

Variational Bounds to Eigenvalues of Self-Adjoint Eigenvalue Problems with Arbitrary Spectrum

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Abstract. In the present paper a method by Lehmann-Maehly and Goerisch is extended to self-adjoint eigenvalue problems with arbitrary essential spectrum. This extension is obtained by consequently making use of the local character of the method. In this way, upper and lower bounds to all isolated eigenvalues are derived. In our proofs, the close relationship to Wielandt's inverse iteration becomes quite obvious.

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0. Introduction

The characterization of eigenvalues of self-adjoint eigenvalue problems by a minimum-maximum principle for the Rayleigh quotient forms the basis for the famous Rayleigh-Ritz method. As is well known, the method allows for a straightforward and efficient computation of non-increasing upper bounds to eigenvalues below the essential spectrum. In [14] and [15], Lehmann and Maehly independently developed complementary eigenvalue characterizations that provide a possibility to calculate corresponding lower bounds (see also [21: Chapter 4.11]). More precisely, they introduce a spectral parameter and construct minimum-maximum as well as maximum-minimum principles for Temple's quotient. From this, upper and lower bounds are deduced. In particular, the Rayleigh-Ritz method can be interpreted as a special application of the Lehmann-Maehly method. In general, however, the Lehmann-Maehly method is applicable only if certain quantities can be determined explicitly. Especially when treating partial differential equations, difficulties may arise. Of great importance for practical applications is therefore a generalized version by Goerisch (see [5, 8, 10]) that combines high flexibility and good convergence properties. Usually, very satisfying inclusion intervals are obtained from a combination of Rayleigh-Ritz and Goerisch calculations (see [2 - 6, 10, 11]).

Some very interesting eigenvalue problems however possess isolated eigenvalues $\lambda \in (\rho^-, \rho^+)$, where both ρ^- and ρ^+ belong to the essential spectrum. Examples are the

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Dirac equation (see [16: p. 952] and [20]), the Schrödinger operator [1] and the linearized equations of ideal magnetohydrodynamics (see [12: Chapters 1.2 and 3.2] and [17]).

So far, no bounds have been constructed for such interior eigenvalues λ . A Rayleigh-Ritz calculation does not necessarily give bounds. Moreover, it was shown that an additional convergence criterium has to be satisfied in order to guarantee convergence of Rayleigh-Ritz approximations [7, 18]. In problems with interior eigenvalues, a spectral parameter appears to be particularly useful. We shall derive here an extension of the Goerisch method that gives bounds also to interior eigenvalues. In our development we generalize the Lehmann-Maehly variational principle in such a way that an appropriate discretization readily produces Goerisch's method – a result which in itself is at least of theoretical interest.

In Section 1, we characterize eigenvalues of a self-adjoint operator by new variational principles. In what follows, we consider two types of variationally posed eigenvalue problems as there exist two non-equivalent versions of the Goerisch method. In Section 2, we apply results of Section 1 to left-definite eigenvalue problems, Section 3 giving a numerical example for the inclusion of interior eigenvalues. In Section 4, the right-definite theory is developed and applied.

1. Characterizations for eigenvalues of self-adjoint operators

Let \mathbf{H} be a separable, complex Hilbert space with inner product $\langle \cdot, \cdot \rangle$ and norm $\| \cdot \|$, and suppose we are given a self-adjoint operator A in \mathbf{H} with domain $\mathcal{D}(A)$. Denote by $\sigma(A)$ and $\sigma_e(A)$ the spectrum and the essential spectrum of A ,

$$\sigma_e(A) = \left\{ \lambda \in \sigma(A) : \lambda \text{ is an accumulation point of } \sigma(A) \right\} \cup \left\{ \lambda \in \sigma(A) : \lambda \text{ is an eigenvalue of infinite multiplicity} \right\}.$$

If the spectrum $\sigma(A)$ of A is bounded from below and begins with isolated eigenvalues of finite multiplicity,

$$\lambda_1^- \leq \lambda_2^- \leq \dots < \inf \sigma_e(A)$$

(eigenvalues are always repeated according to their multiplicities), these eigenvalues can be characterized by well known variational principles, e.g. the Poincaré principle (cf. [22: Chapter 2])

$$\lambda_i^- = \inf_{\substack{v \in \mathcal{D}(A) \\ \dim V=i}} \max_{0 \neq v \in V} \frac{\langle Av, v \rangle}{\langle v, v \rangle}. \tag{1.1}$$

If in addition $\sigma_e(A) \neq \emptyset$, the number of isolated eigenvalues λ_i^- below $\sigma_e(A)$ is either infinite with $\lim_{i \rightarrow \infty} \lambda_i^- = \inf \sigma_e(A)$ or finite, i.e. equal to k for some $k \in \mathbb{N}$. We then set $\lambda_{k+1}^- = \lambda_{k+2}^- = \dots = \inf \sigma_e(A)$. As a consequence, (1.1) holds for $i \in \mathbb{N}$ in both cases (cf. [21: Chapter 3.5]).

If $\sigma(A)$ is bounded from above and begins with isolated eigenvalues λ_i^+ , $\sup \sigma_e(A) < \dots \leq \lambda_2^+ \leq \lambda_1^+$, an analogous principle is obtained, with “inf” and “max” being replaced by “sup” and “min”, respectively.

Evidently, interior eigenvalues are not accessible for this kind of variational principles. Therefore, a spectral parameter ρ is introduced. For fixed $\rho \in \mathbb{R}$ set

$$\rho^- = \sup \{ \lambda \in \sigma_e(A) : \lambda < \rho \} \in \mathbb{R} \cup \{-\infty\}.$$

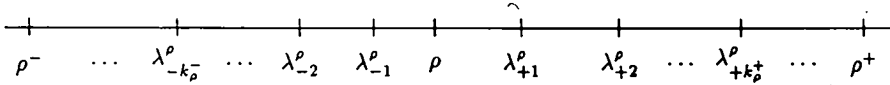
Provided $\rho^- < \rho$ and ρ is no accumulation point of $(\rho^-, \rho) \cap \sigma(A)$, this intersection is either empty or consists of a countable number of eigenvalues of finite multiplicity. It is then possible to introduce the following local notation for these eigenvalues:

$$\rho^- < \dots \leq \lambda_{-k_p^-}^\rho \leq \dots \leq \lambda_{-2}^\rho \leq \lambda_{-1}^\rho < \rho,$$

where $k_p^- \in \mathbb{N}_0$ denotes the number of eigenvalues we are interested in.

The definition of $\rho^+ \in \mathbb{R} \cup \{\infty\}$ and λ_{+i}^ρ for $i = 1, \dots, k_p^+$ is analogous and can thus be omitted.

If $\rho \in \mathbb{R}$ is no accumulation point of $\sigma(A)$ – though it may possibly be an eigenvalue of arbitrary multiplicity of A – we have the following situation:



In applications the boundary points in \mathbb{R} of the essential spectrum are often well known. They represent particularly well suited choices for our spectral parameter ρ (see the examples of Sections 3 and 4).

First, we shall characterize the eigenvalues $\lambda_{\pm i}^\rho$ by three different variational principles which allow to construct lower bounds to λ_{-i}^ρ for $i = 1, \dots, k_p^-$, and upper bounds to λ_{+i}^ρ for $i = 1, \dots, k_p^+$. By an appropriate choice for ρ it is thus possible to reach every isolated eigenvalue of A .

In order to construct variational principles for $\lambda_{\pm i}^\rho$ let

$$D_\rho = \{ u \in \mathcal{D}(A) : Au = \rho u \}$$

and consider the restriction of $A - \rho I$ to

$$\mathcal{D}(A) \ominus D_\rho = \{ u \in \mathcal{D}(A) : \langle u, v \rangle = 0 \text{ for all } v \in D_\rho \}.$$

As D_ρ is a reducing subspace for A , we have the inclusion

$$\sigma(A) \subset \sigma(A|_{\mathcal{D}(A) \ominus D_\rho}) \cup \{ \rho \}.$$

In particular, $(A - \rho I)|_{\mathcal{D}(A) \ominus D_\rho}$ is invertible. We set

$$R_\rho = ((A - \rho I)|_{\mathcal{D}(A) \ominus D_\rho})^{-1}$$

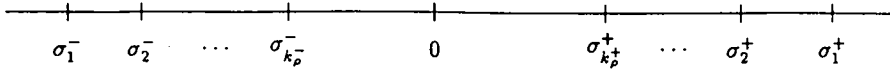
and consider the eigenvalue problem for the self-adjoint operator R_ρ in D_ρ^\perp . Now, $(\lambda, u) \in \mathbb{R} \times \mathcal{D}(A)$ with $\lambda \neq \rho$ is an eigenpair of A if and only if $((\lambda - \rho)^{-1}, u)$ is an eigenpair of R_ρ . Set

$$\sigma_i^- = (\lambda_{-i}^\rho - \rho)^{-1} \quad \text{and} \quad \sigma_i^+ = (\lambda_{+i}^\rho - \rho)^{-1}$$

for $i = 1, \dots, k_\rho^\mp$, so that we arrive at

$$\lambda_{-i}^\rho = \rho + (\sigma_i^-)^{-1} \quad \text{and} \quad \lambda_{+i}^\rho = \rho + (\sigma_i^+)^{-1}, \tag{1.2}$$

respectively.



Poincaré’s principle applied to R_ρ now gives the following characterizations.

Theorem 1.1: *Let $\rho^- < \rho$. Then*

$$\lambda_{-i}^\rho = \rho + \left(\inf_V \max_{0 \neq v \in V} \frac{\langle Av, v \rangle - \rho \langle v, v \rangle}{\|(A - \rho I)v\|^2} \right)^{-1} \tag{1.3}$$

holds for $i = 1, \dots, k_\rho^-$, the infimum being taken over all subspaces V of $\mathcal{D}(A)$ that satisfy

$$\dim V = i \quad \text{and} \quad Av \neq \rho v \text{ for all } 0 \neq v \in V. \tag{1.4}$$

If $\rho < \rho^+$, we have

$$\lambda_{+i}^\rho = \rho + \left(\sup_V \min_{0 \neq v \in V} \frac{\langle Av, v \rangle - \rho \langle v, v \rangle}{\|(A - \rho I)v\|^2} \right)^{-1} \tag{1.5}$$

for $i = 1, \dots, k_\rho^+$ and with V as in (1.4).

Proof: Suppose $k_\rho^- \geq 1$. The spectral mapping theorem [13: Chapter III, Section 6.3] gives

$$\tilde{\sigma}(R_\rho) = (\tilde{\sigma}(A|_{\mathcal{D}(A) \ominus D_\rho}) - \rho)^{-1}$$

where

$$\tilde{\sigma}(S) = \begin{cases} \sigma(S) & \text{if } S \text{ is bounded} \\ \sigma(S) \cup \{\infty\} & \text{otherwise.} \end{cases}$$

The mapping $\lambda \mapsto \sigma = (\lambda - \rho)^{-1}$ transforms the eigenvalues λ_{-i}^ρ ($i = 1, \dots, k_\rho^-$) of A into eigenvalues σ_i^- of R_ρ without changing their multiplicities. The spectrum of R_ρ is bounded from below and begins with the eigenvalues

$$\sigma_1^- \leq \dots \leq \sigma_{k_\rho}^- < \inf \sigma_e(R_\rho).$$

We apply Poincaré’s principle (1.1),

$$\sigma_i^- = \inf_{\substack{U \subset \mathcal{D}(R_\rho) \\ \dim U = i}} \max_{0 \neq u \in U} \frac{\langle R_\rho u, u \rangle}{\langle u, u \rangle} \quad \text{for } i = 1, \dots, k_\rho^-, \tag{1.6}$$

substitute $w = R_\rho u$ and replace $U \subset \mathcal{D}(R_\rho)$ by $W \subset \mathcal{D}(A) \ominus D_\rho$,

$$\sigma_i^- = \inf_{\substack{W \subset \mathcal{D}(A) \ominus D_\rho \\ \dim W = i}} \max_{0 \neq w \in W} \frac{\langle w, (A - \rho I)w \rangle}{\|(A - \rho I)w\|^2} \quad \text{for } i = 1, \dots, k_\rho^-. \tag{1.7}$$

Evidently, every subspace W admissible in (1.7) satisfies (1.4). Conversely, suppose that (1.4) holds for a subspace V of $\mathcal{D}(A)$, and denote by P the orthogonal projector in \mathbf{H} onto the orthogonal complement of D_ρ . Then $W = PV$ is admissible in (1.7), and for $v \in V$ and $w = Pv \in W$ we have

$$(A - \rho I)v = (A - \rho I)w \quad \text{and} \quad \langle v, (A - \rho I)v \rangle = \langle w, (A - \rho I)w \rangle.$$

Thus, we find

$$\sigma_i^- = \inf_{\substack{v \in \mathcal{D}(A) \\ (1.4)}} \max_{0 \neq v \in V} \frac{\langle v, (A - \rho I)v \rangle}{\|(A - \rho I)v\|^2} \quad \text{for } i = 1, \dots, k_\rho^-$$

giving (1.3). Formula (1.5) is obtained in exactly the same way ■

Note that equation (1.6) together with (1.2) already indicates a close relationship to Wielandt's inverse iteration dealing with matrix problems.

Motivated by Goerisch's approach we now generalize the characterisations (1.3) and (1.5) of the eigenvalues in two different ways, in order to enlarge the range of applicability of the derived methods (cf. Remark 2.2).

In addition to our initial assumptions suppose that \mathbf{X} is a complex Hilbert space with inner product $s(\cdot, \cdot)$, and that $T : \mathbf{H} \rightarrow \mathbf{X}$ is an isometry

$$s(Tu, Tv) = \langle u, v \rangle \quad \text{for } u, v \in \mathbf{H}.$$

The choice

$$\mathbf{X} = \mathbf{H}, \quad s(u, v) = \langle u, v \rangle \quad \text{for } u, v \in \mathbf{H}, \quad Tv = v \quad \text{for } v \in \mathbf{H} \quad (1.8)$$

gives the procedure of Lehmann-Maehly (cf. Remark 2.2).

Denote by

$$\mathbf{X}^\circ = \left\{ w \in \mathbf{X} : s(w, Tu) = 0 \quad \text{for } u \in \mathbf{H} \right\}$$

the orthogonal complement of $T\mathbf{H}$ in \mathbf{X} . The quantities \mathbf{X} , s and T provide the following complementary variational principles.

Lemma 1.2: *Suppose S is a self-adjoint operator in \mathbf{H} with domain $\mathcal{D}(S)$. Then for $v \in \mathcal{D}(S)$ given, we have*

$$\begin{aligned} \langle Sv, Sv \rangle &= \max \left\{ \langle Sv, f \rangle + \overline{\langle Sv, f \rangle} - \langle f, f \rangle : f \in \mathbf{H} \right\} \\ \langle Sv, Sv \rangle &= \min \left\{ s(TSv + w^\circ, TSv + w^\circ) : w^\circ \in \mathbf{X}^\circ \right\}. \end{aligned} \quad (1.9)$$

Proof: For $v \in \mathcal{D}(S)$, $f \in \mathbf{H}$ and $z = TSv + w^\circ$ with $w^\circ \in \mathbf{X}^\circ$ one obtains

$$0 \leq \langle Sv - f, Sv - f \rangle = \langle Sv, Sv \rangle - \langle Sv, f \rangle - \overline{\langle Sv, f \rangle} + \langle f, f \rangle$$

and

$$0 \leq s(z - TSv, z - TSv) = s(z, z) - \langle Sv, Sv \rangle$$

which gives the assertion ■

We mention that relation (1.9) still holds when $s(\cdot, \cdot)$ denotes a positive semidefinite sesquilinear form in \mathbf{X} . For simplicity, we restrict our considerations to the positive definite case. It was pointed out in [9] that many well known complementary variational principles can be derived from (1.9) by an appropriate choice of \mathbf{X} , $s(\cdot, \cdot)$ and T , e.g. those given in [19: Chapter 4.1.4].

New characterizations for λ_{\pm}^{ρ} can be deduced by application of (1.9) with $S = A$ and $S = A^{-1}$, respectively.

Theorem 1.3: *Let $\rho^- < \rho$. Then,*

$$\lambda_{-i}^{\rho} = \rho + \left(\inf_{\substack{V \subset \mathcal{D}(A) \\ (1.4)}} \max_{0 \neq v \in V} \min_{w \in TAv + \mathbf{X}^{\circ}} \frac{\langle Av, v \rangle - \rho \langle v, v \rangle}{s(w, w) - 2\rho \langle Av, v \rangle + \rho^2 \langle v, v \rangle} \right)^{-1} \quad (1.10)$$

holds for $i = 1, \dots, k_{\rho}^-$.

Proof: Let $1 \leq i \leq k_{\rho}^-$. We apply Theorem 1.1 and relation (1.9) for $S = A$,

$$\sigma_i^- = \inf_{\substack{V \subset \mathcal{D}(A) \\ (1.4)}} \max_{0 \neq v \in V} \frac{\langle Av, v \rangle - \rho \langle v, v \rangle}{\min_{w \in TAv + \mathbf{X}^{\circ}} \{s(w, w)\} - 2\rho \langle Av, v \rangle + \rho^2 \langle v, v \rangle}. \quad (1.11)$$

Due to $\sigma_i^- < 0$ it suffices to take the infimum over subspaces V of $\mathcal{D}(A)$ satisfying

$$\dim V = i \quad \text{and} \quad \langle Av, v \rangle < \rho \langle v, v \rangle \quad \text{for all } 0 \neq v \in V. \quad (1.12)$$

Inserting this into (1.11) provides

$$\sigma_i^- = \inf_{\substack{V \subset \mathcal{D}(A) \\ (1.4)}} \max_{0 \neq v \in V} \min_{w \in TAv + \mathbf{X}^{\circ}} \frac{\langle Av, v \rangle - \rho \langle v, v \rangle}{s(w, w) - 2\rho \langle Av, v \rangle + \rho^2 \langle v, v \rangle}$$

and thus (1.10) ■

Theorem 1.4: *Suppose $0 \notin \sigma(A)$, $\rho > 0$ and $\rho^- < \rho$. Then we have*

$$\lambda_{-i}^{\rho} = \rho + \rho \left(\inf_V \max_{0 \neq v \in V} \min_{w \in TA^{-1}v + \mathbf{X}^{\circ}} \frac{\langle v, v \rangle - \rho \langle A^{-1}v, v \rangle}{\langle v, v \rangle - 2\rho \langle A^{-1}v, v \rangle + \rho^2 s(w, w)} - 1 \right)^{-1} \quad (1.13)$$

for $1 \leq i \leq k_{\rho}^-$ such that λ_{-i}^{ρ} is positive. Here, V varies over subspaces of \mathbf{H} satisfying

$$\dim V = i \quad \text{and} \quad v \neq \rho A^{-1}v \quad \text{for all } 0 \neq v \in V. \quad (1.14)$$

Proof: Due to our assumptions, $A : \mathcal{D}(A) \rightarrow \mathbf{H}$ is a bijection. We combine Theorem 1.1 and relation (1.9) for $S = A^{-1}$, giving

$$\begin{aligned} \sigma_i^- &= \inf_{\substack{V \subset \mathcal{D}(A) \\ (1.4)}} \max_{0 \neq v \in V} \frac{\langle Av, v \rangle - \rho \langle v, v \rangle}{\langle Av, Av \rangle - 2\rho \langle Av, v \rangle + \rho^2 \langle v, v \rangle} \\ &= \inf_{\substack{V \subset \mathbf{H} \\ (1.14)}} \max_{0 \neq v \in V} \frac{\langle v, A^{-1}v \rangle - \rho \langle A^{-1}v, A^{-1}v \rangle}{\langle v, v \rangle - 2\rho \langle v, A^{-1}v \rangle + \rho^2 \langle A^{-1}v, A^{-1}v \rangle} \\ &= \rho^{-1} \inf_{\substack{V \subset \mathbf{H} \\ (1.14)}} \max_{0 \neq v \in V} \left(\frac{\langle v, v \rangle - \rho \langle v, A^{-1}v \rangle}{\langle v, v \rangle - 2\rho \langle v, A^{-1}v \rangle + \rho^2 \min_{w \in TA^{-1}v + \mathbf{X}^{\circ}} \{s(w, w)\}} - 1 \right). \end{aligned}$$

Here, $\lambda_{-i}^\rho > 0$ implies $\sigma_i^- < -\rho^{-1} < 0$. Thus, we can restrict the infimum to subspaces V of H satisfying

$$\dim V = i \quad \text{and} \quad \langle v, v \rangle < \rho \langle v, A^{-1}v \rangle \quad \text{for all } 0 \neq v \in V. \tag{1.15}$$

Employing (1.15) we find

$$\sigma_i^- = \frac{1}{\rho} \left(\inf_{\substack{V \subset H \\ (1.14)}} \max_{0 \neq v \in V} \min_{w \in TA^{-1}v + X^\circ} \frac{\langle v, v \rangle - \rho \langle v, A^{-1}v \rangle}{\langle v, v \rangle - 2\rho \langle v, A^{-1}v \rangle + \rho^2 s(w, w)} - 1 \right)$$

and the assertion is proved ■

Through a reformulation of Theorems 1.3 and 1.4 one can derive characterizations for λ_{-i}^ρ that involve the following generalized Temple quotients:

$$\frac{s(w, w) - \rho \langle Av, v \rangle}{\langle Av, v \rangle - \rho \langle v, v \rangle} \quad \text{for } v \in \mathcal{D}(A), w \in TAv + X^\circ, \langle v, Av \rangle - \rho \langle v, v \rangle \neq 0 \tag{1.16}$$

$$\frac{\langle v, v \rangle - \rho \langle v, A^{-1}v \rangle}{\langle v, A^{-1}v \rangle - \rho s(w, w)} \quad \text{for } v \in H, w \in TA^{-1}v + X^\circ, \langle v, A^{-1}v \rangle - \rho s(w, w) \neq 0. \tag{1.17}$$

In the situation of (1.8), these yield the classical Temple quotients

$$\frac{\langle Av, (A - \rho I)v \rangle}{\langle v, (A - \rho I)v \rangle} \quad \text{and} \quad \frac{\langle v, (I - \rho A^{-1})v \rangle}{\langle A^{-1}v, (I - \rho A^{-1})v \rangle}$$

Observe that by $\rho \rightarrow \pm\infty$ in (1.16) and $\rho \rightarrow 0$ in (1.17) the Rayleigh quotient of A and the inverse of the Rayleigh quotient of A^{-1} is produced.

Corollary 1.5: *Let $\rho^- < \rho$. Then*

$$\lambda_{-i}^\rho = \sup_{\substack{V \subset \mathcal{D}(A) \\ (1.12)}} \min_{0 \neq v \in V} \max_{w \in TAv + X^\circ} \frac{s(w, w) - \rho \langle Av, v \rangle}{\langle Av, v \rangle - \rho \langle v, v \rangle}$$

holds for $i = 1, \dots, k_\rho^-$.

Corollary 1.6: *Suppose $0 \notin \sigma(A)$, $\rho > 0$ and $\rho^- < \rho$. Then*

$$\lambda_{-i}^\rho = \sup_{\substack{V \subset H \\ (1.15)}} \min_{0 \neq v \in V} \max_{w \in TA^{-1}v + X^\circ} \frac{\langle v, v \rangle - \rho \langle v, A^{-1}v \rangle}{\langle v, A^{-1}v \rangle - \rho s(w, w)} \tag{1.18}$$

holds for $1 \leq i \leq k_\rho^-$ such that λ_{-i}^ρ is positive.

Remark 1.7: In Theorems 1.3 and 1.4, interchanging “inf” and “sup” as well as “min” and “max” gives characterizations for λ_{+i}^ρ , with (1.12) and (1.15) being replaced by

$$\dim V = i \quad \text{and} \quad \langle Av, v \rangle > \rho \langle v, v \rangle \quad \text{for all } 0 \neq v \in V.$$

and

$$\dim V = i \quad \text{and} \quad \langle v, A^{-1}v \rangle > \rho \langle A^{-1}v, A^{-1}v \rangle \quad \text{for all } 0 \neq v \in V,$$

respectively. The same applies to Corollaries 1.5 and 1.6. In Corollary 1.6, however, $w \in TAv + X^\circ$ is then further restricted by

$$\langle v, A^{-1}v \rangle > \rho s(w, w).$$

The proofs follow the exact pattern of the previous ones.

In applications one often finds eigenvalue problems which are given in variational form. In what follows we shall therefore employ the results of this section in order to construct variational bounds to eigenvalues of such problems.

2. The left-definite Goerisch method

Let H_a be a separable, complex Hilbert space with inner product $a(\cdot, \cdot)$ and norm $\|\cdot\|$. Suppose $b(\cdot, \cdot)$ is a continuous, Hermitian, sesquilinear form on $H_a \times H_a$ such that $u \in H_a$ and $b(u, v) = 0$ for all $v \in H_a$ implies $u = 0$. We then consider the left-definite, variationally posed eigenvalue problem

$$\left. \begin{array}{l} \text{Find } \lambda \in \mathbb{R} \text{ and non-trivial } u \in H_a \text{ such} \\ \text{that } a(u, v) = \lambda b(u, v) \text{ holds for all } v \in H_a. \end{array} \right\} \quad (2.1)$$

Denote by B the bounded self-adjoint operator in H_a that satisfies

$$a(Bu, v) = b(u, v) \quad \text{for all } u, v \in H_a.$$

By assumption, B possesses a self-adjoint inverse B^{-1} in H_a . For an application of the results of Section 1, we make the identifications

$$H = H_a, \quad \langle u, v \rangle = a(u, v) \quad \text{for all } u, v \in H_a, \quad A = B^{-1}$$

and we assume that X is a complex Hilbert space with inner product $s(\cdot, \cdot)$ and equipped with an isometry $T : H_a \rightarrow X$. Fix $\rho \in \mathbb{R}$, $\rho > 0$. All assumptions of Section 1 are then satisfied with

$$D_\rho = \left\{ u \in H_a : a(u, v) = \rho b(u, v) \quad \text{for all } v \in H_a \right\}.$$

Evidently, (2.1) is equivalent to the eigenvalue problem for A . Hence, $\sigma(A)$ and $\sigma_e(A)$ represent the spectrum and the essential spectrum of (2.1), with $0 \notin \sigma(A)$. We adopt from Section 1 the definition of ρ^\pm , k_ρ^\pm and $\lambda_{\pm i}^\rho$.

In Theorem 1.4 and Corollary 1.6, we have formulated two equivalent characterizations for λ_{-i}^ρ , each of which is easily discretized. A discretization of (1.18) has the advantage that no transformation of eigenvalues is needed. Furthermore, $\rho \rightarrow 0$ obviously produces the Rayleigh-Ritz method. In general, however, the finite-dimensional eigenvalue problem derived from (1.18) involves two indefinite matrices. For this reason we prefer (1.13) which produces right-definite matrix eigenvalue problems in any case. We thus continue with a reformulation of Theorem 1.4 for left-definite, variationally posed eigenvalue problems.

Corollary 2.1: *Let $\rho^- < \rho$. Then*

$$\lambda_{-i}^\rho = \rho + \rho \left(\inf_{\substack{V \subset \mathbf{H}_a, \dim V = i \\ V \cap D_\rho = \{0\}}} \max_{0 \neq v \in V} \min_w \frac{a(v, v) - \rho b(v, v)}{a(v, v) - 2\rho b(v, v) + \rho^2 s(w, w)} - 1 \right)^{-1} \quad (2.2)$$

holds for $1 \leq i \leq k_\rho^-$ such that λ_{-i}^ρ is positive, the minimum being taken over $w \in \mathbf{X}$ such that

$$s(w, Tu) = b(v, u) \quad \text{for all } u \in \mathbf{H}_a. \quad (2.3)$$

Proof: Since (2.3) is equivalent to the inclusion $w \in TBv + \mathbf{X}^\circ$, the assertion follows immediately from Theorem 1.4 ■

An analogous characterization holds for λ_{+i}^ρ (see Remark 1.7).

For a discretization, suppose $n \in \mathbb{N}$, $m \in \mathbb{N}_0$, and suppose the following:

- (L1) $v_1, \dots, v_n \in \mathbf{H}_a$ are linearly independent.
- (L2) $w_1^*, \dots, w_n^* \in \mathbf{X}$ satisfy $s(w_i^*, Tu) = b(v_i, u)$ for all $u \in \mathbf{H}_a$ ($i = 1, \dots, n$).
- (L3) $w_0^*, \dots, w_m^* \in \mathbf{X}^\circ$, where $w_0^* = 0$ and w_1^*, \dots, w_m^* are linearly independent.

Here v_1, \dots, v_n are Rayleigh-Ritz trial functions. In addition, the Goerisch approach requires w_i^* and w_j^* depending on v_i and the choice of \mathbf{X}, s and T .

Remark 2.2: In the Lehmann-Maehly approach which is characterized by (1.8), one has

$$\mathbf{X} = \mathbf{H}_a, \quad s(\cdot, \cdot) = a(\cdot, \cdot), \quad T = I, \quad \mathbf{X}^\circ = \{0\}$$

so that w_i^* is the unique solution of the linear problem

$$\text{Seek } w \in \mathbf{H}_a \text{ such that } a(w, u) = b(v_i, u) \text{ for all } u \in \mathbf{H}_a. \quad (2.4)$$

Only if this problem is too expensive or even impossible to solve, Goerisch's version is applied. It replaces (2.4) by the linear problem

$$\text{Seek } w \in \mathbf{X} \text{ such that } s(w, Tu) = b(v_i, u) \text{ for all } u \in \mathbf{H}_a. \quad (2.5)$$

Provided $\mathbf{X}^\circ \neq \{0\}$, there exists a manifold of solutions. With X, s and T being chosen appropriately, it is much easier to find a particular solution of problem (2.5) than to solve exactly problem (2.4). Solutions $w_j^* \in \mathbf{X}^\circ$ of the homogeneous equation (2.5) are then used to approximate the residual $Tw_{i,L}^* - w_{i,G}^*$, where $w_{i,L}^*$ and $w_{i,G}^*$ denote solutions of problems (2.4) and (2.5), respectively. Note that the idea of this approach is not to improve Lehmann-Maehly bounds but to enlarge widely the range of applicability of the method.

Set

$$V_n = \text{span}\{v_1, \dots, v_n\} \quad \text{and} \quad \mathbf{X}_m^\circ = \text{span}\{w_0^*, \dots, w_m^*\}.$$

For convenience, we assume

$$V_n \cap D_\rho = \{0\}. \quad (2.6)$$

Whenever (2.6) is violated, one eliminates from V_n all exact eigenvectors corresponding to ρ . Let

$$\begin{aligned}
 a_{ik}^{(1)} &= a(v_k, v_i), & a_{ik}^{(2)} &= b(v_k, v_i), & a_{ik}^{(3)} &= a(Bv_k, Bv_i) & (i, k = 1, \dots, n) \\
 c_{ik}^{(11)} &= s(w_k^*, w_i^*) & (i, k = 1, \dots, n) \\
 c_{ik}^{(12)} &= \begin{cases} s(w_k^\circ, w_i^*) & \text{if } m > 0 \\ 0 & \text{if } m = 0 \end{cases} & (i = 1, \dots, n; k = 1, \dots, \max\{1, m\}) \\
 c_{ik}^{(22)} &= \begin{cases} s(w_k^\circ, w_i^\circ) & \text{if } m > 0 \\ 1 & \text{if } m = 0 \end{cases} & (i, k = 1, \dots, \max\{1, m\}) \\
 A_j &= (a_{ik}^{(j)}) & (j = 1, 2, 3) \\
 C_{jl} &= (c_{ik}^{(jl)}) & (j, l = 1, 2; j \leq l).
 \end{aligned}$$

Lemma 2.3: Let $v = \sum_{i=1}^n x_i v_i$ with $x = (x_1, \dots, x_n)^t \in \mathbf{C}^n$, $w^* = \sum_{i=1}^n x_i w_i^*$, and $w^\circ \in \mathbf{X}_m^\circ$ with

$$s(w^\circ, w) = -s(w^*, w) \quad \text{for all } w \in \mathbf{X}_m^\circ. \tag{2.7}$$

Then

$$\begin{aligned}
 x^H (C_{11} - C_{12} C_{22}^{-1} C_{12}^H) x &= s(w^* + w^\circ, w^* + w^\circ) \\
 &= \min\{s(w^* + w, w^* + w) : w \in \mathbf{X}_m^\circ\} \\
 &\geq x^H A_3 x
 \end{aligned}$$

holds.

Proof: (i) For $m > 0$, there exists $y = (y_1, \dots, y_m)^t \in \mathbf{C}^m$ such that $w^\circ = \sum_{i=1}^m y_i w_i^\circ$. The equation (2.7) is equivalent to $C_{22} y = -C_{12}^H x$, giving

$$s(w^* + w^\circ, w^* + w^\circ) = s(w^* + w^\circ, w^*) = x^H (C_{11} - C_{12} C_{22}^{-1} C_{12}^H) x. \tag{2.8}$$

Trivially, equality holds in (2.8) for $m = 0$.

(ii) For arbitrary $w \in \mathbf{X}_m^\circ$, equation (2.7) gives

$$s(w^* + w, w^* + w) = s(w^* + w^\circ, w^* + w^\circ) + s(w^\circ - w, w^\circ - w).$$

(iii) Since $w^* \in TBv + \mathbf{X}^\circ$, $w^\circ \in \mathbf{X}_m^\circ$, Lemma 1.2 yields

$$s(w^* + w^\circ, w^* + w^\circ) \geq \min\{s(TBv + w, TBv + w) : w \in \mathbf{X}^\circ\} = a(Bv, Bv).$$

Thus the assertion is proved ■

Consider the matrix eigenvalue problem

$$\tau \in \mathbb{R}, x \in \mathbf{C}^n : (A_1 - \rho A_2)x = \tau (A_1 - 2\rho A_2 + \rho^2 (C_{11} - C_{12} C_{22}^{-1} C_{12}^H))x. \tag{2.9}$$

Assumption (2.6) guarantees the positive definiteness of $A_1 - 2\rho A_2 + \rho^2 A_3$. By Lemma 2.3, this implies the existence of n real eigenvalues and corresponding eigenvectors for

problem (2.9). Observe that $C_{11} - C_{12}C_{22}^{-1}C_{12}^H = A_3$ if the Lehmann-Maehly method is applied. Denote all eigenvalues τ of problem (2.9) in $\mathbb{R} \setminus [0, 1]$ by τ_i^- and τ_i^+ and arrange them in the order

$$\tau_1^- \leq \dots \leq \tau_{n^-}^- < 0 < 1 < \tau_{n^+}^+ \leq \dots \leq \tau_1^+$$

with $n^- + n^+ \leq n$. Then we set

$$\begin{aligned} \Lambda_{-i}^{\rho[n,m]} &= \rho + \rho(\tau_i^- - 1)^{-1} && \text{for } i = 1, \dots, n^- \\ \Lambda_{+i}^{\rho[n,m]} &= \rho + \rho(\tau_i^+ - 1)^{-1} && \text{for } i = 1, \dots, n^+. \end{aligned}$$

Theorem 2.4: For $\rho^- < \rho$ and $1 \leq i \leq \min\{n^-, k_\rho^-\}$, we have

$$\Lambda_{-i}^{\rho[n,m]} \leq \lambda_{-i}^\rho. \tag{2.10}$$

For $\rho < \rho^+$ and $1 \leq i \leq \min\{n^+, k_\rho^+\}$, we have

$$\Lambda_{+i}^{\rho[n,m]} \geq \lambda_{+i}^\rho.$$

Proof: We restrict ourselves to the proof of (2.10). With the notation introduced in Lemma 2.3, we find

$$\begin{aligned} 0 > \tau_i^- &= \min_{\substack{V \subset \mathbb{C}^n \\ \dim V=i}} \max_{0 \neq x \in V} \frac{x^H(A_1 - \rho A_2)x}{x^H(A_1 - 2\rho A_2 + \rho^2(C_{11} - C_{12}C_{22}^{-1}C_{12}^H))x} \\ &= \min_{\substack{V \subset V_n \\ \dim V=i}} \max_{0 \neq v \in V} \frac{a(v, v) - \rho b(v, v)}{a(v, v) - 2\rho b(v, v) + \rho^2 s(w^* + w^\circ, w^* + w^\circ)} \\ &= \min_{\substack{V \subset V_n \\ \dim V=i}} \max_{0 \neq v \in V} \min_{w \in w^* + X_m^\circ} \frac{a(v, v) - \rho b(v, v)}{a(v, v) - 2\rho b(v, v) + \rho^2 s(w, w)}. \end{aligned}$$

By construction, we have $V_n \subset H_a$, $V_n \cap D_\rho = \{0\}$ as well as $s(w, Tu) = b(v, u)$ for $u \in H_a$ and $w \in w^* + X_m^\circ$. Thus,

$$\tau_i^- - 1 \geq \inf_{\substack{V \subset H_a, \dim V=i \\ V \cap D_\rho = \{0\}}} \max_{0 \neq v \in V} \min_{(2.3)} \frac{a(v, v) - \rho b(v, v)}{a(v, v) - 2\rho b(v, v) + \rho^2 s(w, w)} - 1.$$

Corollary 2.1 provides 2.10 ■

Remark 2.5: Note that the selfadjoint operator A in H is only of theoretical interest, and that it is not required in the calculation of bounds.

For convergence results, see [24].

3. Numerical example

We shall consider an eigenvalue problem for partial differential equations that was suggested in [7] as a model problem for the study of numerical phenomena in the spectral approximation of the ideal linear magnetohydrodynamics equations ("pollution effect").

Let

$$\Omega = (-1, 1) \times (-1, 1) \subset \mathbb{R}^2 \quad \text{and} \quad \mathbf{H}_\alpha = H_0^1(\Omega) \times (L_2(\Omega) \times L_2(\Omega)).$$

We define

$$a(u, v) = \int_{\Omega} (\alpha \mathbf{grad} u_I \cdot \mathbf{grad} \bar{v}_I + \mathbf{grad} u_I \cdot \bar{\mathbf{v}}_{II} + \mathbf{u}_{II} \cdot \mathbf{grad} \bar{v}_I + u_I \bar{v}_I + 2\mathbf{u}_{II} \cdot \bar{\mathbf{v}}_{II}) d\Omega$$

$$b(u, v) = \int_{\Omega} (u_I \bar{v}_I + \mathbf{u}_{II} \cdot \bar{\mathbf{v}}_{II}) d\Omega$$

for $u = (u_I, \mathbf{u}_{II}), v = (v_I, \mathbf{v}_{II}) \in \mathbf{H}_\alpha$ with $\alpha \in C^\infty(\bar{\Omega}), \alpha(x, y) > 0$ and $(2\alpha - 1)(x, y) > 0$ for $(x, y) \in \Omega$.

The inner product $a(\cdot, \cdot)$ gives a norm in \mathbf{H}_α which is equivalent to the natural norm of this product space. By partial integration one shows that (2.1) is a weak form of the boundary value problem

$$\begin{aligned} -\operatorname{div}(\alpha \mathbf{grad} u_I + \mathbf{u}_{II}) &= (\lambda - 1)u_I && \text{in } \Omega \\ \mathbf{grad} u_I + \mathbf{u}_{II} &= (\lambda - 1)\mathbf{u}_{II} && \text{in } \Omega \\ u_I &= 0 && \text{on } \partial\Omega. \end{aligned}$$

Here, the spectrum is non-negative and depends on the choice of α . In any case, $\lambda = 2$ is an eigenvalue of infinite multiplicity since for each $0 \neq u \in H^2(\Omega)$ an eigenpair of (2.1) is given by

$$(2, (0, (u_y, -u_x))) \in \mathbb{R} \times \mathbf{H}_\alpha.$$

(a) In the Lehmann-Maehly approach (cf. Remark 2.2) we have $\mathbf{X} = \mathbf{H}_\alpha, s(\cdot, \cdot) = a(\cdot, \cdot)$ and $T = I$. Then, for each trial function $v_i \in \mathbf{H}_\alpha$, the equation (2.4) which is a weak form of the following second order boundary value problem has to be solved for $w_i^* = (w_{i,I}^*, \mathbf{w}_{i,II}^*) \in \mathbf{H}_\alpha$:

$$\begin{aligned} -\operatorname{div}(\alpha \mathbf{grad} w_{i,I}^* + \mathbf{w}_{i,II}^*) &= v_{i,I} && \text{in } \Omega \\ \mathbf{grad} w_{i,I}^* + \mathbf{w}_{i,II}^* &= \mathbf{v}_{i,II} && \text{in } \Omega \\ w_{i,I}^* &= 0 && \text{on } \partial\Omega. \end{aligned} \tag{3.1}$$

(b) In Goerisch's approach we set

$$\mathbf{X} = L_2(\Omega) \times (L_2(\Omega))^2 \times (L_2(\Omega))^2$$

$$s(u, v) = \int_{\Omega} (\alpha \mathbf{u}_{III} \cdot \bar{\mathbf{v}}_{III} + \mathbf{u}_{III} \cdot \bar{\mathbf{v}}_{II} + \mathbf{u}_{II} \cdot \bar{\mathbf{v}}_{III} + u_I \bar{v}_I + 2\mathbf{u}_{II} \cdot \bar{\mathbf{v}}_{II}) d\Omega$$

for

$$u = (u_I, u_{II}, u_{III}), v = (v_I, v_{II}, v_{III}) \in X$$

and

$$T : H_a \ni (f_I, f_{II}) \mapsto (f_I, f_{II}, \text{grad } f_I) \in X.$$

The linear equation (2.5) reduces to a weak form of a first order partial differential equation plus an algebraic equation for $w_i^* = (w_{i,I}^*, w_{i,II}^*, w_{i,III}^*) \in X$:

$$\left. \begin{aligned} -\text{div}(\alpha w_{i,III}^* + w_{i,II}^*) + w_{i,I}^* &= v_{i,I} \\ w_{i,III}^* + 2w_{i,II}^* &= v_{i,II} \end{aligned} \right\} \text{ in } \Omega.$$

These are solved by

$$w_i^* = \left(v_{i,I}, \frac{\alpha}{2\alpha - 1} v_{i,II}, \frac{-1}{2\alpha - 1} v_{i,II} \right) \in X.$$

For each $w_j \in H^1(\Omega) \times H^1(\Omega)$ admissible solutions of the homogeneous equation are given by

$$w_j^\circ = \left(\text{div}(2\alpha - 1)w_j, -w_j, 2w_j \right) \in X^\circ. \tag{3.2}$$

In our calculations of bounds we restrict the problem to the symmetry classes (invariant subspace)

$$\begin{aligned} H_a &= H_0^1(\Omega)^{(0,0)} \times \left(L_2(\Omega)^{(1,0)} \times L_2(\Omega)^{(0,1)} \right) \\ X &= L_2(\Omega)^{(0,0)} \times \left(L_2(\Omega)^{(1,0)} \times L_2(\Omega)^{(0,1)} \right) \times \left(L_2(\Omega)^{(1,0)} \times L_2(\Omega)^{(0,1)} \right), \end{aligned}$$

respectively. Here, $L_2(\Omega)^{(a_1, a_2)}$ denotes the subspace of $L_2(\Omega)$ formed of all functions symmetric with respect to the i -th axis, if $a_i = 0$ and antisymmetric with respect to the same axis, if $a_i = 1$ ($i \in \{1, 2\}$).

We use polynomial trial functions $v_1, \dots, v_n \in H_a$, orthogonalized with respect to the inner product $b(\cdot, \cdot)$ and generated via a suitable enumeration by the following functions in H_a :

$$\left. \begin{aligned} (x, y) &\mapsto \left((x^2 - 1)x^{2i}(y^2 - 1)y^{2k}, (0, 0) \right) \\ (x, y) &\mapsto \left(0, (x^{2i+1}y^{2k}, 0) \right) \\ (x, y) &\mapsto \left(0, (0, x^{2i}y^{2k+1}) \right) \end{aligned} \right\} \text{ for } i, k \in \mathbb{N}_0.$$

Furthermore, the solutions $w_1^\circ, \dots, w_m^\circ \in X^\circ$ of the homogeneous equation are constructed according to (3.2) with the following functions in $H^1(\Omega)^{(1,0)} \times H^1(\Omega)^{(0,1)}$:

$$(x, y) \mapsto (x^{2i+1}y^{2k}, 0) \quad \text{and} \quad (x, y) \mapsto (0, x^{2i}y^{2k+1}) \quad \text{for } i, k \in \mathbb{N}_0.$$

Mathematica [23] was employed for the symbolic evaluation of all inner products and for the calculation of all eigenvalues with high precision.

At first consider $\alpha(x, y) = 1$ for $(x, y) \in \Omega$. The corresponding eigenvalue problem can be solved explicitly: A second eigenvalue of infinite multiplicity is given by $\lambda = 1$ with eigenfunctions $(u, -\mathbf{grad} u) \in \mathbf{H}_a$ for every $0 \neq u \in H^2(\Omega) \cap H_0^1(\Omega)$. Let $(\mu_{k,l}, u_{k,l}) \in \mathbb{R} \times H_0^1(\Omega)$ ($k, l \in \mathbb{N}$) be the eigenpairs of the Laplacian operator $-\Delta u = \mu u$ with

$$\mu_{k,l} = (k^2 + l^2) \left(\frac{\pi}{2}\right)^2$$

$$u_{k,l}(x, y) = \sin\left(\frac{k\pi}{2}(x+1)\right) \sin\left(\frac{l\pi}{2}(y+1)\right) \quad \text{for } (x, y) \in \Omega.$$

Then the further eigenpairs of problem (2.1) are given by

$$\left(2 + \mu_{k,l}, (\mu_{k,l} u_{k,l}, \mathbf{grad} u_{k,l})\right) \in \mathbb{R} \times \mathbf{H}_a \quad (k, l \in \mathbb{N}).$$

As the Poincaré principle characterizes $\lambda = 1$ only,

$$1 = \inf_{\substack{V \subset \mathbf{H}_a \\ \dim V = i}} \max_{0 \neq v \in V} \frac{a(v, v)}{b(v, v)} \quad \text{for all } i \in \mathbb{N},$$

the Rayleigh-Ritz approach does not give bounds to eigenvalues in $[2, \infty)$.

From the explicit knowledge of all eigenvalues $\lambda > 2$, we obtain

$$2 < \lambda_1 < 22 < \lambda_2 = \lambda_3 < 40 < \lambda_4 < 60 < \lambda_5 = \lambda_6 < 78 < \lambda_7 = \lambda_8 < 112 < \lambda_9$$

which allows to identify local notations for these eigenvalues with $\rho = 2, 22, 40, 60, 78, 112$. Consequently by Theorem 2.4, the upper [lower] triangle indicated in Table 1 contains upper [lower] bounds to $\lambda_1, \dots, \lambda_8$.

| ρ | $\lambda_1 =$ 6.93480220 | $\lambda_2 =$ 26.6740110 | $\lambda_3 =$ 26.6740110 | $\lambda_4 =$ 46.41322 | $\lambda_5 =$ 66.15243 | $\lambda_6 =$ 66.15243 | $\lambda_7 =$ 85.892 | $\lambda_8 =$ 85.892 |
|--------|-------------------------------------|-------------------------------------|-------------------------------------|-----------------------------------|-----------------------------------|-----------------------------------|---------------------------------|---------------------------------|
| 2 | $\Lambda_{+1}^\rho =$ 6.93480220 | $\Lambda_{+2}^\rho =$ 26.6740112 | $\Lambda_{+3}^\rho =$ 26.6740112 | $\Lambda_{+4}^\rho =$ 46.41562 | $\Lambda_{+5}^\rho =$ 66.15690 | $\Lambda_{+6}^\rho =$ 66.15690 | $\Lambda_{+7}^\rho =$ 86.684 | $\Lambda_{+8}^\rho =$ 86.957 |
| 22 | $\Lambda_{-1}^\rho =$ 6.93480220 | $\Lambda_{+1}^\rho =$ 26.6740138 | $\Lambda_{+2}^\rho =$ 26.6740138 | $\Lambda_{+3}^\rho =$ 46.42534 | $\Lambda_{+4}^\rho =$ 66.17674 | $\Lambda_{+5}^\rho =$ 66.17679 | $\Lambda_{+6}^\rho =$ 88.141 | $\Lambda_{+7}^\rho =$ 88.424 |
| 40 | $\Lambda_{-3}^\rho =$ 6.93480220 | $\Lambda_{-2}^\rho =$ 26.6740098 | $\Lambda_{-1}^\rho =$ 26.6740098 | $\Lambda_{+1}^\rho =$ 46.46094 | $\Lambda_{+2}^\rho =$ 66.20349 | $\Lambda_{+3}^\rho =$ 66.20366 | $\Lambda_{+4}^\rho =$ 89.057 | $\Lambda_{+5}^\rho =$ 89.361 |
| 60 | $\Lambda_{-4}^\rho =$ 6.93480220 | $\Lambda_{-3}^\rho =$ 26.6740105 | $\Lambda_{-2}^\rho =$ 26.6740105 | $\Lambda_{-1}^\rho =$ 46.39284 | $\Lambda_{+1}^\rho =$ 66.39106 | $\Lambda_{+2}^\rho =$ 66.39200 | $\Lambda_{+3}^\rho =$ 91.032 | $\Lambda_{+4}^\rho =$ 91.340 |
| 78 | $\Lambda_{-6}^\rho =$ 6.93480220 | $\Lambda_{-5}^\rho =$ 26.6740106 | $\Lambda_{-4}^\rho =$ 26.6740106 | $\Lambda_{-3}^\rho =$ 46.40672 | $\Lambda_{-2}^\rho =$ 66.02755 | $\Lambda_{-1}^\rho =$ 66.02806 | $\Lambda_{+1}^\rho =$ 99.496 | $\Lambda_{+2}^\rho =$ 99.667 |
| 112 | $\Lambda_{-8}^\rho =$ 6.93480220 | $\Lambda_{-7}^\rho =$ 26.6740108 | $\Lambda_{-6}^\rho =$ 26.6740108 | $\Lambda_{-5}^\rho =$ 46.41070 | $\Lambda_{-4}^\rho =$ 66.12991 | $\Lambda_{-3}^\rho =$ 66.12996 | $\Lambda_{-2}^\rho =$ 82.463 | $\Lambda_{-1}^\rho =$ 82.881 |

Table 1: $\Lambda_{+i}^\rho = \Lambda_{+i}^{\rho[72,56]}$ for $\rho = 2$ and $\Lambda_{\pm i}^\rho = \Lambda_{\pm i}^{\rho[108,56]}$ for $\rho \neq 2$; $\alpha(x, y) = 1$.

We will continue our example setting $\alpha(x, y) = 1 + x^2 y^2$ for $(x, y) \in \Omega$ at the end of the next section.

4. The Right-definite Goerisch method

Let \mathbf{H}_b be a separable, complex Hilbert space with inner product $b(\cdot, \cdot)$, and suppose that $a(\cdot, \cdot)$ is a Hermitian sesquilinear form in \mathbf{H}_b with domain \mathbf{H}_a . Furthermore, $a(\cdot, \cdot)$ is assumed to be bounded from below and closed. We then consider the right-definite eigenvalue problem

$$\left. \begin{aligned} &\text{Seek } \lambda \in \mathbb{R} \text{ and non-trivial } u \in \mathbf{H}_a \text{ such} \\ &\text{that } a(u, v) = \lambda b(u, v) \text{ holds for all } v \in \mathbf{H}_a. \end{aligned} \right\} \quad (4.1)$$

Problem (4.1) is equivalent to the eigenvalue problem for the self-adjoint operator A defined by

$$\mathcal{D}(A) = \left\{ u \in \mathbf{H}_a : a(u, v) = b(\hat{u}, v) \text{ (} v \in \mathbf{H}_b \text{) for some } \hat{u} \in \mathbf{H}_b \right\}$$

$$b(Au, v) = a(u, v) \text{ for all } u \in \mathcal{D}(A), v \in \mathbf{H}_a$$

(see, e.g., [13: Chapter VI, Section 2]). The operator A is bounded from below. With

$$\mathbf{H} = \mathbf{H}_b \quad \text{and} \quad \langle u, v \rangle = b(u, v) \text{ for all } u, v \in \mathbf{H}_b$$

the operator A satisfies the assumptions of Section 1, where

$$D_\rho = \left\{ u \in \mathbf{H}_a : a(u, v) = \rho b(u, v) \text{ for all } v \in \mathbf{H}_a \right\}.$$

Again, we assume that \mathbf{X} is a complex Hilbert space with inner product $s(\cdot, \cdot)$ and isometry $T : \mathbf{H}_b \rightarrow \mathbf{X}$. For fixed $\rho \in \mathbb{R}$, application of Theorem 1.3 gives the following characterization for λ_{-i}^ρ .

Corollary 4.1: *Set*

$$\mathbf{H}_* = \left\{ v \in \mathbf{H}_a : a(v, u) = s(w, Tu) \text{ (} u \in \mathbf{H}_a \text{) for some } w \in \mathbf{X} \right\}.$$

Let $\rho^- < \rho$ and $1 \leq i \leq k_\rho^-$. Then $\mathbf{H}_* = \mathcal{D}(A)$ and

$$\lambda_{-i}^\rho = \rho + \left(\inf_{\substack{v \in \mathbf{H}_* \\ v \cap D_\rho = \{0\}}} \max_{0 \neq v \in V} \min_w \frac{a(v, v) - \rho b(v, v)}{s(w, w) - 2\rho a(v, v) + \rho^2 b(v, v)} \right)^{-1}, \quad (4.2)$$

the minimum being taken over $w \in \mathbf{X}$ such that

$$s(w, Tu) = a(v, u) \quad \text{for all } u \in \mathbf{H}_a. \quad (4.3)$$

Proof: Evidently, $\mathcal{D}(A) \subset \mathbf{H}_*$. To show the inclusion $\mathbf{H}_* \subset \mathcal{D}(A)$, suppose $v \in \mathbf{H}_*$ and let $w \in \mathbf{X}$ be such that condition (4.3) holds. Since $T\mathbf{H}_b$ is a closed subset of \mathbf{X} it possesses an orthogonal projector. Denote by \hat{w} the orthogonal projection of w onto $T\mathbf{H}_b$, $\hat{w} = T\hat{u}$ for some $\hat{u} \in \mathbf{H}_b$. From $s(w - \hat{w}, Tu) = 0$ for all $u \in \mathbf{H}_a \subset \mathbf{H}_b$ then there follows

$$a(v, u) = s(\hat{w}, Tu) = b(\hat{u}, u) \quad \text{for all } u \in \mathbf{H}_a$$

giving the inclusion $v \in \mathcal{D}(A)$. Assertion (4.2) is now an immediate consequence of Theorem 1.3, since the inclusion $w \in TAv + (T\mathbf{H}_b)^\perp$ is equivalent to the equality

$$s(w, Tu) = s(TAv, Tu) \quad \text{for all } u \in \mathbf{H}_b = \overline{\mathbf{H}_a},$$

and also to condition (4.3) ■

We sketch briefly the discretization of (4.2), which is of course analogous to the discretization of (2.2). Suppose $n \in \mathbb{N}$, $m \in \mathbb{N}_0$, and suppose the following:

- (R1) $v_1, \dots, v_n \in \mathcal{D}(A) \subset \mathbf{H}_a$ are linearly independent.
- (R2) $w_1^*, \dots, w_n^* \in \mathbf{X}$ satisfy $s(w_i^*, Tu) = a(v_i, u)$ for all $u \in \mathbf{H}_a$ ($i = 1, \dots, n$).
- (R3) $w_0^\circ, \dots, w_m^\circ \in \mathbf{X}^\circ$, where $w_0^\circ = 0$ and $w_1^\circ, \dots, w_m^\circ$ are linearly independent.

Again, the classical Lehmann-Maehly method is obtained by (1.8), choosing $\mathbf{X} = \mathbf{H}_b$, $s(\cdot, \cdot) = b(\cdot, \cdot)$ and $T = I$. Set

$$\begin{aligned}
 V_n &= \text{span}\{v_1, \dots, v_n\}, & \mathbf{X}_m^\circ &= \text{span}\{w_0^\circ, \dots, w_m^\circ\} \\
 a_{ik}^{(0)} &= b(Av_k, Av_i), & a_{ik}^{(1)} &= a(v_k, v_i) & a_{ik}^{(2)} &= b(v_k, v_i) & (i, k = 1, \dots, n) \\
 b_{ik}^{(11)} &= s(w_k^*, w_i^*) & (i, k = 1, \dots, n) \\
 b_{ik}^{(12)} &= \begin{cases} s(w_k^\circ, w_i^*) & \text{if } m > 0 \\ 0 & \text{if } m = 0 \end{cases} & (i = 1, \dots, n, k = 1, \dots, \max\{1, m\}) \\
 b_{ik}^{(22)} &= \begin{cases} s(w_k^\circ, w_i^\circ) & \text{if } m > 0 \\ 1 & \text{if } m = 0 \end{cases} & (i, k = 1, \dots, \max\{1, m\}) \\
 A_j &= (a_{ik}^{(j)}) & (j = 0, 1, 2) \\
 B_{jl} &= (b_{ik}^{(jl)}) & (j, l = 1, 2; j \leq l).
 \end{aligned}$$

As in Section 2 we suppose $V_n \cap D_\rho = \{0\}$. Here, we have

$$s(w^* + w^\circ, w^* + w^\circ) = x^H (B_{11} - B_{12}B_{22}^{-1}B_{12}^H)x \geq x^H A_0 x$$

for

$$v = \sum_{i=1}^n x_i v_i \quad \text{with } x = (x_1, \dots, x_n)^t \in \mathbf{C}^n \quad \text{and} \quad w^* = \sum_{i=1}^n x_i w_i^*$$

and $w^\circ \in \mathbf{X}_m^\circ$ defined by

$$s(w^\circ, w) = -s(w^*, w) \quad \text{for all } w \in \mathbf{X}_m^\circ.$$

Consider the right-definite matrix eigenvalue problem

$$\tau \in \mathbb{R}, x \in \mathbf{C}^n: \quad (A_1 - \rho A_2)x = \tau (B_{11} - B_{12}B_{22}^{-1}B_{12}^H - 2\rho A_1 + \rho^2 A_2)x.$$

Denote its non-zero eigenvalues by τ_i^- and τ_i^+ ,

$$\tau_1^- \leq \dots \leq \tau_{n^-}^- < 0 < \tau_{n^+}^+ \leq \dots \leq \tau_1^+$$

$n^- + n^+ \leq n$, and set

$$\begin{aligned}
 \Lambda_{-i}^{\rho[n,m]} &= \rho + (\tau_i^-)^{-1} \quad \text{for all } i = 1, \dots, n^- \\
 \Lambda_{+i}^{\rho[n,m]} &= \rho + (\tau_i^+)^{-1} \quad \text{for all } i = 1, \dots, n^+.
 \end{aligned}$$

Theorem 4.2: *Let $\rho^- < \rho$ or $\rho < \rho^+$. Then, we have*

$$\begin{aligned} \Lambda_{-i}^{\rho[n,m]} &\leq \lambda_{-i}^\rho \quad \text{for } i = 1, \dots, \min\{n^-, k_\rho^-\} \\ \Lambda_{+i}^{\rho[n,m]} &\geq \lambda_{+i}^\rho \quad \text{for } i = 1, \dots, \min\{n^+, k_\rho^+\}, \end{aligned} \tag{4.4}$$

respectively.

Proof: We will prove the first inequality of (4.4) only. Here, one has

$$\begin{aligned} \tau_i^- &= \inf_{\substack{V \subset \mathbb{C}^n \\ \dim V = i}} \max_{0 \neq v \in V} \frac{x^H(A_1 - \rho A_2)x}{x^H(B_{11} - B_{12}B_{22}^{-1}B_{12}^H - 2\rho A_1 + \rho^2 A_2)x} \\ &= \inf_{\substack{V \subset V_n \\ \dim V = i}} \max_{0 \neq v \in V} \min_{w \in w^* + X_m^\circ} \frac{a(v, v) - \rho b(v, v)}{s(w, w) - 2\rho a(v, v) + \rho^2 b(v, v)} \end{aligned}$$

with $V_n \subset \mathcal{D}(A)$ and $w^* + X_m^\circ \subset TA v + X^\circ$. Theorem 1.3 yields

$$\Lambda_{-i}^{\rho[n,m]} = \rho + (\tau_i^-)^{-1} \leq \rho + (\sigma_i^-)^{-1} = \lambda_{-i}^\rho$$

and the assertion is proved ■

According to Theorem 4.2, the bound $\Lambda_{\pm i}^{\rho[n,m]}$ is obtained through a shift and inversion process resembling Wielandt's inverse iteration. Here, only one step of the iteration is applied. Instead, the numbers n, m are increased in order to achieve sufficient accuracy. In particular, the Lehmann-Maehly method for right-definite problems can be regarded as one step of Wielandt's inverse iteration for the subspace $V_n \subset \mathcal{D}(A)$. Consequently, Wielandt's eigenvalue approximations are in fact bounds to eigenvalues.

In order to compare the right- and left-definite methods, assume that problem (4.1) satisfies also the assumptions of Section 2. Then, both versions can be applied, each of them involving the Rayleigh-Ritz matrices $A_1 = (a(v_k, v_i))_{ik}$ and $A_2 = (b(v_k, v_i))_{ik}$. In addition,

$$\begin{aligned} B_{11} - B_{12}B_{22}^{-1}B_{12}^H &\simeq A_0 = (b(Av_k, Av_i))_{ik} \\ C_{11} - C_{12}C_{22}^{-1}C_{12}^H &\simeq A_3 = (a(Bv_k, Bv_i))_{ik} \end{aligned}$$

are required in the respective approaches, $B = A^{-1}$.

Consider e.g. the eigenvalue problem given in Section 3. With $v_i \in \mathcal{D}(A)$ the elements Av_i can be obtained through differentiation,

$$Av_i = \left(-\operatorname{div}(\alpha \operatorname{grad} v_{i,I} + \mathbf{v}_{i,II}) + v_{i,I}, \operatorname{grad} v_{i,I} + 2\mathbf{v}_{i,II} \right) \tag{4.5}$$

whereas Bv_i is the (weak) solution of a boundary value problem, namely (3.1). When applying the right-definite version difficulties mostly arise in finding trial functions v_i which are sufficiently smooth and fulfill all boundary conditions.

Continuing the numerical example of Section 3, we set $\alpha(x, y) = 1 + x^2 y^2$ for $(x, y) \in \Omega$. In this case, there exists essential spectrum $\sigma_e = [1, \frac{3}{2}] \cup \{2\}$ and point spectrum in $(2, +\infty)$. Denote the i -th eigenvalue above 2 by λ_i . Since the operator A in \mathbf{H} is well known the right-definite Lehmann-Maehly method can be applied to the trial functions

$v_i \in \mathcal{D}(A)$ of Section 3 and the elements given in (4.5). Again, $\rho = 2 = \max \sigma_e$ provides upper bounds to λ_i . We then successively choose $\rho \approx 0.9 \Lambda_{\pm i}^{2[56,0]}$ for $i = 2, 4, 5, 7, 9$. Under the heuristic assumption that

$$2 < \lambda_1 < 25 < \lambda_2 \leq \lambda_3 < 46 < \lambda_4 < 62 < \lambda_5 \leq \lambda_6 < 82 < \lambda_7 \leq \lambda_8 < 117 < \lambda_9$$

Theorem 4.2 gives upper [lower] bounds to $\lambda_1, \dots, \lambda_8$ shown in the upper [lower] triangle of Table 2.

| ρ | λ_1 | λ_2 | λ_3 | λ_4 | λ_5 | λ_6 | λ_7 | λ_8 |
|--------|-----------------------------------|---------------------------------|---------------------------------|-------------------------------|-------------------------------|-------------------------------|------------------------------|------------------------------|
| 2 | $\Lambda_{+1}^\rho = 7.206691786$ | $\Lambda_{+2}^\rho = 27.473166$ | $\Lambda_{+3}^\rho = 28.305436$ | $\Lambda_{+4}^\rho = 50.5920$ | $\Lambda_{+5}^\rho = 68.6623$ | $\Lambda_{+6}^\rho = 68.7772$ | $\Lambda_{+7}^\rho = 91.718$ | $\Lambda_{+8}^\rho = 96.518$ |
| 25 | $\Lambda_{-1}^\rho = 7.206691735$ | $\Lambda_{+1}^\rho = 27.473373$ | $\Lambda_{+2}^\rho = 28.305513$ | $\Lambda_{+3}^\rho = 50.5943$ | $\Lambda_{+4}^\rho = 68.6639$ | $\Lambda_{+5}^\rho = 68.7786$ | $\Lambda_{+6}^\rho = 91.811$ | $\Lambda_{+7}^\rho = 96.699$ |
| 46 | $\Lambda_{-3}^\rho = 7.206691744$ | $\Lambda_{-2}^\rho = 27.473124$ | $\Lambda_{-1}^\rho = 28.305411$ | $\Lambda_{+1}^\rho = 50.6110$ | $\Lambda_{+2}^\rho = 68.6672$ | $\Lambda_{+3}^\rho = 68.7814$ | $\Lambda_{+4}^\rho = 91.896$ | $\Lambda_{+5}^\rho = 96.896$ |
| 62 | $\Lambda_{-4}^\rho = 7.206691747$ | $\Lambda_{-3}^\rho = 27.473138$ | $\Lambda_{-2}^\rho = 28.305419$ | $\Lambda_{-1}^\rho = 50.5824$ | $\Lambda_{+1}^\rho = 68.6833$ | $\Lambda_{+2}^\rho = 68.7949$ | $\Lambda_{+3}^\rho = 92.036$ | $\Lambda_{+4}^\rho = 97.186$ |
| 82 | $\Lambda_{-6}^\rho = 7.206691749$ | $\Lambda_{-5}^\rho = 27.473144$ | $\Lambda_{-4}^\rho = 28.305422$ | $\Lambda_{-3}^\rho = 50.5877$ | $\Lambda_{-2}^\rho = 68.6493$ | $\Lambda_{-1}^\rho = 68.7658$ | $\Lambda_{+1}^\rho = 92.844$ | $\Lambda_{+2}^\rho = 98.844$ |
| 117 | $\Lambda_{-8}^\rho = 7.206691750$ | $\Lambda_{-7}^\rho = 27.473148$ | $\Lambda_{-6}^\rho = 28.305424$ | $\Lambda_{-5}^\rho = 50.5894$ | $\Lambda_{-4}^\rho = 68.6573$ | $\Lambda_{-3}^\rho = 68.7730$ | $\Lambda_{-2}^\rho = 91.215$ | $\Lambda_{-1}^\rho = 94.973$ |

Table 2: $\Lambda_{\pm i}^\rho = \Lambda_{\pm i}^{\rho[56,0]}$ for $\rho = 2$ and $\Lambda_{\pm i}^\rho = \Lambda_{\pm i}^{\rho[84,0]}$ for $\rho \neq 2$; $\alpha(x, y) = 1 + x^2y^2$; $\Lambda_{+9}^{2[56,0]} = 130.32$.

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