ROOT VECTORS FOR GEOMETRICALLY SIMPLE TWO-PARAMETER EIGENVALUES

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ABSTRACT. A class of two-parameter eigenvalue problems involving generally nonselfadjoint and unbounded operators is studied. A basis for the root subspace at a geometrically simple eigenvalue of Fredholm type is computed in terms of the underlying two-parameter system. Comparison with Faierman's work on two-parameter boundary value problems of Sturm-Liouville type is given as an application.

1. Introduction

We consider a pair of two-parameter equations of the form

$$(1.1) W_i(\lambda)x_i = 0 \neq x_i,$$

where

$$(1.2) W_i(\lambda) = A_{i0}\lambda_0 + A_{i1}\lambda_1 + A_{i2},$$

 A_{ij} (j=1,2) are bounded linear operators acting on Hilbert spaces H_i (i=1,2) over the complex numbers, A_{i0} are closed densely defined with domain $\mathcal{D}(A_{i0}) \subseteq H_i$, and $\boldsymbol{\lambda} = (\lambda_0, \lambda_1) \in \mathbb{C}^2$. Evidently, $W_i(\boldsymbol{\lambda})$ is a linear operator with domain $\mathcal{D}(A_{i0})$ for all such parameters $\boldsymbol{\lambda}$.

The spectral theory of such systems has been developed from various viewpoints; see [1, 8, 22, 24, 26] for books on the subject. In most of the literature, the A_{ij} are selfadjoint and satisfy definiteness conditions guaranteeing that the eigenvalues are simple in certain senses. A standard selfadjoint example is when (1.1) are Sturm-Liouville equations. One application of these ideas is to completeness and expansion theory of elliptic boundary value problems via separation of variables. For work on nonselfadjoint problems we refer to [5, 9, 11, 14, 15].

To introduce our analysis, we make the following constructions. The operators A_{ij} , j=1,2, induce operators A_{ij}^{\dagger} on the Hilbert space tensor product $H=H_1\otimes H_2$ by means of

$$A_{1j}^{\dagger}\left(x_{1}\otimes x_{2}\right)=A_{1j}x_{1}\otimes x_{2}$$
 and $A_{2j}^{\dagger}\left(x_{1}\otimes x_{2}\right)=x_{1}\otimes A_{2j}x_{2}$

on decomposable tensors, extended by linearity and continuity to the whole of H. Similarly, A_{i0} induces an operator A_{i0}^{\dagger} with domain $\mathcal{D}\left(A_{i0}^{\dagger}\right) \subset H$ (see [24, §2.3]).

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We denote by \mathcal{D} the intersection $\mathcal{D}\left(A_{10}^{\dagger}\right) \cap \mathcal{D}\left(A_{20}^{\dagger}\right)$, which is a dense subspace of H. We define the operator

(1.3)
$$\Delta_0 = \begin{vmatrix} A_{11}^{\dagger} & A_{12}^{\dagger} \\ A_{21}^{\dagger} & A_{22}^{\dagger} \end{vmatrix} = A_{11} \otimes A_{22} - A_{12} \otimes A_{21}$$

on H and operators

(1.4)
$$\Delta_1 = -A_{22}^{\dagger} A_{10}^{\dagger} + A_{12}^{\dagger} A_{20}^{\dagger} \text{ and } \Delta_2 = -A_{11}^{\dagger} A_{20}^{\dagger} + A_{21}^{\dagger} A_{10}^{\dagger}$$

on \mathcal{D} . In §2 we shall make assumptions guaranteeing the existence and commutativity of the operators $\Gamma_j = \Delta_2^{-1} \Delta_j$, j = 0, 1, on \mathcal{D} .

The nullspace of an operator A is denoted by $\mathcal{N}(A)$. A pair $\boldsymbol{\lambda} = (\lambda_0, \lambda_1) \in \mathbb{C}^2$ is called an eigenvalue of the pair of commuting linear transformations $\{\Gamma_0, \Gamma_1\}$ if

(1.5)
$$\mathcal{R}_1 = \mathcal{R}_1(\lambda) = \mathcal{N}(\Gamma_0 - \lambda_0 I) \cap \mathcal{N}(\Gamma_1 - \lambda_1 I) \neq \{0\}.$$

The subspace \mathcal{R}_1 is called the *geometric eigenspace* corresponding to $\boldsymbol{\lambda}$. The set of all the eigenvalues is called the *point spectrum* of $\{\Gamma_0, \Gamma_1\}$. We say that $\boldsymbol{\lambda}$ is a *geometrically simple eigenvalue* of $\{\Gamma_0, \Gamma_1\}$ if dim $\mathcal{R}_1 = 1$. Under assumptions to be specified in §2, the eigenvalues $\boldsymbol{\lambda}$ of (1.2) and that of its associated linear transformations Γ_0, Γ_1 coincide and we have

(1.6)
$$\mathcal{R}_{1}(\boldsymbol{\lambda}) = \mathcal{N}\left(W_{1}\left(\boldsymbol{\lambda}\right)\right) \otimes \mathcal{N}\left(W_{2}\left(\boldsymbol{\lambda}\right)\right).$$

(See [5, §4].) We shall simply refer to such λ as eigenvalues from now on. We define the m-th root subspace

(1.7)
$$\mathcal{R}_{m} = \mathcal{R}_{m}(\boldsymbol{\lambda}) = \bigcap_{k=0}^{m} \mathcal{N}\left[(\Gamma_{0} - \lambda_{0}I)^{k} (\Gamma_{1} - \lambda_{1}I)^{m-k} \right],$$

for $m \geq 2$. The least integer M, if it exists, such that $\mathcal{R}_M = \mathcal{R}_{M+1}$ is called the ascent of $\boldsymbol{\lambda}$. We call the subspace

$$\mathcal{R} = \mathcal{R}(\boldsymbol{\lambda}) = igcup_{m=1}^{\infty} \mathcal{R}_m$$

the root subspace of $\mathbf{W}=(W_1,W_2)$ at $\boldsymbol{\lambda}$. Note that each subspace \mathcal{R}_m is an invariant subpace for both Γ_0 and Γ_1 . If the ascent of $\boldsymbol{\lambda}$ is equal to 1, then $\boldsymbol{\lambda}$ is called *semisimple*. If $\boldsymbol{\lambda}$ is a geometrically simple eigenvalue and its ascent is equal to 1, then $\mathcal{R}=\mathcal{R}_1$ is one-dimensional and $\boldsymbol{\lambda}$ is called *algebraically simple*.

The elements of \mathcal{R} are called *root vectors* and it is our purpose to describe bases of root vectors for \mathcal{R} in terms of the A_{ij} . This problem was raised by Atkinson [1] and has been studied in certain definite selfadjoint cases by many authors. See [26] for a review of eigenvector completeness under various conditions including positive definiteness of Δ_2 . In these cases the eigenvalues are real (i.e. $\lambda \in \mathbb{R}^2$) and always dim $\mathcal{R} = \dim \mathcal{R}_1$. Thus if an eigenvalue is geometrically simple, then it is also algebraically simple. Moreover, (1.6) shows that a basis for $\mathcal{R} = \mathcal{R}_1$ can be constructed from the eigenvectors $x_1 \otimes x_2$ for the Γ_j , x_i as in (1.1).

Without definiteness, however, this may fail and various authors have tackled Atkinson's question. We cite [3] on real eigenvalues for selfadjoint problems and [8] on general eigenvalues for Sturm-Liouville (and therefore geometrically simple) problems. Faierman gives bases of \mathcal{R}_l for l=2 and 3 and gives a conjecture for general l. In an earlier paper [5] we gave a basis for \mathcal{R}_2 for general nonselfadjoint

(including geometrically nonsimple) problems. A coalgebraic approach to this problem in [11] describes root subspaces \mathcal{R}_l for all l in terms of A_{ij} , but it is difficult to use it to construct bases explicitly. Here we shall give a general construction for a basis of \mathcal{R} for geometrically simple eigenvalues of problems that are not necessarily selfadjoint. We shall relate our construction in the Sturm-Liouville case to Faierman's work in §5.

Our plan is as follows. Our main assumptions are discussed in $\S 2$, followed by the construction of a basis for the root subspace in the so-called nonderogatory case. In $\S 3$, we lay foundations for the more difficult derogatory case. The main result (Theorem 4.1) is proved by induction on m (see (1.7)) in $\S 4$. $\S 5$ contains the specialization to the Sturm-Liouville case.

2. Preliminaries

We begin with our assumptions on the A_{ij} . We make two regularity assumptions and a solvability assumption.

Assumption I. There exists $\alpha \in \mathbb{C}$ such that the operator $\Delta_2' = \Delta_2 + \alpha \Delta_0 : \mathcal{D} \to H$ has a bounded inverse.

This assumption holds in many interesting cases. We refer to [5, p. 231] for a discussion.

To simplify the notation we assume, unless stated otherwise, that the operator Δ_2 has a bounded inverse. This form of Assumption I can be obtained by a shift in parameters.

Let $\mathcal{A}_{ij}^{\dagger}$ denote the restriction of A_{ij}^{\dagger} to \mathcal{D} . The array

$$\mathcal{A} = \left[egin{array}{ccc} \mathcal{A}_{10}^\dagger & \mathcal{A}_{11}^\dagger & \mathcal{A}_{12}^\dagger \ \mathcal{A}_{20}^\dagger & \mathcal{A}_{21}^\dagger & \mathcal{A}_{22}^\dagger \end{array}
ight]$$

then defines a linear map $\mathcal{A}: \mathcal{D}^3 \longrightarrow H^2$. Here H^k is the direct sum of k copies of H. Omitting the j-th column we get a transformation \mathcal{A}_j acting on the (algebraic) direct sum \mathcal{D}^2 for j = 0, 1, 2. Note that $\Delta_j = (-1)^j \det \mathcal{A}_j$. We denote by \mathcal{C}_j the j-th column of \mathcal{A} . Now we state the solvability assumption.

Assumption II. The equation $A_2 \mathbf{y} = C_2 x$ has a solution $\mathbf{y} \in \mathcal{D}^2$ for all $x \in \mathcal{D}$.

Again see [5, p. 236] for a discussion of this assumption.

The linear transformations $\Gamma_j = \Delta_2^{-1} \Delta_j$, (j=0,1) act on \mathcal{D} . These are called the associated linear transformations of the two-parameter system (1.2). Under Assumption II it follows from [5, Thm. 3.2] that the linear transformations Γ_0 and Γ_1 commute on \mathcal{D} , i.e., $\Gamma_0 \Gamma_1 x = \Gamma_1 \Gamma_0 x$ for all $x \in \mathcal{D}$, and that

(2.1)
$$A_{i0}^{\dagger} \Gamma_0 x + A_{i1}^{\dagger} \Gamma_1 x + A_{i2}^{\dagger} x = 0$$

for i=1,2 and all $x\in\mathcal{D}$. Relation (2.1) is a consequence of Assumption II that will be often used in the proofs.

Next we define the notions of eigenvalues and point spectra for a two-parameter system. A point $\lambda \in \mathbb{C}^2$ is called an eigenvalue of the two-parameter system (1.2) if both nullspaces $\mathcal{N}(W_i(\lambda))$, i=1,2, are nontrivial. The set of all eigenvalues is called the point spectrum of (1.2), and it is denoted by $\sigma(\mathbf{W})$. An eigenvalue $\lambda \in \sigma(\mathbf{W})$ is called geometrically simple if dim $\mathcal{N}(W_i(\lambda)) = 1$ for both i. Thus λ is a geometrically simple eigenvalue of (1.2) if and only if it is a geometrically

simple eigenvalue for the pair of associated linear transformations (Γ_0, Γ_1) although perhaps not for either of Γ_j separately.

Assumption III. $\lambda = (\lambda_0, \lambda_1)$ is a geometrically simple eigenvalue of finite ascent M and $W_i(\lambda)$, i = 1, 2, are Fredholm of index 0 [25].

This assumption is satisfied, for example, in several cases arising from boundary value problems, e.g., of Sturm-Liouville type (see §5). Assumption III is in terms of the maps A_{ij} except for the assumption of finite ascent. This is defined in terms of the associated linear transformations. It follows from our results in §4 that an equivalent definition of the ascent is possible in terms of the A_{ij} – it is the least m such that $\mathcal{N}(\mathcal{S}_m) = 0$.

In view of Assumption III we have $\dim \mathcal{N}(W_i(\boldsymbol{\lambda})) = \dim \mathcal{N}(W_i(\boldsymbol{\lambda})^*) = 1$ for both $i, \mathcal{R}_1 \subseteq \mathcal{R}_2 \subseteq \cdots \subseteq \mathcal{R}_M$ and $\mathcal{R}_{M+k} = \mathcal{R}_M$ for all $k \geq 1$.

Finally, we remark that Assumptions I–III of our present setup are stronger than Assumptions I–III of [5] or [21] in the case n=2. Assumption I here and in [5, 21] is the same. Here and in [21] we use the same solvability assumption while in [5] a weaker solvability assumption is used. In [21] it is assumed further that λ is geometrically simple. In both [5, 21] it is assumed that the $W_i(\lambda)$ are Fredholm operators not necessarily of index 0. The assumption of Fredholm index 0 here is made simply for ease of notation.

Next we state two-parameter versions of some results of [5, 18, 19, 20] on root subspaces at geometrically simple eigenvalues of commuting linear maps and of multiparameter systems. These versions are used in what follows.

We write

$$D_m = \dim \mathcal{R}_m \text{ and } d_m = D_m - D_{m-1}$$

for m > 1. Here $D_0 = 0$.

Lemma 2.1. If λ is a geometrically simple eigenvalue, then for each $m \geq 2$ the subspace \mathcal{R}_m is finite-dimensional.

Proof. The argument is similar to the one used in the proof of [5, Lemma 5.2]. It is by induction on m. Assume $m \geq 2$. The subspace \mathcal{R}_m is invariant for both $\Gamma_j - \lambda_j I$, j = 0, 1. The range $\mathcal{R}\left((\Gamma_j - \lambda_j I)|_{\mathcal{R}_m}\right)$ is contained in \mathcal{R}_{m-1} and thus is finite-dimensional because of the inductive assumption. Each kernel $\mathcal{N}\left((\Gamma_j - \lambda_j I)|_{\mathcal{R}_m}\right)$ has finite codimension in \mathcal{R}_m , i.e., the orthogonal complement \mathcal{Q}_j of $\mathcal{N}\left((\Gamma_j - \lambda_j I)|_{\mathcal{R}_m}\right)$ in \mathcal{R}_m is finite-dimensional. Hence the linear span \mathcal{Q} of the \mathcal{Q}_j , j = 0, 1, is finite-dimensional. Then $\mathcal{R}_m = \mathcal{R}_{m-1} + \mathcal{Q}$ also has finite dimension.

Since λ is geometrically simple, (1.6) implies that $d_1 = 1$. The restricted maps of Γ_0 and Γ_1 to \mathcal{R}_m commute, and so it follows that $d_i \leq i$ for all $i \geq 1$. Suppose that

$$(2.2) x_{i1} \in \mathcal{N}\left(W_i\left(\boldsymbol{\lambda}\right)\right), \quad i = 1, 2,$$

are nonzero vectors. Then the vector

$$(2.3) z_1 = x_{11} \otimes x_{21}$$

spans \mathcal{R}_1 . Further, we choose nonzero vectors $y_{i1} \in \mathcal{N}\left(W_i\left(\boldsymbol{\lambda}\right)^*\right)$ (i=1,2) and we write

(2.4)
$$S_2 = \begin{bmatrix} \langle A_{10}x_{11}, y_{11} \rangle & \langle A_{11}x_{11}, y_{11} \rangle \\ \langle A_{20}x_{21}, y_{21} \rangle & \langle A_{21}x_{21}, y_{21} \rangle \end{bmatrix}.$$

Here we use $\langle x,y \rangle$ to denote the scalar product of vectors x,y. It is clear from the context which Hilbert space H_i or H is meant. Now [5, Cor. 6.4] states that $d_2 = \dim \mathcal{N}(\mathcal{S}_2)$. If $d_2 \leq 1$, then we call λ nonderogatory. Otherwise, $d_2 = 2$, and we call λ derogatory. Bases for the root subspaces \mathcal{R}_m corresponding to a nonderogatory eigenvalue of an n-parameter system (for $n \geq 2$) are constructed in [19] in the finite-dimensional setup. For completeness we will give an infinite-dimensional two-parameter version of the main result of [19] (see Theorem 2.2 below).

For $\mathbf{c} = \left[c_1, c_2\right]^T \in \mathbb{C}^2$ and i = 1, 2, we write

$$U_i(\mathbf{c}) = A_{i0}c_1 + A_{i1}c_2.$$

If M = 1, then z_1 spans \mathcal{R}_M . Assume from now on that $M \geq 2$. The following is a two-parameter analogue of the main result of [19]. It gives a basis for the root subspace at a nonderogatory eigenvalue.

Theorem 2.2. Suppose that $\lambda \in \mathbb{C}^2$ is a nonderogatory eigenvalue for a two-parameter system (1.2) and suppose that $2 \leq m \leq M$. Then there exist $\mathbf{c}_2, \mathbf{c}_3, \ldots, \mathbf{c}_m \in \mathbb{C}^2, \mathbf{c}_2 \neq 0$, and $x_{i2}, x_{i3}, \ldots, x_{im} \in \mathcal{D}(A_{i0})$, i = 1, 2, such that

(2.5)
$$\sum_{j=1}^{k-1} U_i(\mathbf{c}_{k+1-j}) x_{ij} + W_i(\lambda) x_{ik} = 0 \quad \text{for } k = 2, 3, \dots, m, \ i = 1, 2.$$

Moreover, for k = 1, 2, ..., m, the vectors

(2.6)
$$z_k = \sum_{j=1}^k x_{1j} \otimes x_{2,k+1-j}$$

are such that

(2.7)
$$(\Gamma_{i-1} - \lambda_{i-1}I) z_k = \sum_{j=1}^{k-1} c_{k+1-j,i} z_j, \quad i = 1, 2,$$

where $\mathbf{c}_k = \begin{bmatrix} c_{k1} & c_{k2} \end{bmatrix}^T$ and $\{z_1, z_2, \dots, z_m\}$ is a basis for the m-th root subspace \mathcal{R}_m .

Proof. The theorem is proved by induction on m. For each m we choose a vector $z_m \in \mathcal{R}_m \backslash \mathcal{R}_{m-1}$ and use Assumptions II and III to prove the existence of vectors $x_{ij} \in \mathcal{D}(A_{i0})$ such that (2.6) and (2.7) hold. Here the arguments follow closely those in the proof of [5, Thm. 6.2]. For the converse, (2.7) follows from (2.5) and (2.6) by a direct calculation (see [19, Lemma 13]).

3. Derogatory eigenvalues

We assume now that λ is a geometrically simple eigenvalue which is derogatory, i.e., $S_2 = 0$ or equivalently $d_2 = 2$. First we recall some results on the structure of nilpotent commutative matrices with one-dimensional joint kernel and $d_2 = 2$ (see [18, 20]).

3.1. Structure of nilpotent commutative matrices. We denote by $\mathcal{R} = \mathcal{R}_M$ the root subspace corresponding to λ and we write $N = \dim \mathcal{R}$. The restricted linear maps $C_i = (\Gamma_i - \lambda_i I)|_{\mathcal{R}}$, i = 0, 1, are commuting and nilpotent. There exists a basis

$$(3.1) \mathcal{B} = \{z_1, z_{21}, z_{22}, z_{31}, \dots, z_{3d_3}, \dots, z_{M1}, z_{M2}, \dots, z_{Md_M}\}\$$

for \mathcal{R} such that for $m = 1, 2, \ldots, M$, the set

$$\mathcal{B}_m = \{z_1, z_{21}, z_{22}, z_{31}, \dots, z_{3d_3}, \dots, z_{m1}, z_{m2}, \dots, z_{md_m}\}\$$

is a basis for \mathcal{R}_m . We call such a basis \mathcal{B} a filtered basis (for \mathcal{R}). In a filtered basis the matrices for $\mathbf{C} = \{C_0, C_1\}$ are simultaneously reduced to a special upper-triangular form and we view them as a cubic array of dimensions $N \times N \times 2$

(3.2)
$$\mathbf{C} = \begin{bmatrix} \mathbf{0} & \mathbf{C}^{12} & \mathbf{C}^{13} & \cdots & \mathbf{C}^{1,M} \\ \mathbf{0} & \mathbf{0} & \mathbf{C}^{23} & \cdots & \mathbf{C}^{2,M} \\ \vdots & \vdots & \ddots & \ddots & \vdots \\ \mathbf{0} & \mathbf{0} & \mathbf{0} & \cdots & \mathbf{C}^{M-1,M} \\ \mathbf{0} & \mathbf{0} & \mathbf{0} & \cdots & \mathbf{0} \end{bmatrix},$$

where

$$\mathbf{C}^{km} = egin{bmatrix} \mathbf{c}_{11}^{km} & \mathbf{c}_{12}^{km} & \cdots & \mathbf{c}_{1,d_m}^{km} \ \mathbf{c}_{21}^{km} & \mathbf{c}_{22}^{km} & \cdots & \mathbf{c}_{2,d_m}^{km} \ dots & dots & dots \ \mathbf{c}_{d_{k},1}^{km} & \mathbf{c}_{d_{k},2}^{km} & \cdots & \mathbf{c}_{d_{k},d_m}^{km} \end{bmatrix}$$

is a cubic array of dimensions $d_k \times d_m \times 2$ and $\mathbf{c}_{gf}^{km} = \begin{bmatrix} c_{1gf}^{km} & c_{2gf}^{km} \end{bmatrix}^T \in \mathbb{C}^2$. The column cross-sections

(3.3)
$$C_f^{km} = \begin{bmatrix} \mathbf{c}_{1f}^{km} & \mathbf{c}_{2f}^{km} & \cdots & \mathbf{c}_{d_k,f}^{km} \end{bmatrix}, \ f = 1, 2, \dots, d_m,$$

of \mathbf{C}^{km} are $2 \times d_k$ matrices. The constants c_{iqf}^{km} are determined by the relations

$$C_{i-1}z_{mf} = \sum_{k=1}^{m-1} \sum_{g=1}^{d_k} c_{igf}^{km} z_{kg}, \quad i = 1, 2,$$

where we write $z_{11} = z_1$. (For further details see [18, 20].)

For $m \geq 4$ we denote by Φ_m the set of multi-indices

$$\left\{ (l_1, l_2, l_3); \ l_j \ge 2, \ \sum_{j=1}^3 l_j \le m+2 \right\},\,$$

and for $\mathbf{l} = (l_1, l_2, l_3) \in \Phi_m$ we denote by χ_1 the set of multi-indices

$$\{(h_1, h_2, h_3); 1 \le h_j \le d_{l_j} \text{ for } j = 1, 2, 3\}.$$

We call a matrix C symmetric if $C = C^T$.

For use in our proofs we state a version of some of results of [18, 20] for pairs of commutative matrices. Since \mathbf{c}_{11}^{12} and \mathbf{c}_{12}^{12} are linearly independent we can choose a basis for \mathcal{R}_2 such that $C_{i-1}z_{2j} = \delta_{ij}z_1$, where δ_{ij} is the Kronecker symbol. By adding appropriate linear combinations of z_{21} and z_{22} to other vectors z_{mf} , $m \geq 3$, we find a basis \mathcal{B} such that the following result holds.

Theorem 3.1. There exists a filtered basis \mathcal{B} for \mathcal{R} and associated with it a set of symmetric matrices

(3.4)
$$C_f^m = \left[C_f^{klm} \right]_{k,l=2}^{m-1}, \ m = 3, 4, \dots, M, \ f = 1, 2, \dots, d_m,$$

where $C_f^{klm} = \left[c_{ghf}^{klm}\right]_{g=1,\,h=1}^{d_k}$ are $d_k \times d_l$ matrices, such that for the corresponding commutative array (3.2) we have:

- (a) $C_1^{12} = \begin{bmatrix} 1 & 0 \end{bmatrix}^T$, $C_2^{12} = \begin{bmatrix} 0 & 1 \end{bmatrix}^T$ and $C_f^{1m} = \begin{bmatrix} 0 & 0 \end{bmatrix}^T$ for m = 3, 4, ..., M and $f = 1, 2, ..., d_m$, (b) $C_f^{2km} = C_f^{km}$ are the column cross-sections (3.3),
- (c) for m = 3, 4, ..., M the matrices $C_f^{2,m-1,m}$, $f = 1, 2, ..., d_m$, are linearly
- (d) $C_f^{klm} = 0 \text{ if } k + l > m + 1,$
- (e) for m = 4, 5, ..., M and $f = 1, 2, ..., d_m$ the relations

$$(3.5) \qquad \sum_{k=l_1+l_2-1}^{m-l_3+1} \sum_{g=1}^{d_k} c_{h_1 h_2 g}^{l_1 l_2 k} c_{h_3 g f}^{l_3 k m} = \sum_{k=l_2+l_3-1}^{m-l_1+1} \sum_{g=1}^{d_k} c_{h_1 h_3 g}^{l_1 l_3 k} c_{h_2 g f}^{l_2 k m}$$

hold for all
$$\mathbf{l} = (l_1, l_2, l_3) \in \Phi_m$$
 and $(h_1, h_2, h_3) \in \chi_{\mathbf{l}}$.

We remark that property (e) in the above theorem together with the symmetry of matrices C_f^m is equivalent to the commutativity of the matrices C_0 and C_1 and that c_{ghf}^{klm} are the structure constants for multiplication in the commutative algebra generated by C_0, C_1 and the identity matrix (see [20]).

In the rest of the section we consider, unless stated otherwise, only filtered bases $\mathcal B$ such that the associated matrices C_m^f satisfy the properties (a)–(e) of Theorem 3.1.

3.2. Bases for \mathcal{R}_m for m=2,3. We denote by \mathcal{D}_i' the (vector space) intersection $\mathcal{D}(A_{i0}) \cap (\mathcal{N}(W_i(\boldsymbol{\lambda})))^{\perp}$, where $(\mathcal{N}(W_i(\boldsymbol{\lambda})))^{\perp}$ is the orthogonal complement of $\mathcal{N}\left(W_i\left(\boldsymbol{\lambda}\right)\right)$ in H_i .

Lemma 3.2. For i = 1, 2, the subspace \mathcal{D}'_i is an infinite-dimensional vector subspace of $\mathcal{D}_i = \mathcal{D}_i(A_{i0})$.

Proof. Since the lattice of vector subspaces is modular [7, p. 13] and since $\mathcal{N}(W_i(\boldsymbol{\lambda}))$ $\subset \mathcal{D}_i$, it follows that

$$\mathcal{D}_{i} = \left(\mathcal{N}\left(W_{i}(\boldsymbol{\lambda})\right) \oplus \left(\mathcal{N}\left(W_{i}(\boldsymbol{\lambda})\right) \right)^{\perp} \right) \cap \mathcal{D}_{i} = \mathcal{N}\left(W_{i}(\boldsymbol{\lambda})\right) + \left(\mathcal{N}\left(W_{i}(\boldsymbol{\lambda})\right)^{\perp} \cap \mathcal{D}_{i} \right).$$

Since $\mathcal{N}(W_i(\boldsymbol{\lambda}))$ is finite-dimensional and \mathcal{D}_i infinite-dimensional, \mathcal{D}'_i is also an infinite-dimensional vector space.

The following is a technical observation, which enables us to shorten proofs but has no significance otherwise. In actual calculations of basis vectors we can choose vectors $x_{ik}^g \in \mathcal{D}(A_{i0})$ (i.e. not necessarily $x_{ik}^g \in \mathcal{D}_i'$) that satisfy all other required

Lemma 3.3. Suppose that a vector $x_i \in (\mathcal{N}(W_i(\boldsymbol{\lambda})^*))^{\perp}$. Then there exists a vector $v_i \in \mathcal{D}'_i$ such that $x_i = W_i(\lambda) v_i$.

Proof. Suppose that $x_i \in (\mathcal{N}(W_i(\boldsymbol{\lambda})^*))^{\perp}$. Since $W_i(\boldsymbol{\lambda})$ is a Fredholm operator, its range is a closed subspace of H_i [25] and therefore $(\mathcal{N}(W_i(\boldsymbol{\lambda})^*))^{\perp} = \mathcal{R}(W_i(\boldsymbol{\lambda}))$. Thus there exists a vector $u_i \in \mathcal{D}(A_{i0})$ such that $x_i = W_i(\boldsymbol{\lambda})u_i$. Since $\mathcal{D}(W_i(\boldsymbol{\lambda})) = \mathcal{D}(A_{i0})$ we can find vectors $v_i \in \mathcal{D}'_i$ and $w_i \in \mathcal{N}(W_i(\boldsymbol{\lambda}))$ such that $v_i + w_i = u_i$. Then it follows that $x_i = W_i(\boldsymbol{\lambda})v_i$.

The following result is a special case of [5, Thm. 6.3]. Note that the additional condition that $x_{i2}^k \in \mathcal{D}'_i$ is a consequence of Lemma 3.3.

Theorem 3.4. Assume that $\lambda \in \sigma(\mathbf{W})$ is a derogatory eigenvalue, i.e., $\mathcal{S}_2 = 0$, and that $\mathbf{c}_{11}^{12} = \begin{bmatrix} 1 & 0 \end{bmatrix}^T$ and $\mathbf{c}_{12}^{12} = \begin{bmatrix} 0 & 1 \end{bmatrix}^T$. Then there exist vectors $x_{i2}^k \in \mathcal{D}_i'$ such that

(3.6)
$$U_i \left(\mathbf{c}_{1k}^{12} \right) x_{i1} + W_i \left(\lambda \right) x_{i2}^k = 0$$

for i, k = 1, 2. Furthermore, the vectors

$$(3.7) z_{2k} = x_{11} \otimes x_{22}^k + x_{12}^k \otimes x_{21}$$

are such that

$$(\Gamma_{i-1} - \lambda_{i-1}I) z_{2k} = \delta_{ik} z_1$$

for i, k = 1, 2, and $\{z_1, z_{21}, z_{22}\}$ is a basis for \mathcal{R}_2 . Here δ_{ik} is the Kronecker symbol.

Next we define the matrix S_3 by

$$\mathcal{S}_{3} = \left[\begin{array}{c|c} \left\langle A_{10}x_{12}^{1}, y_{11} \right\rangle & \left\langle A_{10}x_{12}^{2}, y_{11} \right\rangle + \left\langle A_{11}x_{12}^{1}, y_{11} \right\rangle & \left\langle A_{11}x_{12}^{2}, y_{11} \right\rangle \\ \left\langle A_{20}x_{22}^{2}, y_{21} \right\rangle + \left\langle A_{21}x_{22}^{2}, y_{21} \right\rangle & \left\langle A_{21}x_{22}^{2}, y_{21} \right\rangle \end{array} \right].$$

We identify the vector space Θ of all symmetric 2×2 matrices with the vector space \mathbb{C}^3 via the isomorphism $\psi : \Theta \longrightarrow \mathbb{C}^3$ defined by

(3.8)
$$\psi(C) = \begin{bmatrix} c_{11} & c_{12} & c_{22} \end{bmatrix}^T$$

for

$$C = \left[\begin{array}{cc} c_{11} & c_{12} \\ c_{12} & c_{22} \end{array} \right] \in \Theta.$$

The following theorem gives the structure of basis vectors of \mathcal{R}_3 .

Theorem 3.5. Suppose that $\{\mathbf{c}_l: 1 \leq l \leq d\}$ is a basis for $\mathcal{N}(\mathcal{S}_3)$ and that $C_l^3 = \left[c_{jkl}^{223}\right]_{j,k=1}^2 = \psi^{-1}(\mathbf{c}_l) \in \Theta$ for $1 \leq l \leq d$. Then we have $d = d_3$ and there exist vectors $x_{i3}^l \in \mathcal{D}_i'$, $i = 1, 2, l = 1, 2, \ldots, d$, such that

(3.9)
$$\sum_{k=1}^{2} U_i \left(\mathbf{c}_{kl}^{23} \right) x_{i2}^k + W_i \left(\mathbf{\lambda} \right) x_{i3}^l = 0,$$

where $\mathbf{c}_{kl}^{23} = \begin{bmatrix} c_{1kl}^{223} & c_{2kl}^{223} \end{bmatrix}^T$. The vector

(3.10)
$$z_{31} = x_{11} \otimes x_{23}^l + x_{13}^l \otimes x_{21} + \sum_{i,k=1}^2 c_{jkl}^{223} x_{12}^j \otimes x_{22}^k$$

is in $\mathcal{R}_3 \backslash \mathcal{R}_2$, and

(3.11)
$$(\Gamma_{i-1} - \lambda_{i-1}I) z_{3l} = \sum_{k=1}^{2} c_{ikl}^{223} z_{2k}$$

for i = 1, 2 and $1 \le l \le d$. Moreover, $\{z_1, z_{21}, z_{22}, z_{31}, z_{32}, \dots, z_{3d}\}$ is a basis for \mathcal{R}_3 .

Conversely, choose vectors $z_{31}, z_{32}, \ldots, z_{3d_2} \in \mathcal{R}_3 \backslash \mathcal{R}_2$ such that $\{z_1, z_{21}, z_{22}, z_{31}, z_{32}, \ldots, z_{3d_3}\}$ is a basis for \mathcal{R}_3 and such that (3.11) holds for both i and all l. Then

$$C_l^3 = \left[\begin{array}{cc} c_{11l}^{223} & c_{12l}^{223} \\ c_{21l}^{223} & c_{22l}^{223} \end{array} \right]$$

are symmetric, $\{\psi\left(C_l^3\right): 1 \leq l \leq d_3\}$ is a basis for $\mathcal{N}\left(\mathcal{S}_3\right)$, and there exist vectors $x_{i3}^l \in \mathcal{D}_i', i = 1, 2, \text{ such that } (3.9) \text{ and } (3.10) \text{ hold.}$

Proof. Because $\mathbf{c}_l \in \mathcal{N}\left(\mathcal{S}_3\right)$ and $C_l^3 = \psi^{-1}\left(\mathbf{c}_l\right)$ it follows that

$$\sum_{j=1}^{2} \sum_{k=1}^{2} c_{jkl}^{223} \left\langle A_{i,j-1} x_{i2}^{l}, y_{i1} \right\rangle = 0$$

for i = 1, 2. Thus $\sum_{k=1}^{2} U_i\left(\mathbf{c}_{kl}^{23}\right) x_{i2}^k \in \left(\mathcal{N}\left(W_i\left(\boldsymbol{\lambda}\right)^*\right)\right)^{\perp}$ and by Lemma 3.3 there exist vectors $x_{i3} \in \mathcal{D}_i'$ such that (3.9) hold for i = 1, 2. Next we form the vector z_{3l} as in (3.10). In the following calculation we assume that i = 0. For i = 1 the calculation is similar and we omit it. We have that

$$(\Delta_{0} - \lambda_{0}\Delta_{2}) z_{3l}$$

$$= (A_{11} \otimes W_{2}(\lambda) - W_{1}(\lambda) \otimes A_{21}) \left(x_{11} \otimes x_{23}^{l} + x_{13}^{l} \otimes x_{21} + \sum_{j,k=1}^{2} c_{jkl}^{223} x_{12}^{j} \otimes x_{22}^{k} \right)$$

$$= \sum_{j,k=1}^{2} c_{jkl}^{223} \left(A_{1,k-1} x_{12}^{j} \otimes A_{21} x_{21} - A_{11} x_{11} \otimes A_{2,j-1} x_{22}^{k} \right)$$

$$+ \sum_{j,k=1}^{2} c_{jkl}^{223} \left(-A_{11} x_{12}^{j} \otimes A_{2,k-1} x_{21} + A_{1,j-1} x_{11} \otimes A_{21} x_{22}^{k} \right)$$

$$= \sum_{j,k=1}^{2} c_{jkl}^{223} \left(-A_{11} \otimes A_{2,j-1} + A_{1,j-1} \otimes A_{21} \right) \left(x_{12}^{k} + x_{11} \otimes x_{22}^{k} \right)$$

$$= \Delta_{2} \left(\sum_{k=1}^{2} c_{1kl}^{223} z_{2k} \right).$$

Conversely, suppose that $z_3 \in \mathcal{R}_3 \setminus \mathcal{R}_2$. Then there exist a symmetric matrix

$$C = \left[\begin{array}{cc} c_{11} & c_{12} \\ c_{12} & c_{22} \end{array} \right] \in \Theta$$

and a vector $\mathbf{a} = \begin{bmatrix} a_0 & a_1 \end{bmatrix}^T \in \mathbb{C}^2$ such that for i = 0, 1, we have

$$(\Gamma_i - \lambda_i I) z_3 = \sum_{k=1}^2 c_{i+1,k} z_{2k} + a_i z_1.$$

We may and will assume that $\mathbf{a} = 0$. This is achieved by substituting $z_3 - a_0 z_{21} - a_1 z_{22}$ for z_3 . Thus we have

(3.12)
$$(\Gamma_i - \lambda_i I) z_3 = \sum_{k=1}^2 c_{i+1,k} z_{2k}.$$

Next it follows from (2.1) that

$$\sum_{i=0}^{1} A_{ij}^{\dagger} \left(\Gamma_{j} - \lambda_{j} I \right) z_{3} + W_{i} \left(\boldsymbol{\lambda} \right)^{\dagger} z_{3} = 0$$

and from (3.12) that

(3.13)
$$\sum_{j=0}^{1} A_{ij}^{\dagger} \left(\sum_{k=1}^{2} c_{j+1,k} z_{2k} \right) + W_{i} \left(\boldsymbol{\lambda} \right)^{\dagger} z_{3} = 0.$$

For i=1,2, we choose vectors $v_i \in H_i$ so that $\langle x_{i1}, v_i \rangle = 1$ and $\langle x_{i2}^k, v_i \rangle = 0$ for k=1,2. This is possible because Span $\{x_{i1}\} \cap \mathcal{D}'_i = \{0\}$. For i=1 and i=2, respectively, we form the scalar product of (3.13) by $y_{11} \otimes v_2$ and $v_1 \otimes y_{21}$, respectively, to get

(3.14)
$$\sum_{j=0}^{1} \left\langle A_{ij} \sum_{k=1}^{2} c_{j+1,k} x_{i2}^{k}, y_{i1} \right\rangle = 0.$$

Hence, by Lemma 3.3, there exist vectors $x_{i3} \in \mathcal{D}'_i$ such that

(3.15)
$$\sum_{k=1}^{2} U_{i}(\mathbf{c}_{k}) x_{i2}^{k} + W_{i}(\boldsymbol{\lambda}) x_{i3} = 0,$$

where $\mathbf{c}_k = \begin{bmatrix} c_{1k} & c_{2k} \end{bmatrix}^T$. Now we form the vector

$$z_{31} = x_{11} \otimes x_{23} + x_{13} \otimes x_{21} + \sum_{j,k=1}^{2} c_{jk} x_{12}^{j} \otimes x_{22}^{k}.$$

The same calculation as in the first part of the proof shows that

$$(\Gamma_i - \lambda_i I) z_{31} = \sum_{k=1}^{2} c_{i+1,k} z_{2k}$$

for i=0,1. It follows that $z_{31}-z_3 \in \mathcal{R}_1$ and so there exists a number $\delta \in \mathbb{C}$ such that $z_3=z_{31}+\delta z_1$. Without loss, we can use the vector $x_{13}+\delta x_{11}$ in place of x_{13} . Then it follows that

$$z_3 = x_{11} \otimes x_{23} + x_{13} \otimes x_{21} + \sum_{j,k=1}^{2} c_{jk} x_{12}^{j} \otimes x_{22}^{k}.$$

The final step is to verify that $\psi(T) \in \mathcal{N}(S_3)$, which follows immediately from the equalities (3.14).

Note that Theorem 3.1 implies the existence of vectors $z_{3l} \in \mathcal{R}_3 \backslash \mathcal{R}_2$ $(1 \le l \le d_3)$ such that the following hold:

- $\{z_1, z_{21}, z_{22}, z_{31}, z_{32}, \dots, z_{3d_3}\}$ is a basis for \mathcal{R}_3 ,
- equality (3.11) holds for both i and all l,
- $\{C_l^3: 1 \le l \le d_3\}$ are linearly independent.

Since Γ_0 and Γ_1 commute it follows that $c_{12l}^{223}=c_{21l}^{223}$ $(1 \leq l \leq d_3)$ and therefore the matrices C_l^3 are symmetric.

4. A basis for root subspace \mathcal{R}_m , $m \geq 3$, at a derogatory eigenvalue

Our goal in this section is to prove, by induction on m, the following theorem which is our main result.

Theorem 4.1. Suppose that $\left\{\mathbf{c}_{1f}^{m}, f=1,2,\ldots,d\right\}$ is a basis for the kernel of \mathcal{S}_m . Then $d=d_m$ and there exist vectors $x_{im}^f \in \mathcal{D}_i'$ and matrices C_f^m such that conditions (i) to (iii) below hold when m-1 is replaced by m. In particular, the union $\mathcal{B}_m = \mathcal{B}_{m-1} \cup \{z_{m1}, z_{m2}, \dots, z_{md}\}, \text{ where }$

$$(4.1) z_{mf} = x_{1m}^f \otimes x_{21} + x_{11} \otimes x_{2m}^f + \sum_{k=2}^{\left[\frac{m}{2}\right]} \sum_{l=k}^{m+1-k} \sum_{g=1}^{d_k} \sum_{h=1}^{d_l} c_{ghf}^{klm} \left(x_{1k}^g \otimes x_{2l}^h + x_{1l}^h \otimes x_{2k}^g \right)$$

for f = 1, 2, ..., d, is a basis for the root subspace \mathcal{R}_m .

Here $\left\lceil \frac{m}{2} \right\rceil$ denotes the integer part of $\frac{m}{2}$.

We will prove Theorem 4.1 at the end of the section. Let us first introduce the inductive assumptions. For $m = 1, 2, 3, \dots$ we suppose that we have vectors $x_{i1} \in$ $\mathcal{D}(A_{i0}), x_{il}^h \in \mathcal{D}'_i, i = 1, 2, l = 2, 3, \dots, m-1, h = 1, 2, \dots, d_l, \text{ and symmetric}$ (structure) matrices (see (3.4))

$$(4.2) C_h^l = \begin{bmatrix} C_h^{22l} & C_h^{23l} & \cdots & C_h^{2,l-2,l-1} & C_h^{2,l-1,l} \\ \left(C_h^{23l}\right)^T & C_h^{33l} & \cdots & C_h^{3,l-2,l} & 0 \\ \vdots & \vdots & & \vdots & \vdots \\ \left(C_h^{2,l-2,l}\right)^T & \left(C_h^{3,l-2,l}\right)^T & \cdots & 0 & 0 \\ \left(C_h^{2,l-1,l}\right)^T & 0 & \cdots & 0 & 0 \end{bmatrix}$$

such that

(i) the matrices $C_h^{2,l-1,l}$, $h=1,2,\ldots,d_l$, are linearly independent for $l=1,2,\ldots,d_l$ $3, 4, \ldots, m-1$, the column cross-sections (see (3.3)) \mathcal{C}_h^{jl} are equal to \mathcal{C}_h^{2jl} , and the entries of matrices C_h^l satisfy the relations

$$(4.3) \qquad \sum_{k=l_1+l_2-1}^{l-l_3+1} \sum_{g=1}^{d_k} c_{h_1 h_2 g}^{l_1 l_2 k} c_{h_3 g h}^{l_3 k l} = \sum_{k=l_2+l_3-1}^{l-l_1+1} \sum_{g=1}^{d_k} c_{h_1 h_3 g}^{l_1 l_3 k} c_{h_2 g h}^{l_2 k l}$$

for $l = 4, 5, \ldots, m - 1, \mathbf{l} \in \Phi_l$ and $\mathbf{h} \in \chi_{\mathbf{l}}$,

(ii) the relations (2.2), (3.6) and

(4.4)
$$\sum_{k=2}^{l-1} \sum_{g=1}^{d_k} U_i \left(\mathbf{c}_{gh}^{kl} \right) x_{ik}^g + W_i \left(\mathbf{\lambda} \right) x_{il}^h = 0$$

for $i=1,2,\ l=3,4,\ldots,m-1$ and $h=1,2,\ldots,d_l$ hold (here $\mathbf{c}_{gh}^{kl}=1,\ldots,d_l$ $\begin{bmatrix} c_{1gh}^{2kl} & c_{2gh}^{2kl} \end{bmatrix}^T$), (iii) the vectors z_1 , z_{21} , z_{22} , of (2.3) and (3.7) together with vectors z_{31} , z_{32} ,

 $\ldots, z_{3d_3}, \ldots, z_{m-1,1}, z_{m-1,2}, \ldots, z_{m-1,d_{m-1}},$ that are defined by

$$z_{lh} = x_{1l}^h \otimes x_{21} + x_{11} \otimes x_{2l}^h + \sum_{k=2}^{\left[\frac{l}{2}\right]} \sum_{r=k}^{l+1-k} \sum_{g=1}^{d_k} \sum_{e=1}^{d_r} c_{geh}^{krl} \left(x_{1k}^g \otimes x_{2r}^e + x_{1r}^e \otimes x_{2k}^g \right)$$

for l = 3, 4, ..., m-1 and $h = 1, 2, ..., d_l$, form a filtered basis for the root subspace \mathcal{R}_{m-1} . We also write

$$\mathcal{B}_{m-1} = \left\{ z_1, \ z_{21}, z_{22}, \dots, z_{m-1,1}, z_{m-1,2}, \dots, z_{m-1,d_{m-1}} \right\}.$$

By Theorem 3.5 it follows that there exist vectors $x_{i1}, x_{i2}^h, x_{i3}^h$ and matrices C_2^h such that conditions (i) to (iii) are satisfied for m-1=3. Now we assume that the above conditions are satisfied for m-1 (with $3 \le m-1 \le M-1$) and we prove that we can find vectors $x_{im}^f \in \mathcal{D}'_i$, $i = 1, 2, f = 1, 2, \ldots, d_m$, and matrices C_f^m , $f = 1, 2, ..., d_m$, such that (i) to (iii) hold with m - 1 replaced by m.

Our next step is to introduce matrices \mathcal{S}_m . We do this in two stages. First we

use matrices $C_q^{\bar{k}}$, k < m, and second, certain orthogonality relations.

For the purpose of calculation we write $T^m = \begin{bmatrix} T^{klm} \end{bmatrix}_{k,l=2}^{m-1}$ for an unknown symmetric matrix in the form (4.2). The entries of the $d_k \times d_l$ matrix T^{klm} are denoted by t_{ah}^{klm} . They must satisfy the conditions

$$(4.5) \qquad \sum_{k=l_1+l_2-1}^{m-l_3+1} \sum_{g=1}^{d_k} c_{h_1 h_2 g}^{l_1 l_2 k} t_{h_3 g}^{l_3 k m} - \sum_{k=l_1+l_3-1}^{m-l_2+1} \sum_{g=1}^{d_k} c_{h_1 h_3 g}^{l_1 l_3 k} t_{h_2 g}^{l_2 k m} = 0$$

for $\mathbf{l} \in \Phi_m$ and $\mathbf{h} \in \chi_1$. We write the $2 \times d_{m-1}$ matrix $T^{1,m-1,m}$ also as a column

$$\mathbf{t}_{1}^{m} = \left[\begin{array}{ccccc} t_{11}^{2,m-1,m} & t_{21}^{2,m-1,m} & t_{12}^{2,m-1,m} & t_{22}^{2,m-1,m} & \cdots & t_{1d_{m-1}}^{2,m-1,m} & t_{2d_{m-1}}^{2,m-1,m} \end{array} \right]^{T}.$$

We denote the sum $\sum_{i=1}^{l} d_i$ by \widehat{d}_l . For every column of T^m we define a column vector \mathbf{t}_{kq}^m , k = 2, 3, ..., m - 1, $g = 1, 2, ..., d_k$, of size $\nu = \min \left\{ \widehat{d}_{k-1} + g, \ \widehat{d}_{m-k} \right\}$

by taking the first ν entries in the $(\widehat{d}_{k-1}+g)$ -th column of T^m . Note that \mathbf{t}_{kq}^m are defined so that they consist of all the entries above and including the main diagonal of the matrix T^m omitting those labelled 0 in (4.2). We define a column vector \mathbf{t}_2^m

(4.7)
$$\mathbf{t}_{2}^{m} = \begin{bmatrix} (\mathbf{t}_{11}^{m})^{T} & (\mathbf{t}_{12}^{m})^{T} & (\mathbf{t}_{21}^{m})^{T} & \cdots & (\mathbf{t}_{2d_{2}}^{m})^{T} & \cdots \\ (\mathbf{t}_{m-2,1}^{m})^{T} & \cdots & (\mathbf{t}_{m-2,d_{m-2}}^{m})^{T} \end{bmatrix}^{T}.$$

We split the entries of the matrix T^m into two column vectors \mathbf{t}_1^m and \mathbf{t}_2^m . The mapping $\psi_m: T^m \mapsto (\mathbf{t}_1^m, \mathbf{t}_2^m)$, defined on all symmetric matrices of the form (4.2), is a generalization of the transformation ψ defined by (3.8). It is bijective and therefore it has an inverse, which maps two vectors \mathbf{t}_1^m and \mathbf{t}_2^m into a matrix T^m of the form (4.2). We use this inverse mapping in Lemma 4.3. We write the system of equations (4.5) in matrix form as

$$(4.8) S_m^{21} \mathbf{t}_1^m + S_m^{22} \mathbf{t}_2^m = 0$$

where the entries of the matrices S_m^{21} and S_m^{22} are determined by the system (4.5). These entries are given because we assumed that the matrices C_g^k were known for

We also want the entries of the matrices T^{2km} to satisfy the scalar relations

(4.9)
$$\sum_{k=2}^{m-1} \sum_{g=1}^{d_k} \left\langle \left(t_{1g}^{2km} A_{i0} + t_{2g}^{2km} A_{i1} \right) x_{ik}^g, y_{i1} \right\rangle = 0$$

for both i. These can be written equivalently in matrix form

$$(4.10) S_m^{11} \mathbf{t}_1^m + S_m^{12} \mathbf{t}_2^m = 0.$$

Again the entries of the matrices S_m^{11} and S_m^{12} are determined by the equations (4.9).

Now we are prepared to define the matrices S_m . First we choose a matrix B_m so that its columns form a basis for the kernel of

$$[(S_m^{12})^* (S_m^{22})^*].$$

Then we define a matrix

$$\mathcal{S}_m = B_m^* \left[\begin{array}{c} S_m^{11} \\ S_m^{21} \end{array} \right],$$

so equations (4.8) and (4.10) yield $S_m \mathbf{t}_1^m = 0$.

We now choose vectors $z_{mf} \in \mathcal{D}$, $f = 1, 2, ..., d_m$, so that $\mathcal{B}_{m-1} \cup \{z_{m1}, z_{m2}, ..., z_{md_m}\}$, is a basis for the space \mathcal{R}_m . By Theorem 3.1 there exist matrices C_f^m for $f = 1, 2, ..., d_m$ such that (i) holds and that we have

(4.11)
$$(\Gamma_{i-1} - \lambda_{i-1}I) z_{mf} = \sum_{k=2}^{m-1} \sum_{q=1}^{d_k} c_{igf}^{2km} z_{kg}$$

for i = 1, 2 and $f = 1, 2, ..., d_m$. Next we prove three auxiliary results.

Lemma 4.2. In the above setting it follows that $\dim \mathcal{N}(\mathcal{S}_m) \geq d_m$.

Proof. Let $\psi_m\left(C_f^m\right) = \left(\mathbf{c}_{1f}^m, \mathbf{c}_{2f}^m\right)$. By Theorem 3.1(e) and the definition of ψ_m it follows that $S_m^{21}\mathbf{c}_{1f}^m + S_m^{22}\mathbf{c}_{2f}^m = 0$. Relation (2.1) implies

$$\sum_{j=0}^{1} A_{ij}^{\dagger} \left(\Gamma_{j} - \lambda_{j} I \right) z_{mf} + W_{i} \left(\boldsymbol{\lambda} \right)^{\dagger} z_{mf} = 0$$

for both i. By relations (4.11) we get that

(4.12)
$$\sum_{k=2}^{m-1} \sum_{g=1}^{d_k} U_i \left(\mathbf{c}_{gf}^{km} \right)^{\dagger} z_{kg} + W_i \left(\boldsymbol{\lambda} \right)^{\dagger} z_{mf} = 0.$$

Here $\mathbf{c}_{gf}^{km} = \begin{bmatrix} c_{1gf}^{2km} & c_{2gf}^{2km} \end{bmatrix}^T$. Because we assume that $x_{ik}^g \in \mathcal{D}_i'$ and $\mathcal{D}_i' \cap \operatorname{Span}\{x_{i1}\} = \{0\}$ it follows that for each i there exists a vector $v_i \in H_i$ such that $\langle x_{i1}, v_i \rangle = 1$ and $\langle x_{ik}^g, v_i \rangle = 0$ for $k = 2, 3, \ldots, m-1$ and $g = 1, 2, \ldots, d_k$. Then the scalar products of (4.12) by $y_{11} \otimes v_2$ and $v_1 \otimes y_{21}$, for i = 1 and i = 2 respectively, are equal to 0. It follows, using the structure of the vectors z_{kg} in condition (iii), that

(4.13)
$$\sum_{k=2}^{m-1} \sum_{g=1}^{d_k} \langle U_i \left(\mathbf{c}_{gf}^{km} \right) x_{ik}^g, y_{i1} \rangle = 0$$

for all f, and so

$$\sum_{i=1}^{2} \sum_{k=2}^{m-1} \sum_{a=1}^{d_k} c_{jgf}^{2km} \langle A_{i,j-1} x_{ik}^g, y_{i1} \rangle = 0.$$

Therefore \mathbf{c}_{1f}^m and \mathbf{c}_{2f}^m solve equation (4.10) and \mathbf{c}_{1f}^m , $f = 1, 2, ..., d_m$, are elements of the kernel of the matrix \mathcal{S}_m . Since these vectors are linearly independent, we have $d_m \leq \dim \mathcal{N}(\mathcal{S}_m)$.

Lemma 4.3. Suppose that \mathbf{c}_{11}^m is an element of the kernel $\mathcal{N}(\mathcal{S}_m)$. Then there exists a vector \mathbf{c}_{21}^m such that (4.8) and (4.10) hold. Furthermore, there exist vectors $x_{im}^1 \in \mathcal{D}'_i$, i = 1, 2, such that

(4.14)
$$\sum_{k=2}^{m-1} \sum_{g=1}^{d_k} U_i \left(\mathbf{c}_{g1}^{km} \right) x_{ik}^g + W_i \left(\mathbf{\lambda} \right) x_{im}^1 = 0.$$

Here \mathbf{c}_{q1}^{km} are determined by the first block column of $C_1^m = \psi_m^{-1}(\mathbf{c}_1^m, \mathbf{c}_2^m)$.

Proof. From the structure of the matrix S_m it follows that for an element $\mathbf{c}_{11}^m \in \mathcal{N}(S_m)$ there exists a vector \mathbf{c}_{21}^m such that relations (4.8) and (4.10) hold. Using the inverse of the isomorphism ψ_m^{-1} we associate with the pair of vectors \mathbf{c}_{11}^m and \mathbf{c}_{21}^m a matrix C_1^m . The relations (4.10) can be written equivalently in the form (4.13). Then it follows for both i that

$$\sum_{k=2}^{m-1} \sum_{g=1}^{d_k} U_i\left(\mathbf{c}_{g1}^{km}\right) x_{ik}^g \in \left(\mathcal{N}\left(W_i\left(\boldsymbol{\lambda}\right)^*\right)\right)^{\perp} = \mathcal{R}\left(W_i(\boldsymbol{\lambda})\right).$$

The latter equality holds since $W_i(\lambda)$ is a Fredholm operator. By Lemma 3.3 there exists a vector $x_{im}^1 \in \mathcal{D}'_i$ such that (4.14) holds.

Lemma 4.4. Suppose that we have the same setting as in Lemma 4.3. We construct a vector

$$z_{m1} = x_{1m}^1 \otimes x_{21} + x_{11} \otimes x_{2m}^1 + \sum_{k=1}^{\left[\frac{m}{2}\right]} \sum_{l=k}^{m+1-k} \sum_{g=1}^{d_k} \sum_{h=1}^{d_l} c_{gh1}^{klm} \left(x_{1k}^g \otimes x_{2l}^h + x_{1l}^h \otimes x_{2k}^g \right).$$

Then it follows that

(4.15)
$$(\Gamma_{j-1} - \lambda_{j-1}I) z_{m1} = \sum_{k=2}^{m-1} \sum_{q=1}^{d_k} c_{jg1}^{2km} z_{kq}$$

for j = 1, 2.

Proof. We use a direct calculation to show (4.15). We consider only the case j=1. The case j=2 is proved in the same way. By the basic properties of operator determinants we have

$$(4.16) \qquad (\Delta_0 - \lambda_0 \Delta_2) z_{m1} - \sum_{k=2}^{m-1} \sum_{g=1}^{d_k} c_{1g1}^{2km} \Delta_2 z_{kg}$$

$$(4.17) \qquad = \begin{vmatrix} -W_1(\boldsymbol{\lambda})^{\dagger} & A_{11}^{\dagger} \\ -W_2(\boldsymbol{\lambda})^{\dagger} & A_{21}^{\dagger} \end{vmatrix} \left(x_{1m}^{1} \otimes x_{21} + x_{11} \otimes x_{2m}^{1} \right)$$

$$(4.18) -\sum_{k=2}^{m-1} \sum_{g=1}^{d_k} \left| \begin{array}{c} U_1 \left(\mathbf{c}_{g1}^{km} \right)^{\dagger} & A_{11}^{\dagger} \\ U_2 \left(\mathbf{c}_{g1}^{km} \right)^{\dagger} & A_{21}^{\dagger} \end{array} \right| (x_{1k}^g \otimes x_{21} + x_{11} \otimes x_{2k}^g)$$

$$(4.19) \qquad + \left| \begin{array}{cc} -W_{1}\left(\boldsymbol{\lambda}\right)^{\dagger} & A_{11}^{\dagger} \\ -W_{2}\left(\boldsymbol{\lambda}\right)^{\dagger} & A_{21}^{\dagger} \end{array} \right| \sum_{k=2}^{\left[\frac{m}{2}\right]} \sum_{l=k}^{m+1-k} \sum_{g=1}^{d_{k}} \sum_{h=1}^{d_{l}} c_{gh1}^{klm} \left(x_{1k}^{g} \otimes x_{2l}^{h} + x_{1l}^{h} \otimes x_{2k}^{g} \right)$$

$$(4.20) -\sum_{k=3}^{m-1} \sum_{g=1}^{d_k} c_{1g1}^{2km} \sum_{l=2}^{\left[\frac{k}{2}\right]} \sum_{r=l}^{k+1-l} \sum_{h=1}^{d_l} \sum_{e=1}^{d_r} c_{ehg}^{rlk} \Delta_2 \left(x_{1r}^e \otimes x_{2l}^h + x_{1l}^h \otimes x_{2r}^e \right).$$

Relation (4.14) implies that the sum of (4.17) and (4.18) is equal to 0. Since $c_{ehg}^{rlk}=c_{heg}^{lrk}$ it follows that (4.20) is equal to

$$-\sum_{k=3}^{m-1}\sum_{l=1}^{k-1}\sum_{r=l}^{k+1-l}\sum_{g=1}^{d_k}\sum_{h=1}^{d_l}\sum_{e=1}^{d_r}c_{1g1}^{2km}c_{ehg}^{rlk}\Delta_2x_{1r}^e\otimes x_{2l}^h$$

$$\xrightarrow{m-2}\sum_{r=1}^{r-1}\sum_{d_l}\sum_{d_r}\sum_{h=1}^{m-1}\sum_{e=1}^{d_k}\sum_{ehg}\sum_{h=1}^{d_k}\sum_{ehg}\sum_{h=1}^{d_k}\sum_{h=1}$$

$$(4.21) = -\sum_{r=2}^{m-2} \sum_{l=2}^{r-1} \sum_{h=1}^{d_l} \sum_{e=1}^{d_r} \left[\sum_{k=l+r-1}^{m-1} \sum_{g=1}^{d_k} c_{ehg}^{2km} c_{ehg}^{rlk} \right] \Delta_2 x_{1r}^e \otimes x_{2l}^h.$$

We use relations (3.5) and the definition of Δ_2 , and we rearrange the order of summation to show that (4.21) is equal to (4.22)

$$-\sum_{k=2}^{\left[\frac{m}{2}\right]}\sum_{l=k}^{m+1-k}\sum_{g=1}^{d_{k}}\sum_{h=1}^{d_{l}}\sum_{r=k}^{l-1}\sum_{e=1}^{d_{r}}c_{gh1}^{klm}\left|\begin{array}{c}U_{1}\left(\mathbf{c}_{eg}^{rk}\right)^{\dagger}&A_{11}^{\dagger}\\U_{2}\left(\mathbf{c}_{eg}^{rk}\right)^{\dagger}&A_{21}^{\dagger}\end{array}\right|\left(x_{1r}^{e}\otimes x_{2l}^{h}+x_{1l}^{h}\otimes x_{2r}^{e}\right).$$

If we use (4.4) and (4.14) we see that the sum of (4.19) and (4.22) is 0. Therefore (4.16) is equal to 0, and (4.15) follows.

Now we are ready to prove our main result.

Proof of Theorem 4.1. Suppose that we are given a basis $\left\{\mathbf{c}_{1f}^m, f=1,2,\ldots,d\right\}$ for $\mathcal{N}\left(\mathcal{S}_m\right)$. Then by Lemma 4.3 it follows that we can find vectors $x_{im}^f \in \mathcal{D}_i'$ and matrices C_f^m of the form (4.2), where $C_f^m = \psi_m^{-1}\left(\mathbf{c}_{1f}^m, \mathbf{c}_{2f}^m\right)$ for some \mathbf{c}_{2f}^m , such that (i) and (ii) hold when m-1 is replaced by m. We apply Lemma 4.4 to show that for $f=1,2,\ldots,d$ the vectors z_{mf} are in \mathcal{R}_m but not in \mathcal{R}_{m-1} . They are linearly independent because \mathbf{c}_{1f}^m , $f=1,2,\ldots,d$, are linearly independent. It follows that $d\leq d_m$ and, because $d\geq d_m$ by Lemma 4.2, we have $d=d_m$. Thus (iii) holds also.

5. The Sturm-Liouville case

Faierman [8] considers a two-parameter eigenvalue problem involving a class of coupled Sturm-Liouville boundary value problems

(5.1)

$$\frac{d}{dt_{i}}\left(p_{i}\left(t_{i}\right)\frac{dy_{i}}{dt_{i}}\right) + \left((-1)^{i+1}\mu_{1}a_{i1}\left(t_{i}\right) + (-1)^{i}\mu_{2}a_{i2}\left(t_{i}\right) - q_{i}\left(t_{i}\right)\right)y_{i} = 0, \ i = 1, 2,$$

(5.2)
$$y_i(0)\cos\alpha_i - p_i(0)\frac{dy_i}{dt_i}(0)\sin\alpha_i = 0, \ 0 \le \alpha_i < \pi, \ i = 1, 2,$$

and

(5.3)
$$y_i(1)\cos\beta_i - p_i(1)\frac{dy_i}{dt_i}(1)\sin\beta_i = 0, \ 0 < \beta_i \le \pi, \ i = 1, 2,$$

where $\boldsymbol{\mu} = (\mu_1, \mu_2)$ are parameters, $\mathcal{I} = [0, 1]$ and $t_i \in \mathcal{I}$. As in [8, pp. 2 and 10] we assume that

- (i) for i = 1, 2 the functions $p_i, q_i, a_{ij}, j = 1, 2$, are real-valued, p_i, a_{ij} are Lipschitz continuous, p_i is positive and q_i is essentially bounded,
- (ii) the function $\omega(t_1, t_2) = a_{11}(t_1)a_{22}(t_2) a_{12}(t_1)a_{21}(t_2)$ on \mathcal{I}^2 is not identically 0.
- (iii) $a_{i1}(t_i) > 0$ for $t_i \in \mathcal{I}$ and i = 1, 2.

Weaker assumptions would be possible (cf. [4, 5, 26]) but we adhere to Faierman's assumptions for the purpose of comparison.

Two-parameter problems which satisfy condition (iii) are called *elliptic*. In the literature also a formally stronger condition

(iii') each $a_{ij}(t_i)$ is positively bounded below on \mathcal{I}

is considered (e.g. [3]). Two-parameter problems which satisfy condition (iii') are called *uniformly elliptic*. We show in Lemma 5.1 that we can transform an elliptic problem to a uniformly elliptic one by a linear substitution in parameters.

To introduce the operators A_{ij} we take $H_i = L_2(\mathcal{I})$. We have

(5.4)
$$A_{ij}y_i(t_i) = a_{ij}(t_i)y_i(t_i), \quad i, j = 1, 2,$$

for $y_i \in H_i$ and

$$(5.5) A_{i0}y_i(t_i) = \frac{d}{dt_i} \left(p_i(t_i) \frac{dy_i}{dt_i} \right) - q_i(t_i) y_i, \quad i = 1, 2,$$

where the domain $\mathcal{D}(A_{i0})$ consists of all functions $y_i \in H_i$ such that y_i and $\frac{dy_i}{dt_i}$ are absolutely continuous, $\frac{d^2y_i}{dt_i^2} \in H_i$ and the boundary conditions (5.2) and (5.3) hold.

Lemma 5.1. Assume that the problem (5.1)–(5.3) is such that (i)–(iii) hold. Then there exists an invertible linear substitution of parameters μ such that the problem is uniformly elliptic after the substitution.

Proof. The numerical range

$$\kappa_i = \{(c_{i1}, c_{i2}) \in \mathbb{R}^2 : \text{ for } j = 1, 2 \ c_{ij} = (A_{ij}x_i, x_i) \text{ for some } x_i \in H_i, ||x_i|| = 1\}$$

is a bounded convex set for i=1,2. By conditions (i) and (iii) it follows that for some $\epsilon > 0$ we have $c_{11} > \epsilon$ for all elements of κ_1 and $c_{21} < -\epsilon$ for all elements of κ_2 . Then there exist two lines through the origin in \mathbb{R}^2 that separate κ_1 and κ_2 and we can find an invertible linear map τ on \mathbb{R}^2 that maps κ_1 and κ_2 into the first and the third quadrant, respectively. If we apply τ to the parameters μ it follows that the transformed problem is uniformly elliptic.

We hereafter assume, without loss of generality, that the parameters μ are chosen so that condition (iii') holds, i.e., that the two-parameter problem is uniformly elliptic. In order to apply the results of the previous sections we need to verify Assumptions I–III. By [8, Thm. 2.4] it follows that Δ_2 is a selfadjoint operator with dense domain \mathcal{D} of those functions in the Sobolev space $W_2^2\left(\mathcal{I}^2\right)$ that satisfy the boundary condition on \mathcal{I}^2 associated with (5.2) and (5.3). For details see [8, pp. 10–30]. Next we have:

Lemma 5.2.
$$\mathcal{N}(\Delta_0) \cap \mathcal{N}(\Delta_2) = \{0\}.$$

Proof. Assume that $\Delta_0 u = \Delta_2 u = 0$ for some $u \in \mathcal{D}$. By (ii) it follows that ω is nonzero on some open set $\Omega \subset \mathcal{I}^2$. Hence $u|_{\Omega} \equiv 0$. By the unique continuation property for Δ_2 [8, Prop. 3.1] it follows that $u \equiv 0$ on \mathcal{I}^2 .

Now it follows from [5, Prop. 2.1] that Assumption I holds for an $\alpha \in \mathbb{R}$. As in §2 we shall hereafter assume, without loss of generality, that $\alpha = 0$, i.e., that Δ_2 has a bounded inverse. This is achieved by a shift in parameters. Assumption II is a consequence of the following lemma. This is given in a special case in [5, §9.3]. Since the argument there appears to be incomplete we provide a more detailed proof here.

Lemma 5.3. The system

$$\begin{bmatrix} A_{10}^{\dagger} & A_{11}^{\dagger} \\ A_{20}^{\dagger} & A_{21}^{\dagger} \end{bmatrix} \begin{bmatrix} u_1 \\ u_2 \end{bmatrix} = \begin{bmatrix} A_{12}^{\dagger} z \\ A_{22}^{\dagger} z \end{bmatrix}$$

has a solution $u = (u_1, u_2) \in \mathcal{D}^2$ for all $z \in \mathcal{D}$.

Proof. By eliminating u_2 in (5.6) we have

$$(5.7) -\Delta_2 u_1 = \Delta_0 z.$$

By [8, Thm. 2.2] there exists $u_1 \in \mathcal{D}$ which solves (5.7). Similarly, the elimination of u_2 leads to $-\Delta_2 u_2 = \Delta_1 z$, which again is soluble for $u_2 \in \mathcal{D}$ by [8, Thm. 2.2]. Thus it remains to show that u_1 and u_2 satisfy (5.6).

We shall do this first under the assumption that each A_{i0} is positive definite. For $N \in \mathbb{N}$ we define bounded selfadjoint operators $A_{i0}^{(N)}$ sharing the same eigenvectors and the lowest N eigenvalues with A_{i0} . The remaining eigenvalues of $A_{i0}^{(N)}$ are defined as 1. The definition of $A_{i0}^{(N)}$ is completed by linearity and continuity to H_i . We define $\Delta_j^{(N)}$, j=1,2, analogously to Δ_j and we note that $A_{i0}^{(N)}$ and $\Delta_j^{(N)}$ are uniformly bounded with inverses bounded uniformly in N for i,j=1,2.

Now consider the system of equations

(5.8)
$$A_{i0}^{(N)\dagger} z_0 = A_{i1}^{\dagger} z_1 + A_{i2}^{\dagger} z_2, \quad i = 1, 2,$$

where $z_2 = \left(\Delta_1^{(N)}\right)^{-1} w$ for a fixed $w \in H$. Proceeding as in the proof of [2, Thms. 4.2 and 6.1] we obtain

$$z_0 = -\left(\Delta_2^{(N)}\right)^{-1} \Delta_0 \left(\Delta_1^{(N)}\right)^{-1} w \text{ and } z_1 = -\left(\Delta_2^{(N)}\right)^{-1} w,$$

so (5.8) yields

(5.9)
$$A_{i0}^{(N)\dagger} x_N = -A_{i1}^{\dagger} \left(\Delta_2^{(N)} \right)^{-1} w - A_{i2}^{\dagger} \left(\Delta_1^{(N)} \right)^{-1} w, \quad i = 1, 2,$$

where $x_N = \left(\Delta_2^{(N)}\right)^{-1} \Delta_0 \left(\Delta_1^{(N)}\right)^{-1} w$. Now [2, Lemma 4.1] shows that the right-hand side in (5.9) has limit $-A_{i1}^{\dagger} \Delta_2^{-1} w - A_{i2}^{\dagger} \Delta_1^{-1} w = y_i$, say, as $N \to \infty$. Thus $A_{i0}^{(N)\dagger} x_N$ converges to y_i . Since the inverses of $A_{i0}^{(N)}$ are bounded uniformly in N it follows that $x_N - \left(\left(A_{i0}^{(N)}\right)^{-1}\right)^{\dagger} y_i \to 0$ as $N \to \infty$. Similarly, since the inverses of $\left(\Delta_2^{(N)}\right)^{-1}$ are bounded uniformly in N it follows that $x_N \to \Delta_2^{-1} \Delta_0 \Delta_1^{-1} w$ as $N \to \infty$ (cf. [2, Lemma 4.1]). Thus

$$\Delta_2^{-1} \Delta_0 \Delta_1^{-1} w = \left(A_{i0}^{-1} \right)^{\dagger} y_i,$$

whence

$$A_{i0}^{\dagger} \Delta_2^{-1} \Delta_0 \Delta_1^{-1} w = -A_{i1}^{\dagger} \Delta_2^{-1} w - A_{i2}^{\dagger} \Delta_1^{-1} w.$$

Now put $w = \Delta_1 z$ to obtain (5.6) for $u_1 = -\Delta_2^{-1} \Delta_0 z$ and $u_2 = -\Delta_2^{-1} \Delta_1 z$. To remove the restriction that A_{i0} are uniformly positive we use the technique of [6, Thms. 3.1 and 3.2].

In the rest of this section we discuss the eigenvalues and associated root subspaces of (5.1)–(5.3). If there exist nonzero functions $y_i \in \mathcal{D}(A_{i0})$, i=1,2, which solve (5.1)–(5.3) for a pair $\boldsymbol{\mu}=(\mu_1,\mu_2)\in\mathbb{C}^2$, then $\boldsymbol{\mu}$ is called an eigenvalue for the problem (5.1)–(5.3). The function $w(t_1,t_2)=y_1(t_1)y_2(t_2)$ is such that $(\Delta_i-\mu_i\Delta_0)\,w=0$ for i=1,2. Because we assume that Δ_2 has a bounded inverse, and in particular, is one-to-one, we have that $\mu_2\neq 0$. Then it follows for each eigenvalue $\boldsymbol{\mu}$ of (5.1)–(5.3) that $\boldsymbol{\lambda}=(\lambda_0,\lambda_1)=\left(\frac{1}{\mu_2},\frac{\mu_1}{\mu_2}\right)$ is an eigenvalue of the associated two-parameter system

$$(5.10) W_i(\lambda) = A_{i0}\lambda_0 + A_{i1}\lambda_1 + A_{i2}, \quad i = 1, 2.$$

Here A_{ij} are given by (5.4) and (5.5). We shall hereafter keep the notation $\boldsymbol{\mu}$ and $\boldsymbol{\lambda}$ as above, so if $\boldsymbol{\lambda} = (\lambda_0, \lambda_1)$ is an eigenvalue of (5.10), then $\boldsymbol{\mu} = (\mu_1, \mu_2) = \left(\frac{\lambda_1}{\lambda_0}, \frac{1}{\lambda_0}\right)$ is the corresponding eigenvalue of (5.1)–(5.3). Note that μ_2 and λ_0 are always nonzero.

The eigenvalues of (5.1)–(5.3) form a countably infinite set of points without accumulation. There are only finitely many nonsemisimple eigenvalues and finitely many nonreal eigenvalues. Nonreal eigenvalues appear in conjugate pairs. See [8, Thms. 6.1 and 6.2] for these results.

The problem (5.1)–(5.3) for i=1 or i=2 has for a fixed value of $\boldsymbol{\mu}$ at most a one-dimensional space of solutions. If $\boldsymbol{\mu} \in \mathbb{R}^2$, then $W_i(\boldsymbol{\mu})$ is a selfadjoint operator. If $\boldsymbol{\mu} \in \mathbb{C}^2 \backslash \mathbb{R}^2$, then $W_i(\boldsymbol{\mu})^* = W_i(\overline{\boldsymbol{\mu}})$, where $\bar{}$ denotes complex conjugation. Since $W_i(\boldsymbol{\mu}) y_i = 0$ holds if and only if $W_i(\overline{\boldsymbol{\mu}}) \overline{y}_i = 0$, it follows that $\dim \mathcal{N}(W_i(\boldsymbol{\mu})) = \dim \mathcal{N}(W_i(\boldsymbol{\mu})^*) = 1$. We are now ready to establish Assumption III.

Lemma 5.4. Each eigenvalue μ of the problem (5.1)–(5.3) is geometrically simple of finite ascent, and both $W_i(\lambda)$, i = 1, 2, are Fredholm operators of index 0.

Proof. The remarks preceding this lemma together with [16, Example III.5.14] establish everything except finite ascent, and the latter follows from [8, Thm. 6.6].

There are various differences between the approach of [8] and the one here. In [8] the eigenvalues $\boldsymbol{\mu}=(\mu_1,\mu_2)$ are considered. We use instead the eigenvalues $\boldsymbol{\lambda}=(\lambda_0,\lambda_1)$. In [8] root subspaces are considered only with respect to the pencil $\Delta_2-\mu_2\Delta_0$, while we consider joint root subspaces for pencils $\Delta_i-\lambda_i\Delta_2$, i=0,1. When bases for root subspaces are constructed our approach significantly reduces technical difficulties since all the geometric eigenspaces are now one-dimensional, i.e., all the joint eigenvalues are geometrically simple.

In order to be able to compare our results with those of [8] we need the following auxiliary results:

Lemma 5.5. If $\lambda_0 \neq 0$ and $z_1, z_2, \dots, z_m \in \mathcal{D}$ are such that

$$(5.11) \qquad (\Delta_0 - \lambda_0 \Delta_2) z_j = \Delta_2 z_{j-1}, \ j = 1, 2, \dots, m,$$

where $z_0 = 0$, then

$$u_j = \sum_{k=0}^{j-1} (-1)^j \binom{j}{k} \lambda_0^{j+k} z_{k+1}$$

are such that

(5.12)
$$(\Delta_2 - \mu_2 \Delta_0) u_j = \Delta_0 u_{j-1}, \ j = 1, 2, \dots, m,$$
 where $u_0 = 0$ and $\mu_2 = \lambda_0^{-1}$.

Proof. By induction on j we prove that (5.11) implies

$$(\Delta_2 - \mu_2 \Delta_0) z_j = \sum_{k=1}^{j-1} (-1)^k \mu_2^{k+1} \Delta_0 z_{j-k}, \quad j = 1, 2, \dots, m.$$

Then (5.12) follows by a straightforward calculation which we omit.

We call a chain of vectors u_1, u_2, \ldots, u_m such that (5.12) holds a *Jordan chain* for the pencil $\Delta_2 - \mu_2 \Delta_0$.

Corollary 5.6. Let $\mu_2 = \lambda_0^{-1}$. Then μ_2 is a semi-simple eigenvalue for the pencil $\Delta_2 - \mu_2 \Delta_0$ if and only if λ_0 is a semi-simple eigenvalue for $\Delta_0 - \lambda_0 \Delta_2$.

Lemma 5.7. Suppose that \mathcal{R}_m is the m-th root subspace for (5.10) at an eigenvalue $\boldsymbol{\lambda} = (\lambda_0, \lambda_1)$ (see (1.7)). Then it follows that $(\Delta_2 - \mu_2 \Delta_0) \mathcal{R}_m \subset \Delta_0 \mathcal{R}_{m-1}$. In particular, if \mathcal{B}_m is the basis for \mathcal{R}_m given in Theorem 4.1, then

$$(\Delta_2 - \mu_2 \Delta_0) z_{mj}$$

$$= \sum_{l=1}^{m-1} \left((-1)^l \mu_2^l \sum_{k=1}^l \sum_{g=1}^{d_k} \gamma_{jg}^{mk} \Delta_0 z_{kg} \right),$$

where γ_{iq}^{mk} are defined inductively by

$$\gamma_{jg}^{k+1,k} = c_{1gj}^{2,k,k+1}$$
 and $\gamma_{jg}^{k+r,k} = \sum_{l=k+1}^{k+r-1} \sum_{h=1}^{d_l} c_{1hj}^{2,l,k+r} \gamma_{hg}^{lk}$

for r = 2, 3, ...

Proof. If we multiply the relation $(\Gamma_0 - \lambda_0 I) \mathcal{R}_m \subset \mathcal{R}_{m-1}$ by $-\lambda_0^{-1} \Delta_2$ on the left, we have that $(\Delta_2 - \mu_2 \Delta_0) \mathcal{R}_m \subset \Delta_2 \mathcal{R}_{m-1}$. Then we prove by induction on m that $(\Delta_2 - \mu_2 \Delta_0) \mathcal{R}_m \subset \Delta_0 \mathcal{R}_{m-1}$. The second part of the lemma also follows by induction on m. We omit the details.

Now we compare our results with those of [8]. First we consider the case of a real eigenvalue. We assume that $\lambda \in \mathbb{R}^2$ is an eigenvalue of (5.10), that $z_1 = x_{11} \otimes x_{21}$ is a corresponding eigenvector and that

$$S_2 = \left[\begin{array}{cc} \langle A_{11}x_{11}, x_{11} \rangle & \langle A_{11}x_{11}, x_{11} \rangle \\ \langle A_{20}x_{21}, x_{21} \rangle & \langle A_{21}x_{21}, x_{21} \rangle \end{array} \right].$$

Proposition 5.8. For $\lambda \in \mathbb{R}^2$ the following are equivalent:

- (i) **\(\lambda \)** is a semi-simple eigenvalue,
- (ii) $\langle \Delta_0 z_1, z_1 \rangle \neq 0$,
- (iii) the rank of the matrix S_2 is equal to 2.

Proof. Conditions (i) and (ii) are equivalent by [8, Thm. 6.15] and Corollary 5.6. The equivalence of (i) and (iii) follows from [5, Thm. 6.3]. \Box

Proposition 5.9. For $\lambda \in \mathbb{R}^2$ the following are equivalent:

- (i) λ is not a semi-simple eigenvalue,
- (ii) λ is a nonderogatory eigenvalue, but not algebraically simple,
- (iii) $\langle \Delta_0 z_1, z_1 \rangle = 0$,
- (iv) the rank of the matrix S_2 is less than or equal to 1.

Proof. If rank S_2 is less than 2, then it has to be equal to 1. Namely, it follows from (iii') that $\langle A_{i1}x_{i1}, x_{i1}\rangle \neq 0$ for i=1,2, and therefore $S_2 \neq 0$. Then if λ is not semi-simple, it has to be nonderogatory. The equivalence of conditions (i)-(iv) now follows from Proposition 5.8.

By the above propositions it follows that a basis \mathcal{B} for a root subspace \mathcal{R} at a real eigenvalue λ is described by Theorem 2.2. Then a root subspace \mathcal{M}_{μ_2} (in the notation of [8]) at a real eigenvalue μ_2 is equal to the direct sum of the root subspaces corresponding to the eigenvalues $\lambda = (\lambda_0, \lambda_1)$ with $\lambda_0 = \frac{1}{\mu_2}$. Lemma 5.5 implies that a basis for \mathcal{M}_{μ_2} is thus obtained by Theorem 2.2 as the union of all the bases of the root subspaces $\mathcal{R}(\lambda)$ with $\lambda_0 = \frac{1}{\mu_2}$. This yields Theorems 3.5 and 6.18 of [8] and also the Sturm-Liouville case of [3, Thm. 3.1]. The interested reader will find a numerical implementation of the algorithm for computation of a basis at a nonderogatory eigenvalue in [23].

Assume now that $\lambda = (\lambda_0, \lambda_1)$ is a nonreal eigenvalue of (5.10). Then λ_0 is nonreal [8, p. 99]. We write $z_1 = x_{11} \otimes x_{21}$ for an eigenvector corresponding to λ and

$$S_2 = \left[\begin{array}{cc} \langle A_{10}x_{11}, \overline{x}_{11} \rangle & \langle A_{11}x_{11}, \overline{x}_{11} \rangle \\ \langle A_{20}x_{21}, \overline{x}_{21} \rangle & \langle A_{21}x_{21}, \overline{x}_{21} \rangle \end{array} \right].$$

Proposition 5.10. For a nonreal eigenvalue λ the following are equivalent:

- (i) λ is semi-simple,
- (ii) $\langle \Delta_0 z_1, \overline{z}_1 \rangle \neq 0$,
- (iii) the rank of the matrix S_2 is equal to 2.

Proof. Conditions (i) and (ii) are equivalent by [8, Thm. 6.14] and Corollary 5.6. The equivalence of (i) and (iii) follows from [5, Thm. 6.3].

If λ is not semi-simple and rank S_2 is equal to 1, then λ is nonderogatory and Theorem 2.2 gives a basis for the corresponding root subspace. By Lemma 5.7 it follows that the union of bases, of root subspace at the eigenvalues λ for which $\lambda_0 = \frac{1}{\mu_2}$, is a basis for the root subspace \mathcal{M}_{μ_2} (in notation of [8]). If \mathcal{M}_{μ_2} is a root subspace at a nonreal eigenvalue μ_2 , it is equal to the direct sum of the root subspaces corresponding to the eigenvalues $\lambda = (\lambda_0, \lambda_1)$ with $\lambda_0 = \frac{1}{\mu_2}$. If all these eigenvalues λ are nonderogatory, then Lemma 5.5 implies that a basis for \mathcal{M}_{μ_2} is obtained by Theorem 2.2. This yields Theorems 6.9 and 6.19 of [8] and it also gives a basis for those j in [8, pp. 118–122] for which $s_2 < j \le s_1$. The eigenvalues λ that correspond to the remaining j (i.e., such that $1 \leq j \leq s_2$) are derogatory, i.e., $S_2 = 0$. Now we use Theorem 4.1 to obtain a basis \mathcal{B} for the root subspace at λ . The main result about bases for \mathcal{M}_{μ_2} in [8] is Theorem 6.20. The basis that Faierman gives there is a union of bases of root subspaces $\mathcal{R}(\lambda)$ over all λ with $\lambda_0 = \frac{1}{\mu_2}$. These are of two types: for $s_2 < j \le s_1$ the eigenvalues λ are nonderogatory and our Theorem 2.2 gives the required properties described in [8, Lemma 6.9]. The other eigenvalues with $1 \leq j < s_2$ are derogatory. Bases for $\mathcal{R}(\lambda)$ are described in Conjecture 6.1 of [8], which is established in [8] only when the ascent of Γ_0 is at most 3. Our basis in Theorem 4.1 is constructed in all cases of finite ascent. We hope to discuss the precise relationship with [8, Conj. 6.1] in a separate publication.

Finally, we remark that bases for $\mathcal{R}(\lambda)$ can be used to obtain completeness and expansion results as carried out in [8, Thms. 6.7 and 6.8] and [8, §6.7], respectively.

Then orthonormal bases are important. These can be constructed by a Gram-Schmidt type orthogonalization applied inductively for each m in the construction of a basis for $\mathcal{R}(\lambda)$. Additional care needs to be taken for the root subspaces at pairs of conjugate eigenvalues; for other pairs of eigenvalues the root subspaces are orthogonal by [8, Thm. 6.1].

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