6, still hold.

The analogy between (13a) with $a_{1} a_{3}=a_{2}^{2}$ and $I B$ with the convective model discussed in section 2 is manifest. With the stellar model used [3], all sufficiently differentiable functions satisfying the boundary conditions with $\psi$ vanishing throughout the envelope (where $A \neq 0$ ) but not throughout the core, and with $\chi^{\prime}=\psi^{\prime}+\left(p^{\prime} \mid \Gamma_{1} p\right) \psi$ throughout the star, are eigenfunctions of (2) and hence of $I B$. The corresponding eigenvalue is zero, and again has infinite multiplicity. Results of Lebovitz [19] show that these functions are also eigenfunctions of $I A$ for this model. Results of [11] (where the zero eigenvalues are not mentioned) show that with this model $I B$ has, as well as the zero eigenvalues, the normal $p$-modes, $f$-mode and $g$-modes, although the eigenvalues obtained there for $g$-modes were small
compared with those obtained by others for models where $A<0$ throughout.
The relatively complicated spectrum when $A$ vanishes in parts of the star could help explain why the method was less successful in [3] than in [14] and [15], although in practice rounding errors have introduced some (very small) non-zero values of $A$ into the core of the model used in [3]. This should make the model resemble a real star still more closely [8]. However the spurious eigenfunctions bear no resemblance to the true eigenfunctions corresponding to zero eigenvalues.

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