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2009

MIMS EPrint: 2007.117

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ISSN 1749-9097
DETECTING AND SOLVING HYPERBOLIC QUADRATIC EIGENVALUE PROBLEMS∗

CHUN-HUA GUO†, NICHOLAS J. HIGHAM‡, AND FRANÇOISE TISSEUR‡

Abstract. Hyperbolic quadratic matrix polynomials $Q(\lambda) = \lambda^2 A + \lambda B + C$ are an important class of Hermitian matrix polynomials with real eigenvalues, among which the overdamped quadratics are those with nonpositive eigenvalues. Neither the definition of overdamped nor any of the standard characterizations provides an efficient way to test if a given $Q$ has this property. We show that a quadratically convergent matrix iteration based on cyclic reduction, previously studied by Guo and Lancaster, provides necessary and sufficient conditions for $Q$ to be overdamped. For weakly overdamped $Q$ the iteration is shown to be generically linearly convergent with constant at worst $1/2$, which implies that the convergence of the iteration is reasonably fast in almost all cases of practical interest. We show that the matrix iteration can be implemented in such a way that when overdamping is detected a scalar $\mu < 0$ is provided that lies in the gap between the $n$ largest and $n$ smallest eigenvalues of the $n \times n$ quadratic eigenvalue problem (QEP) $Q(\lambda)x = 0$. Once such a $\mu$ is known, the QEP can be solved by linearizing to a definite pencil that can be reduced, using already available Cholesky factorizations, to a standard Hermitian eigenproblem. By incorporating an initial preprocessing stage that shifts a hyperbolic $Q$ so that it is overdamped, we obtain an efficient algorithm that identifies and solves a hyperbolic or overdamped QEP maintaining symmetry throughout and guaranteeing real computed eigenvalues.

Key words. quadratic eigenvalue problem, hyperbolic, overdamped, weakly overdamped, quadratic matrix polynomial, quadratic matrix equation, solvent, cyclic reduction, doubling algorithm

AMS subject classifications. 15A18, 15A24, 65F15, 65F30

DOI. 10.1137/070704058

1. Introduction. The quadratic eigenvalue problem (QEP) is to find scalars $\lambda$ and nonzero vectors $x$ and $y$ satisfying $Q(\lambda)x = 0$ and $y^*Q(\lambda) = 0$, where

$$Q(\lambda) = \lambda^2 A + \lambda B + C, \quad A, B, C \in \mathbb{C}^{n \times n}$$

is a quadratic matrix polynomial. The vectors $x$ and $y$ are right and left eigenvectors, respectively, corresponding to the eigenvalue $\lambda$. The many applications of the QEP, as well as its theory and algorithms for solving it, are surveyed by Tisseur and Meerbergen [27].

Our interest in this work is in Hermitian quadratic matrix polynomials—those with Hermitian $A$, $B$, and $C$—and more specifically those that are hyperbolic. Hyperbolic quadratics, and the subclass of overdamped quadratics, are defined as follows. For Hermitian $X$ and $Y$ we write $X > Y$ ($X \geq Y$) if $X - Y$ is positive definite (positive semidefinite).

∗Received by the editors October 1, 2007; accepted for publication (in revised form) by Q. Ye September 15, 2008; published electronically January 16, 2009.
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Definition 1.1. $Q(\lambda)$ is hyperbolic if $A$, $B$, and $C$ are Hermitian, $A > 0$, and

$$\left(x^* B x\right)^2 > 4(x^* A x)(x^* C x)$$

for all nonzero $x \in \mathbb{C}^n$.  

Definition 1.2. $Q(\lambda)$ is overdamped if it is hyperbolic with $B > 0$ and $C \geq 0$.

Overdamped quadratics arise in overdamped systems in structural mechanics [22, section 7.6].

Any eigenpair of $Q$ satisfies $x^* Q(\lambda) x = 0$ and hence

$$\lambda = \frac{-x^* B x \pm \sqrt{(x^* B x)^2 - 4(x^* A x)(x^* C x)}}{2x^* A x}.$$  

Therefore the eigenvalues of a hyperbolic $Q$ are real and those of an overdamped $Q$ are real and nonpositive. Both classes of quadratics have other important spectral properties, which we summarize in section 2.

We have two aims. The first is to devise an efficient and reliable numerical test for hyperbolicity or overdamping of a given Hermitian quadratic. The second aim is to build upon an affirmative test result an efficient algorithm for solving the QEP that exploits hyperbolicity and in particular that guarantees real computed eigenvalues in floating point arithmetic.

Part of the motivation for testing overdamping concerns the stability of gyroscopic systems. It is known that a gyroscopic system $G(\lambda) = \lambda^2 A_g + \lambda B_g + C_g$ with $A_g, C_g > 0$ and $B_g$ Hermitian indefinite and nonsingular is stable whenever the quadratic $Q_g(\lambda) = \lambda^2 A_g + \lambda |B_g| + C_g$ is overdamped [9]. Here $|B_g|$ is the Hermitian positive definite square root of $B_g^2$ (i.e., the Hermitian polar factor of the Hermitian matrix $B_g$) [12].

Checking the hyperbolicity condition (1.2) is a nontrivial task, and plausible sufficient conditions for hyperbolicity may be incorrect. For example, it is claimed in [21] that when $A = I$, $B > 0$, and $C \geq 0$, $Q$ is hyperbolic if $B > 2C^{1/2}$. That this claim is false has been shown by Barkwell and Lancaster [1].

Guo and Lancaster [9] propose testing overdamping by using a matrix iteration based on cyclic reduction to compute two solvents (solutions) of the quadratic matrix equation

$$Q(X) = AX^2 + BX + C = 0$$

and then computing an extremal eigenvalue of each solvent. A definiteness test on $Q$ evaluated at the average of the two extremal eigenvalues finally determines whether $Q$ is overdamped. We show that the same iteration can be used to test overdamping in a much more efficient way that does not necessarily require the iteration to be run to convergence, even for a positive test result. Our test is based on a more complete understanding of the behavior of the matrix iteration, developed in section 3.

In section 4 we extend the convergence analysis to weakly overdamped quadratics, for which the strict inequality in (1.2) is replaced by a weak inequality ($\geq$). The key idea is to use an interpretation of the matrix iteration as a doubling algorithm. We show that for weakly overdamped $Q$ with equality in (1.2) for some nonzero $x$, the iteration is linearly convergent with constant at worst $1/2$ in the generic case. A reasonable speed of convergence can therefore be expected in almost all practically important cases.

In section 5 we turn to algorithmic matters. We first show how a hyperbolic $Q$ can be shifted to make it overdamped. Then we specify our test for overdamping, which
requires only the building blocks of Cholesky factorization, matrix multiplication, and the solution of triangular systems. We then show how after a successful test the eigensystem of an overdamped $Q$ can be efficiently computed in a way that exploits the symmetry and definiteness and guarantees real computed eigenvalues.

Veselić [28] and Higham, Tisseur, and Van Dooren [19] have previously shown that every hyperbolic quadratic can be reformulated as a definite pencil $L(\lambda) = \lambda X + Y$ of twice the dimension, and this connection is explored in detail and in more generality by Higham, Mackey, and Tisseur [16]. However, the algorithm developed here is the first practical procedure for arranging that $X$ or $Y$ is a definite matrix and hence allowing symmetry and definiteness to be fully exploited.

Section 6 concludes the paper with a numerical experiment that provides further insight into the theory and algorithms.

2. Preliminaries. We first recall the definition of a definite pencil.

**Definition 2.1.** A Hermitian pencil $L(\lambda) = \lambda X + Y$ is definite (or, equivalently, the matrices $X, Y$ form a definite pair) if $(z^*Xz)^2 + (z^*Yz)^2 > 0$ for all nonzero $z \in \mathbb{C}^n$.

Definite pairs have the desirable properties that they are simultaneously diagonalizable under congruence and, in the associated eigenproblem $L(\lambda)x = 0$, the eigenvalues are real and semisimple.

The next result gives three conditions each equivalent to the condition (1.2) in the definition of hyperbolic quadratic.

**Theorem 2.2.** Let the $n \times n$ quadratic $Q(\lambda) = \lambda^2 A + \lambda B + C$ be Hermitian with $A > 0$ and let

\[
\gamma = \min_{\|x\|_2 = 1} [(x^*Bx)^2 - 4(x^*Ax)(x^*Cx)].
\]

The following statements are equivalent:

(a) $Q$ is hyperbolic.

(b) $\gamma > 0$.

(c) $x'Q(\lambda)x = 0$ has two distinct real zeros for all nonzero $x \in \mathbb{C}^n$.

(d) $Q(\mu) < 0$ for some $\mu \in \mathbb{R}$.

Proof. (a) $\iff$ (b) $\iff$ (c) is immediate. (c) $\iff$ (d) follows from Markus [25, Lemma 31.15].

Hyperbolic quadratics have many interesting properties [25, section 31].

**Theorem 2.3.** Let the $n \times n$ quadratic $Q(\lambda) = \lambda^2 A + \lambda B + C$ be hyperbolic.

(a) The $2n$ eigenvalues of $Q(\lambda)$ are all real and semisimple.

(b) There is a gap between the $n$ largest and $n$ smallest eigenvalues, that is, the eigenvalues can be ordered $\lambda_1 \geq \cdots \geq \lambda_n > \lambda_{n+1} \geq \cdots \geq \lambda_{2n}$.

(c) $Q(\mu) < 0$ for all $\mu \in (\lambda_{n+1}, \lambda_n)$ and $Q(\mu) > 0$ for all $\mu \in (-\infty, \lambda_{2n}) \cup (\lambda_1, \infty)$.

(d) There are $n$ linearly independent eigenvectors associated with the $n$ largest eigenvalues and likewise for the $n$ smallest eigenvalues.

(e) The quadratic matrix equation $Q(X) = 0$ in (1.4) has a solvent $S^{(1)}$ with eigenvalues $\lambda_1, \ldots, \lambda_n$ and a solvent $S^{(2)}$ with eigenvalues $\lambda_{n+1}, \ldots, \lambda_{2n}$. Moreover,

$Q(\lambda) = (I - S^{(2)*})A(I - S^{(1)}) = (I - S^{(1)*})A(I - S^{(2)})$.

The $n$ largest eigenvalues of a hyperbolic quadratic are called the primary eigenvalues and the $n$ smallest eigenvalues are the secondary eigenvalues. The solvents

\footnote{An eigenvalue of a matrix polynomial $P(\lambda) = \sum_{k=0}^f \lambda^k P_k$ is semisimple if it appears only in $1 \times 1$ Jordan blocks in a Jordan triple for $P$ [7].}
$S^{(1)}$ and $S^{(2)}$ having as their eigenvalues the primary eigenvalues and the secondary eigenvalues, respectively, are referred to as the primary and secondary solvents.

Hyperbolicity can also be defined for matrix polynomials $P$ of arbitrary degree [25, section 31]. The notion has recently been extended in [16] by replacing the assumption of a positive definite leading coefficient matrix with $P(\omega) > 0$ for some $\omega \in \mathbb{R} \cup \{\infty\}$.

The next result gives some characterizations of an overdamped quadratic. First, we need a simple lemma.

**Lemma 2.4.** Let $Q(\lambda) = \lambda^2 A + \lambda B + C$ be Hermitian and let $\mu > 0$. Then $Q(-\mu) < 0$ if and only if $B > \mu A + \mu^{-1} C$.

**Proof.** The proof is immediate from $Q(-\mu) = \mu^2 A - \mu B + C < 0 \iff \mu A - B + \mu^{-1} C < 0$. \qed

**Theorem 2.5.** Let $Q(\lambda) = \lambda^2 A + \lambda B + C$ be Hermitian with $A > 0$. Then the following statements are equivalent:

(a) $Q(\lambda)$ is overdamped.

(b) $Q(\lambda)$ is hyperbolic and all of its eigenvalues are real and nonpositive.

(c) $B > 0$, $C \geq 0$, and $B > \mu A + \mu^{-1} C$ for some $\mu > 0$.

**Proof.** (a) $\iff$ (b) is proved in [9, Theorem 5]. (b) $\Rightarrow$ (c): By Theorem 2.3(c), $Q(\mu) < 0$ for some $\mu < 0$; (c) $\Rightarrow$ (a): $B > \mu A + \mu^{-1} C$ with $\mu > 0$ implies $Q(-\mu) < 0$ by Lemma 2.4, which implies $Q$ is hyperbolic by Theorem 2.2(d) and hence overdamped since $B > 0$ and $C \geq 0$. \qed

It follows from (b) in Theorem 2.5 that if we know an upper bound, say, $\theta$, on the largest eigenvalue $\lambda_1$ of a hyperbolic quadratic $Q$, then, with $\lambda = \mu + \theta$, the quadratic $Q_{\theta}$ defined by

\[
Q(\lambda) = Q(\mu + \theta) = \mu^2 A + \mu(B + 2\theta A) + C + \theta B + \theta^2 A
= \mu^2 A_{\theta} + \mu B_{\theta} + C_{\theta}
=: Q_{\theta}(\mu)
\]

is overdamped. Thus any hyperbolic quadratic can be transformed into an overdamped quadratic by an appropriate shifting of the eigenvalues. Hence for the purposes of testing hyperbolicity and overdamping it suffices to consider overdamping. We make this restriction in the next two sections and consider in section 5 how to implement the shifting in practice.

**3. An iteration for testing overdamping.** Suppose we have a Hermitian quadratic $Q(\lambda) = \lambda^2 A + \lambda B + C$, where we assume throughout this section that $A > 0$, $B > 0$, and $C \geq 0$. The challenge is how to test the hyperbolicity (or, equivalently, the overdamping) condition (1.2) or, equivalently, condition (c) in Theorem 2.5.

The primary and secondary solvents $S^{(1)}$ and $S^{(2)}$ of an overdamped quadratic can be found efficiently by applying an iteration based on cyclic reduction [2], [9]. The iteration is

\[
S_0 = B, \quad A_0 = A, \quad B_0 = B, \quad C_0 = C,
S_{k+1} = S_k - A_k B_k^{-1} C_k, \quad A_{k+1} = A_k B_k^{-1} A_k,
B_{k+1} = B_k - A_k B_k^{-1} C_k - C_k B_k^{-1} A_k, \quad C_{k+1} = C_k B_k^{-1} C_k.
\]

The next theorem summarizes properties of the iteration proved in [9, Lemma 6, Theorem 7 and proof].
THEOREM 3.1. Let \( Q(\lambda) = \lambda^2 A + \lambda B + C \) be an \( n \times n \) overdamped quadratic with eigenvalues \( \lambda_1 \geq \cdots \geq \lambda_n > \lambda_{n+1} \geq \cdots \geq \lambda_{2n} \). Consider iteration (3.1) and any matrix norm \( \| \cdot \| \).

(a) The iterates satisfy \( A_k > 0, C_k \geq 0, \) and \( B_k > 0 \) for all \( k \geq 0 \).

(b) \( \| A_k \| \| C_k \| \) converges quadratically to zero with

\[
\limsup_{k \to \infty} \frac{1}{2^k} \sqrt{\| A_k \| \| C_k \|} \leq \frac{\lambda_n}{\lambda_{n+1}} < 1.
\]

(c) \( S_k \) converges quadratically to a nonsingular matrix \( \hat{S} \) with

\[
\limsup_{k \to \infty} \frac{1}{2^k} \sqrt{\| S_k - \hat{S} \|} \leq \frac{\lambda_n}{\lambda_{n+1}} < 1.
\]

(d) The primary and secondary solvents of \( Q(X) \), \( S^{(1)} \) and \( S^{(2)} \), respectively, are given by

\[
S^{(1)} = -\hat{S}^{-1} C, \quad S^{(2)} = -A^{-1} \hat{S}^2.
\]

The next lemma reveals a crucial property of iteration (3.1) for overdamped quadratics. The “only if” part of the result is [9, Lemma 6].

LEMMA 3.2. Let \( \mu > 0 \) and assume \( A_k > 0 \) and \( C_k \geq 0 \). In (3.1), \( B_k > \mu^2 A_k + \mu^{-2} C_k \) if and only if \( A_{k+1} > 0, C_{k+1} \geq 0, \) and \( B_{k+1} > \mu^{2k+1} A_{k+1} + \mu^{-2k+1} C_{k+1} \).

Proof. “\( \Rightarrow \)”: We have

\[
B_{k+1} = B_k - A_k B_k^{-1} C_k - C_k B_k^{-1} A_k \\
= B_k - (\mu^2 A_k + \mu^{-2} C_k) B_k^{-1} (\mu^2 A_k + \mu^{-2} C_k) \\
+ \mu^{2k+1} A_k B_k^{-1} A_k + \mu^{-2k+1} C_k B_k^{-1} C_k \\
> \mu^{2k+1} A_k B_k^{-1} A_k + \mu^{-2k+1} C_k B_k^{-1} C_k,
\]

where we have used the fact that \( X - YX^{-1}Y > Y - YY^{-1}Y = 0 \) when \( X > Y > 0 \). Clearly, \( A_{k+1} > 0 \) and \( C_{k+1} \geq 0 \) since \( B_{k+1} > 0 \).

“\( \Leftarrow \)”: As in the first part we have

\[
B_{k+1} = B_k - F_k B_k^{-1} F_k + F_{k+1},
\]

where \( F_k = \mu^2 A_k + \mu^{-2} C_k \). Now if \( B_{k+1} > \mu^{2k+1} A_{k+1} + \mu^{-2k+1} C_{k+1} = F_{k+1} \) then (3.4) gives \( B_k - F_k B_k^{-1} F_k > 0 \). Note that \( B_k - F_k B_k^{-1} F_k \) is the Schur complement of \( B_k > 0 \) in

\[
T = \begin{bmatrix} B_k & F_k \\ F_k & B_k \end{bmatrix}.
\]

So we have \( T > 0 \), and it follows that \( B_k - F_k > 0 \) (for example, by looking at the \((1,1)\) block of the congruence \( \begin{bmatrix} 1 & -I \end{bmatrix} T \begin{bmatrix} 1 & 0 \end{bmatrix} \)). Therefore \( B_k > F_k = \mu^2 A_k + \mu^{-2} C_k \). \( \square \)

In view of Theorem 2.5(c), Lemma 3.2 implies that \( Q \) is overdamped if and only if any one of the quadratics

\[
Q_k(\lambda) = \lambda^2 A_k + \lambda B_k + C_k
\]
generated during the iteration is overdamped, assuming that \( A_k > 0 \) and \( C_k \geq 0 \) for all \( k \). Note that the latter assumption holds if \( B_k > 0 \) for all \( k \).

**Corollary 3.3.** Let \( Q \) be a Hermitian quadratic with \( A, B > 0 \) and \( C \geq 0 \). For iteration (3.1) and any fixed \( m \geq 0 \), if \( B_k > 0 \) for \( k = 1: m - 1 \) and

\[
B_m > \mu^{2m} A_m + \mu^{-2m} C_m
\]

for some \( \mu > 0 \), then \( B > \mu A + \mu^{-1} C \) and \( Q \) is overdamped.

Intuitively, we can think of the scalars \( \mu^{2m} \) and \( \mu^{-2m} \) in (3.6) as trying to balance \( A_m \) and \( C_m \). This suggests that (3.6) could be replaced by \( B_m > \bar{A}_m + \bar{C}_m \) if the iteration is scaled so that \( \bar{A}_m \) and \( \bar{C}_m \) are balanced. Normwise balancing is included in the following scaled version of (3.1), introduced in [9]; it generates iterates \( \bar{A}_k, B_k \) (unchanged from (3.1)), and \( \bar{C}_k \) according to

\[
\begin{align*}
\alpha_0 &= \sqrt{\|C\|/\|A\|}, \\
\bar{A}_0 &= \alpha_0 A, \quad B_0 = B, \quad \bar{C}_0 = \alpha_0^{-1} C, \\
A_{k+1} &= \bar{A}_k B_k^{-1} \bar{A}_k, \\
B_{k+1} &= B_k - \bar{A}_k B_k^{-1} \bar{C}_k - \bar{C}_k B_k^{-1} \bar{A}_k, \\
C_{k+1} &= \bar{C}_k B_k^{-1} \bar{C}_k, \\
\alpha_{k+1} &= \sqrt{\|C_{k+1}\|/\|A_{k+1}\|}, \\
\bar{A}_{k+1} &= \alpha_{k+1} A_{k+1}, \quad \bar{C}_{k+1} = \alpha_{k+1}^{-1} C_{k+1}.
\end{align*}
\]

Here we have assumed that \( C \neq 0 \) (the overdamping condition holds for the trivial case \( C = 0 \) by (1.2)); thus \( \alpha_k > 0 \) for each \( k \geq 0 \). The scaling procedure ensures that

\[
\|\bar{A}_k\| = \|\bar{C}_k\| \quad \text{and} \quad \|A_k\| \|\bar{C}_k\| = \|A_k\| \|C_k\|.
\]

The next result describes the behavior of the scaled iteration.

**Theorem 3.4.** A Hermitian quadratic \( Q \) with \( A, B > 0 \) and \( 0 \neq C \geq 0 \) is overdamped if and only if in (3.7)

\[
B_k > 0 \quad \text{for all} \quad k, \quad \lim_{k \to \infty} \bar{A}_k = 0, \quad \lim_{k \to \infty} \bar{C}_k = 0, \quad \lim_{k \to \infty} B_k > 0,
\]

and in this case

\[
\begin{align*}
\limsup_{k \to \infty} \frac{1}{k}\sqrt[k]{\|\bar{A}_k\|} &= \limsup_{k \to \infty} \frac{1}{k}\sqrt[k]{\|\bar{C}_k\|} \leq \left( \frac{\lambda_n}{\lambda_{n+1}} \right)^{1/2}, \\
\limsup_{k \to \infty} \frac{1}{k}\sqrt[k]{\|B_k - \bar{B}\|} &\leq \frac{\lambda_n}{\lambda_{n+1}},
\end{align*}
\]

with \( \bar{B} = A(S^{(1)} - S^{(2)}) \).

**Proof.** Assume that the conditions in (3.8) hold. Then \( B_m > \bar{A}_m + \bar{C}_m \) for some \( m \geq 0 \). It is easy to see that the iterates \( \bar{A}_k \) and \( \bar{C}_k \) defined in (3.7) are related to \( A_k \) and \( C_k \) in (3.1) by

\[
\begin{align*}
\bar{A}_k &= \alpha_0^2 \alpha_1^{-k-1} \ldots \alpha_{k-1}^{-2} \alpha_k A_k, \\
\bar{C}_k &= \alpha_0^{-2} \alpha_1^{-2} \ldots \alpha_{k-1}^{-2} \alpha_k^{-1} C_k, \quad k \geq 0.
\end{align*}
\]

So \( B_m > \bar{A}_m + \bar{C}_m \) implies \( B_m > \mu^{2m} A_m + \mu^{-2m} C_m \) with \( \mu = \alpha_0 \alpha_1^2 \alpha_2^{-2} \ldots \alpha_{m-1}^{-2} \alpha_m^{-2} \), which implies \( Q \) is overdamped by Corollary 3.3.
Now assume the QEP is overdamped. Then, from Theorem 3.1(a), $B_k > 0$ for each $k \geq 0$, while, since $\|A_k\| = \|C_k\| = (\|A_k\| \|C_k\|)^{1/2}$, Theorem 3.1(b) implies $\lim A_k = \lim C_k = 0$ and that (3.9) holds. To show the convergence of $B_k$, we note that from (3.1), $B_{k+1} = B_k - (S_k - S_{k+1}) - (S_k - S_{k+1})^*$, which implies

$$B_k = B_0 - (S_0 - S_k) - (S_0 - S_k)^* = -B + S_k + S_k^*.$$  

In view of (3.2), (3.3), and $B_k \geq 0$, (3.10) holds with $\hat{B} = -B + \hat{S} + \hat{S}^* = A(S^{(1)} - S^{(2)}) \geq 0$. Since the sequence $\{\|B_k^{-}\|\}$ is known to be bounded (see the proof of [9, Theorem 7]), we have $\hat{B} > 0$.  

The next result confirms that $\mu$ can be removed from (3.6) for the scaled iteration. It follows readily from Theorem 3.4 and its proof.

**COROLLARY 3.5.** A Hermitian quadratic $Q$ with $A, B > 0$ and $0 \neq C \geq 0$ is overdamped if and only if, for some $m \geq 0$, $B_k > 0$ for $k = 1: m - 1$ in (3.7) and $B_m \geq A_m + C_m$.

The corollary is important for two reasons. First, it provides a basis for an elegant, practical test for overdamping, as definiteness of a matrix is easily tested numerically. Second, in the case of an affirmative test result a $\mu$ with $Q(-\mu) < 0$ can be identified, and such a $\mu$ is very useful when we go on to solve the QEP, as we will show in section 5.

From a numerical point of view it is preferable to work with the original data as much as possible. The following variant of Corollary 3.5 tests the overdamping condition using the original quadratic $Q$ and will be the basis of the algorithm in section 5. It follows readily from Corollary 3.3 and Theorem 3.4 and its proof.

**COROLLARY 3.6.** A Hermitian quadratic $Q$ with $A, B > 0$ and $0 \neq C \geq 0$ is overdamped if and only if, for some $m \geq 0$, $B_k > 0$ for $k = 1: m - 1$ in (3.7) and $Q(-\mu_m) < 0$, where $\mu_m = \alpha_0 \alpha_1^{m-1} \alpha_2^{m-2} \ldots \alpha_{m-1}^{m-1} > 0$ and the $\alpha_k$ are defined in (3.7).

Usually, only a few iterations of the cyclic reduction algorithm (3.7) will be necessary. To illustrate, we consider a quadratic $Q(\lambda)$ of dimension $n = 100$ defined by

$$(3.11) \quad A = I, \quad B = \beta \begin{bmatrix} 20 & -10 \\ -10 & 30 \\ \vdots & \vdots \\ \vdots & \vdots \\ 30 & -10 \\ -10 & 20 \end{bmatrix}, \quad C = \begin{bmatrix} 15 & -5 \\ -5 & 15 \\ \vdots & \vdots \\ \vdots & \vdots \\ -5 & 15 \end{bmatrix},$$

where $\beta > 0$ is a real parameter. This example, which comes from a damped mass-spring system, is used in [13] with $\beta = 1$. We use the 1-norm in (3.7). Tables 3.1 and 3.2 report the number of iterations required to demonstrate that $Q$ is over-

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damped, through verification of the conditions in Corollary 3.6, or that it is not overdamped, through generation of a non-positive definite iterate $B_m$. Note that when $Q$ is “strongly” overdamped and when $Q$ is far from being overdamped, the overdamping condition is shown to hold or not after just a few iterations.

4. Convergence analysis for weakly overdamped quadratics. For the example at the end of section 3 and some $\beta_0 \in (0.5196152422, 0.5196152423)$, the inequality (1.2) holds as a weak inequality with equality attained for some nonzero $x$. We have seen that the overdamping test requires a very small number of iterations when $\beta$ is not close to $\beta_0$. When $\beta \approx \beta_0$, the number of iterations increases but is still under 20 in our experiments. The purpose of this section is to explain this behavior by showing that the convergence of iteration (3.1) is reasonably fast even when the QEP is weakly overdamped in the sense defined as follows.

**Definition 4.1.** $Q(\lambda)$ is weakly hyperbolic if $A$, $B$, and $C$ are Hermitian, $A > 0$, and

$$\gamma = \min_{\|x\| = 1} [(x^* Bx)^2 - 4(x^* Ax)(x^* Cx)] \geq 0.$$  

**Definition 4.2.** $Q(\lambda)$ is weakly overdamped if it is weakly hyperbolic with $B > 0$ and $C \succeq 0$.

The eigenvalues of a weakly hyperbolic $Q$ are real and those of a weakly overdamped $Q$ are real and nonpositive. The following result collects further properties of a weakly overdamped quadratic [25, section 31].

**Theorem 4.3.** Let $Q(\lambda) = \lambda^2 A + \lambda B + C$ be a weakly overdamped $n \times n$ quadratic.

(a) If $\gamma = 0$ in (4.1), then $Q(\lambda)$ has $2n$ real eigenvalues that can be ordered $\lambda_1 \geq \cdots \geq \lambda_n = \lambda_{n+1} \geq \cdots \geq \lambda_{2n}$. The partial multiplicities$^2$ of $\lambda_n$ are at most 2, and the eigenvalues other than $\lambda_n$ are semisimple.

(b) $Q(\lambda_n) \preceq 0$.

(c) The quadratic matrix equation $Q(X) = 0$ in (1.4) has a solvent $S^{(1)}$ with eigenvalues $\lambda_1, \ldots, \lambda_n$ and a solvent $S^{(2)}$ with eigenvalues $\lambda_{n+1}, \ldots, \lambda_{2n}$.

In the overdamped case considered in the previous section, convergence results for the iteration (3.1) are established using matrix identities obtained from the cyclic reduction method. Those identities do not contain enough information about (3.1) to allow a proof of convergence for weakly overdamped quadratics with $\gamma = 0$, for which $\lambda_{n+1} = \lambda_n$. We now study this critical case and thereby obtain a better understanding of the convergence of the iteration for overdamped QEPs with $\lambda_n \approx \lambda_{n+1}$. The next lemma shows that (3.1) remains well defined in the critical case, which is the setting for the rest of this section.

**Lemma 4.4.** For a weakly overdamped quadratic $Q(\lambda) = \lambda^2 A + \lambda B + C$ with $\gamma = 0$ in (4.1), there is a positive real number $\mu$ such that for the iteration (3.1)

$$A_k > 0, \quad C_k > 0, \quad B_k \geq \mu^{2^k} A_k + \mu^{-2^k} C_k$$  

for all $k \geq 0$.

**Proof.** We have $\lambda_n \leq \lambda_1 \leq 0$. If $\lambda_n = 0$ then, from Theorem 4.3, $C = Q(\lambda_n) \preceq 0$. Since $C \succeq 0$ we must have $C = 0$. However, $\gamma > 0$ for the trivial case $C = 0$. Therefore $\lambda_n < 0$ since $\gamma = 0$. It then follows from $Q(\lambda_n) \preceq 0$ that $B \geq \mu A + \mu^{-1} C$.

$^2$The partial multiplicities of an eigenvalue of $Q$ are the sizes of the Jordan blocks in which it appears in a Jordan triple for $Q$ [7].
for \( \mu = -\lambda_n > 0 \). The inequalities in (4.2) are then proved inductively using the technique from the proof of the first part of Lemma 3.2.

Lin and Xu [24] recently showed that Meini’s iterations based on cyclic reduction for the matrix equation \( X + A^*X^{-1}A = Q \) [26] can also be derived from a structure-preserving doubling algorithm. Following their approach we show that the iteration (3.1) is related to a doubling algorithm, and we use this observation to study the convergence of (3.1) for weakly overdamped quadratics. The rate of convergence will be shown to be at least linear with constant 1/2 in the generic case, which is the case where \( \lambda_n = \lambda_{n+1} \) is a multiple eigenvalue with partial multiplicities all equal to 2 (that is, \( \lambda_n \) occurs only in \( 2 \times 2 \) Jordan blocks). This rate and constant are to be expected in view of the results of Guo in [8].

**Lemma 4.5.** Let \( X = \begin{bmatrix} A & 0 \\ H & -I \end{bmatrix} \) and \( Y = \begin{bmatrix} G & I \\ C & 0 \end{bmatrix} \) be block \( 2 \times 2 \) matrices with \( n \times n \) blocks. When \( H + G \) is nonsingular there exist \( 2n \times 2n \) matrices \( \tilde{X} \) and \( \tilde{Y} \) such that \( (a) \tilde{X}Y = \tilde{Y}X \) and \( (b) \tilde{X}X, \tilde{Y}Y \) have the same zero and identity blocks as \( X \) and \( Y \), respectively.

**Proof.** Applying block row permutations and block Gaussian elimination to \( \begin{bmatrix} X \\ Y \end{bmatrix} \) yields \( P \begin{bmatrix} X \\ Y \end{bmatrix} = \begin{bmatrix} U \\ 0 \end{bmatrix} \), where \( U = \begin{bmatrix} G & I \\ G + H & I \end{bmatrix} \) and

\[
P = \begin{bmatrix} P_{11} & P_{12} \\ P_{21} & P_{22} \end{bmatrix} = \begin{bmatrix} 0 & 0 & I & 0 \\ 0 & I & -A(G+H)^{-1} & 0 \\ I & 0 & -A(G+H)^{-1} & 0 \\ 0 & C(G+H)^{-1} & C(G+H)^{-1} & -I \end{bmatrix}.
\]

Since \( [p_{21} \ p_{22}] \begin{bmatrix} X \\ Y \end{bmatrix} = 0 \), the required equality \( \tilde{Y}X = \tilde{X}Y \) is satisfied with \( \tilde{X} := -P_{22} \) and \( \tilde{Y} := P_{21} \). Furthermore,

\[
\tilde{X}X = \begin{bmatrix} A(G+H)^{-1}A & 0 \\ H - C(G+H)^{-1}A & -I \end{bmatrix}, \quad \tilde{Y}Y = \begin{bmatrix} G - A(G+H)^{-1}C & I \\ C(G+H)^{-1}C & 0 \end{bmatrix}.
\]

Lemma 4.5 and its proof suggest the recurrence

\[
X_{k+1} = \tilde{X}_k X_k, \quad Y_{k+1} = \tilde{Y}_k Y_k, \quad k \geq 0,
\]

with

\[
X_k = \begin{bmatrix} A_k & 0 \\ H_k & -I \end{bmatrix}, \quad Y_k = \begin{bmatrix} G_k & I \\ C_k & 0 \end{bmatrix}
\]

and

\[
\tilde{X}_k = \begin{bmatrix} A_k(G_k + H_k)^{-1} & 0 \\ -C_k(G_k + H_k)^{-1} & I \end{bmatrix}, \quad \tilde{Y}_k = \begin{bmatrix} I & -A_k(G_k + H_k)^{-1} \\ 0 & C_k(G_k + H_k)^{-1} \end{bmatrix},
\]

which leads to

\[
A_{k+1} = A_k(G_k + H_k)^{-1}A_k, \quad G_{k+1} = G_k - A_k(G_k + H_k)^{-1}C_k, \quad H_{k+1} = H_k - C_k(G_k + H_k)^{-1}A_k, \quad C_{k+1} = C_k(G_k + H_k)^{-1}C_k.
\]

With

\[
A_0 = A, \quad C_0 = C, \quad G_0 = 0, \quad H_0 = B,
\]
the iteration (3.1) is recovered from (4.5) by letting $B_k = G_k + H_k$ and $S_k = H^*_k$. By Lemma 4.4, $B_k > 0$ for all $k \geq 0$. Therefore with the starting matrices (4.6), iteration (4.5) is well defined. Note that $X_k$ in (4.4) is nonsingular for all $k \geq 0$ and, from (4.3) and the property that $\tilde{X}_kY_k = \tilde{Y}_kX_k$,

\begin{equation}
X_{k+1}^{-1}Y_{k+1} = (\tilde{X}_kX_k)^{-1}\tilde{Y}_kY_k = X_k^{-1}\tilde{X}_k^{-1}\tilde{Y}_kY_k = X_k^{-1}Y_kX_k^{-1}Y_k = (X_k^{-1}Y_k)^2.
\end{equation}

It follows from (4.7) that for all $k \geq 0$,

\begin{equation}
X_k^{-1}Y_k = (X_0^{-1}Y_0)^{2^k}.
\end{equation}

The identity (4.8) is what we need to prove the convergence of (4.5) with (4.6) and hence the convergence of (3.1).

The next result describes the convergence behavior in the generic case.

**Theorem 4.6.** Let $Q(\lambda)$ be weakly overdamped with eigenvalues $\lambda_1 \geq \cdots \geq \lambda_n = \lambda_{n+1} \geq \cdots \geq \lambda_{2n}$, and assume that the partial multiplicities of $\lambda_n$ are all equal to 2. Let $S^{(1)}$ and $S^{(2)}$ be the primary and secondary solvents of $Q(X) = 0$, respectively, and assume that $\lambda_n$ is a semisimple eigenvalue of $S^{(1)}$ and $S^{(2)}$. Then the iterates $G_k$, $H_k$, $A_k$, and $C_k$ defined by (4.5) and (4.6) satisfy

\[
\limsup_{k \to \infty} \sqrt[k]{\|G_k - AS^{(1)}\|} \leq \frac{1}{2}, \quad \limsup_{k \to \infty} \sqrt[k]{\|H_k + AS^{(2)}\|} \leq \frac{1}{2}, \quad \limsup_{k \to \infty} \sqrt[k]{\|A_k\|C_k\|} \leq \frac{1}{4}.
\]

**Proof.** We start by making the change of variables (or scaling) $\lambda = \mu \theta$, where $\theta = |\lambda_n| > 0$ (see the proof of Lemma 4.4) so that $\mu_k = \mu_{n+1} = -1$, and we define $\tilde{Q}(\mu) = \mu^2\tilde{A} + \mu\tilde{B} + \tilde{C}$ with $(\tilde{A}, \tilde{B}, \tilde{C}) = (\theta A, B, \theta^{-1}C)$. For this triple denote the iterates of (4.5) by $\tilde{A}_k$, $\tilde{G}_k$, $\tilde{H}_k$, and $\tilde{C}_k$. It is easy to see that for all $k \geq 0$, $\tilde{G}_k = G_k$, $\tilde{H}_k = H_k$, $\tilde{A}_k = \theta^k A_k$, and $\tilde{C}_k = \theta^{-2k} C_k$ so that $\|A_k\|C_k\| = \|\tilde{A}_k\|\|\tilde{C}_k\|$. The primary and secondary solvents of $\tilde{A}S^{(2)} + \tilde{B}S + \tilde{C} = 0$ are $\tilde{S}^{(1)} = \theta^{-1}S^{(1)}$ and $\tilde{S}^{(2)} = \theta^{-1}S^{(2)}$, respectively. Note that $\tilde{A}S^{(i)} = AS^{(i)}$, $i = 1, 2$. To avoid notational clutter, we omit the hats on matrices in the rest of the proof.

We now consider the iterations for the block $2 \times 2$ matrices $X_k$ and $Y_k$ in (4.4). With $A_0 = A$, $C_0 = C$, $G_0 = 0$, and $H_0 = B$, the pencil

\begin{equation}
\mu X_0 + Y_0 = \mu \begin{bmatrix} A & 0 \\ B & -I_n \end{bmatrix} + \begin{bmatrix} 0 & I_n \\ C & 0 \end{bmatrix}
\end{equation}

is a linearization of $Q(\mu)$ [7]. Hence $-X_0^{-1}Y_0$ and $Q(\mu)$ have the same eigenvalues, with the same partial multiplicities. Suppose there are $r$ $2 \times 2$ Jordan blocks associated with eigenvalues equal to $\mu_n = -1$, where $r \geq 1$ by assumption. Rearranging the Jordan canonical form of $-X_0^{-1}Y_0$ appropriately yields

\begin{equation}
V^{-1}(X_0^{-1}Y_0)V = \begin{bmatrix} D_2 \oplus I_r & 0 \oplus I_r \\ 0 & D_1 \oplus I_r \end{bmatrix} =: D_V,
\end{equation}

\begin{equation}
W^{-1}(X_0^{-1}Y_0)W = \begin{bmatrix} D_2 \oplus I_r & 0 \\ 0 \oplus I_r & D_1 \oplus I_r \end{bmatrix} =: D_W,
\end{equation}

where $V$ and $W$ are nonsingular, $D_1$ and $D_2$ are $(n-r) \times (n-r)$ diagonal matrices containing the (semisimple) eigenvalues less than 1 and greater than 1 in modulus,
respectively, and $M \oplus N$ denotes $\begin{bmatrix} M & 0 \\ 0 & N \end{bmatrix}$. Now partition $V$ and $W$ as block $2 \times 2$ matrices with $n \times n$ blocks

\[ V = \begin{bmatrix} V_1 & V_3 \\ V_2 & V_4 \end{bmatrix}, \quad W = \begin{bmatrix} W_1 & W_3 \\ W_2 & W_4 \end{bmatrix}, \]

and note from (4.10)–(4.11) that

\[ (4.12) \quad X_0^{-1}Y_0 \begin{bmatrix} V_1 \\ V_2 \end{bmatrix} = \begin{bmatrix} V_1 \\ V_2 \end{bmatrix}(D_2 \oplus I_r), \quad X_0^{-1}Y_0 \begin{bmatrix} W_3 \\ W_4 \end{bmatrix} = \begin{bmatrix} W_3 \\ W_4 \end{bmatrix}(D_1 \oplus I_r). \]

By Theorem 4.3 and our assumption on $S^{(1)}$ and $S^{(2)}$ there exist nonsingular $U_1$ and $U_2$ such that

\[ (4.13) \quad -S^{(1)} = U_1(D_1 \oplus I_r)U_1^{-1}, \quad -S^{(2)} = U_2(D_2 \oplus I_r)U_2^{-1}. \]

Since $S^{(i)}$, $i = 1, 2$, is a solution of $Q(X) = 0$, from (4.9) we obtain

\[ (4.14) \quad X_0^{-1}Y_0 \begin{bmatrix} I_n \\ -AS^{(i)} \end{bmatrix} = \begin{bmatrix} I_n \\ -AS^{(i)} \end{bmatrix}(-S^{(i)}), \quad i = 1, 2. \]

On comparing with the invariant subspaces in (4.12) and using (4.13) we deduce that

\[ \begin{bmatrix} V_1 \\ V_2 \end{bmatrix} = \begin{bmatrix} U_2 \\ -AS^{(2)}U_2 \end{bmatrix}Z_1, \quad \begin{bmatrix} W_3 \\ W_4 \end{bmatrix} = \begin{bmatrix} U_1 \\ -AS^{(1)}U_1 \end{bmatrix}Z_2, \]

with $Z_1$ and $Z_2$ nonsingular, where we have also used the fact that there are exactly $r$ eigenvectors of $X_0^{-1}Y_0$ corresponding to the eigenvalue 1. Hence $V_1$ and $W_3$ are nonsingular and

\[ (4.15) \quad -AS^{(2)} = V_2V_1^{-1}, \quad -AS^{(1)} = W_4W_3^{-1}. \]

By (4.8)–(4.11) we have $V^{-1}(X_k^{-1}Y_k)V = D_2^{2k}$ and $W^{-1}(X_k^{-1}Y_k)W = D_1^{2k}$, so that

\[ (4.16) \quad G_kV_1 + V_2 = A_kV_1(D_2^{2k} \oplus I_r), \]
\[ (4.17) \quad G_kV_3 + V_4 = A_kV_1(0 \oplus 2^{2k}I_r) + A_kV_3(D_1^{2k} \oplus I_r), \]
\[ (4.18) \quad C_kV_1 = (H_kV_1 - V_2)(D_2^{2k} \oplus I_r), \]
\[ (4.19) \quad C_kV_3 = (H_kV_1 - V_2)(0 \oplus 2^{2k}I_r) + (H_kV_3 - V_4)(D_1^{2k} \oplus I_r) \]

and

\[ (4.20) \quad G_kW_1 + W_2 = A_kW_1(D_2^{2k} \oplus I_r) + A_kW_3(0 \oplus 2^{2k}I_r), \]
\[ (4.21) \quad G_kW_3 + W_4 = A_kW_3(D_1^{2k} \oplus I_r), \]
\[ (4.22) \quad C_kW_1 = (H_kW_1 - W_2)(D_2^{2k} \oplus I_r) + (H_kW_3 - W_4)(0 \oplus 2^{2k}I_r), \]
\[ (4.23) \quad C_kW_3 = (H_kW_3 - W_4)(D_1^{2k} \oplus I_r). \]
By (4.22) and (4.23) we have

\[(4.24) \quad C_k(W_3 - W_1(0 \oplus 2^{-k} I_r)) = (H_k W_3 - W_4) (D_1^{2k} \oplus 0) - (H_k W_1 - W_2) (0 \oplus 2^{-k} I_r).\]

By (4.18) we have

\[(4.25) \quad H_k = V_2 V_1^{-1} + C_k V_1 (D_2^{-2k} \oplus I_r) V_1^{-1}.\]

Inserting (4.25) in (4.24) we obtain

\[C_k \left( W_3 - W_1(0 \oplus 2^{-k} I_r) - V_1 (D_2^{-2k} \oplus I_r) V_1^{-1} (W_3(D_1^{2k} \oplus 0) - W_1(0 \oplus 2^{-k} I_r)) \right) = (V_2 V_1^{-1} W_3 - W_4)(D_1^{2k} \oplus 0) - (V_2 V_1^{-1} W_1 - W_2)(0 \oplus 2^{-k} I_r),\]

from which it follows, since $D_1$ and $D_2$ are diagonal with diagonal elements of magnitude less than 1 and greater than 1, respectively, that

\[(4.26) \quad C_k = O(2^{-k});\]

the latter notation means that $\|C_k\| = O(2^{-k})$. It then follows from (4.25) and (4.14) that

\[(4.27) \quad H_k + AS^{(2)} = H_k - V_2 V_1^{-1} = O(2^{-k}).\]

By (4.20) and (4.21),

\[(4.28) \quad G_k W_3 + W_4 - (G_k W_1 + W_2)(0 \oplus 2^{-k} I_r) = A_k (W_3(D_1^{2k} \oplus 0) - W_1(0 \oplus 2^{-k} I_r)).\]

By (4.16),

\[(4.29) \quad A_k = (G_k V_1 + V_2)(D_2^{-2k} \oplus I_r) V_1^{-1}.\]

Inserting (4.29) in (4.28) we obtain

\[G_k W_3 + W_4 - (G_k W_1 + W_2)(0 \oplus 2^{-k} I_r) = (G_k V_1 + V_2) M_k,\]

with $M_k = O(2^{-k})$. Thus

\[-G_k (W_3 - W_1(0 \oplus 2^{-k} I_r) - V_1 M_k) = W_4 - W_2(0 \oplus 2^{-k} I_r) - V_2 M_k.\]

It follows from (4.14) that

\[(4.30) \quad G_k - AS^{(1)} = G_k + W_4 W_3^{-1} = O(2^{-k}).\]

Postmultiplying (4.16) by $D_2^{-2k} \oplus 0$ gives

\[(4.31) \quad (G_k V_1 + V_2)(D_2^{-2k} \oplus 0) = A_k V_1 (I_r \oplus 0),\]

while postmultiplying (4.17) by $0 \oplus 2^{-k} I_r$ gives

\[(4.32) \quad (G_k V_3 + V_4)(0 \oplus 2^{-k} I_r) = A_k V_3(0 \oplus I_r) + A_k V_3(0 \oplus 2^{-k} I_r).\]

Adding (4.31) and (4.32) we get

\[A_k (V_1 + V_3(0 \oplus 2^{-k} I_r)) = (G_k V_1 + V_2)(D_2^{-2k} \oplus 0) + (G_k V_3 + V_4)(0 \oplus 2^{-k} I_r).\]
It follows that
\[ A_k = O(2^{-k}), \]
since \( \{G_k\} \) has been shown to be bounded. Equations (4.26), (4.27), (4.30), and (4.33) yield the required convergence results.

For \( S_k \) and \( B_k \) in iteration (3.1) we obtain the following convergence result.

**Corollary 4.7.** Under the conditions of Theorem 4.6, the iterates \( S_k \) and \( B_k \) in (3.1) satisfy
\[
\limsup_{k \to \infty} \sqrt[k]{\|S_k - \hat{S}\|} \leq \frac{1}{2}, \quad \limsup_{k \to \infty} \sqrt[k]{\|B_k - \hat{B}\|} \leq \frac{1}{2},
\]
where \( \hat{S} = -S^{(2)} A \) is nonsingular and \( \hat{B} = A(S^{(1)} - S^{(2)}) \geq 0 \) is singular.

**Proof.** The convergence results follow from Theorem 4.6 by noting \( B_k = H_k + G_k \) and \( S_k = H_k \). By (4.27) and (4.30), \( \hat{B} = A(S^{(1)} - S^{(2)}) \). We have \( \hat{B} \geq 0 \) since \( B_k > 0 \) for each \( k \), by Lemma 4.4. We now show that \( \hat{B} \) is singular. Using (4.9) it is easy to check that
\[
(-X_0^{-1}Y_0) \left[ \begin{array}{cc} I & I \\ -AS^{(1)} & -AS^{(2)} \end{array} \right] = \left[ \begin{array}{cc} I & I \\ -AS^{(1)} & -AS^{(2)} \end{array} \right] (S^{(1)} \oplus S^{(2)}),
\]
and \( S^{(1)} \oplus S^{(2)} \) is diagonalizable. Now \( -X_0^{-1}Y_0 \) is not diagonalizable, by assumption, since it has at least one eigenvalue of partial multiplicity 2. Thus (4.34) can only hold if \( \left[ \begin{array}{cc} I & I \\ -AS^{(1)} & -AS^{(2)} \end{array} \right] \) is singular. Thus the Schur complement \( \hat{B} = A(S^{(1)} - S^{(2)}) \) is singular.

In the generic case for a weakly overdamped \( Q \) with \( \gamma = 0 \), in which all of the partial multiplicities of \( \lambda_n \) are 2, \( Q \) is in some sense irreducible or coupled. The next example shows that this condition is necessary for the conclusions in Theorem 4.6 and Corollary 4.7 (and at the same time answers an open question from [9, section 4]).

Consider
\[
Q(\lambda) = \lambda^2 A + \lambda B + C = \lambda^2 \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} + \lambda \begin{bmatrix} 3 & 0 \\ 0 & 1 \end{bmatrix} + \begin{bmatrix} 2 & 0 \\ 0 & 0 \end{bmatrix}.
\]

It is easy to see that \( \gamma = 0 \), so \( Q(\lambda) \) is weakly overdamped with eigenvalues \( \{0, -1, -1, -2\} \) with \( \lambda_2 = \lambda_3 = -1 \) semisimple. In (3.1) and (4.5), (4.6),
\[
\lim_{k \to \infty} A_k = \begin{bmatrix} 0 & 0 \\ 0 & 1 \end{bmatrix}, \quad \lim_{k \to \infty} B_k = I_2, \quad \lim_{k \to \infty} C_k = \begin{bmatrix} 1 & 0 \\ 0 & 0 \end{bmatrix}, \quad \lim_{k \to \infty} G_k = \begin{bmatrix} -1 & 0 \\ 0 & 0 \end{bmatrix}, \quad \lim_{k \to \infty} H_k = \begin{bmatrix} 2 & 0 \\ 0 & 1 \end{bmatrix}.
\]

Neither \( A_k \) nor \( C_k \) converges to zero. We also note that the convergence is quadratic for \( B_k, G_k, \) and \( H_k \). Moreover, \( B_k \) converges to a nonsingular matrix. This does not come as a surprise, since \( Q(\lambda) \) can be decomposed into the direct sum of two scalar quadratics
\[
Q_1(\lambda) = \lambda^2 + 3\lambda + 2, \quad Q_2(\lambda) = \lambda^2 + \lambda.
\]

It is readily seen that \( Q_1 \) is overdamped with eigenvalues \( -1, -2 \) and that \( Q_2 \) is overdamped with eigenvalues \( 0, -1 \). Thus the convergence of \( B_k \) to a positive definite matrix is guaranteed by Theorem 3.4 applied to each component of the direct sum.
5. Algorithm for the detection and numerical solution. Let \( Q(\lambda) = \lambda^2 A + \lambda B + C \) be Hermitian with \( A > 0 \). We develop in this section an efficient algorithm that checks if \( Q \) is hyperbolic and, if it is, computes some or all of the eigenvalues and associated eigenvectors, exploiting the symmetry and hyperbolicity and thereby preserving the spectral properties.

Our algorithm consists of three steps:

1. **Preprocessing.** This step forms \( Q_\theta(\lambda) = Q(\lambda + \theta) = \lambda^2 A_\theta + \lambda B_\theta + C_\theta \) with \( \theta \) such that \( B_\theta > 0 \) and \( C_\theta \geq 0 \) or concludes that \( Q \) is not hyperbolic and terminates the algorithm.

2. **Overdamping test.** This step checks the overdamping condition for \( Q_\theta \). If \( Q_\theta \) is overdamped, a \( \mu \in \mathbb{R} \) such that \( Q_\theta(\mu) = Q(\mu + \theta) < 0 \) is also computed; otherwise the algorithm terminates.

3. **Solution.** The quadratic \( Q_\theta \) is converted into a definite pencil \( \lambda X + Y \in \mathbb{C}^{2n \times 2n} \) with \( X > 0 \) or \( Y > 0 \). The eigenvalues and eigenvectors of \( Q(\lambda) \) are then obtained from the eigendecomposition of a \( 2n \times 2n \) Hermitian matrix obtained by transforming \( \lambda X + Y \) and exploiting the definiteness of \( X \) or \( Y \) and the block structure of \( X \) and \( Y \).

We now detail each of these three steps and compare the cost and stability of our solution process with that of three alternative ways of solving the QEP: The QZ algorithm applied to a linearization of \( Q(\lambda) \), the \( J \)-orthogonal Jacobi algorithm [28] also applied to a linearization of \( Q(\lambda) \), and the method of computing the eigenpairs of the primary and secondary solvents obtained via the cyclic reduction method [9].

At different stages our algorithm needs to test the (semi)definiteness of a matrix. This is done by attempting a Cholesky factorization, with complete pivoting in the case of semidefiniteness: Completion of the factorization means the matrix is (semi)definite. This is a numerically stable test, as shown in [10].

5.1. **Preprocessing step.** The preprocessing step aims to eliminate, by simple tests, quadratics that are not hyperbolic and to produce, if possible, a shifted quadratic \( Q_\theta(\lambda) = Q(\lambda + \theta) \) (with \( \theta = 0 \) is possible) for which the necessary condition

\[
B > 0, \quad C \geq 0
\]

for overdamping is satisfied.

If \( B \) is singular then, by (1.2), \( Q \) cannot be hyperbolic. Assume now that \( B \) is nonsingular but not positive definite or \( C \) is not positive semidefinite. Since \( A > 0 \), for \( \theta > 0 \) large enough the matrices

\[
B_\theta = B + 2\theta A, \quad C_\theta = C + \theta B + \theta^2 A
\]

defining the shifted quadratic \( Q_\theta(\lambda) = Q(\lambda + \theta) \) with \( A_\theta = A \) (see (2.2)) satisfy (5.1).

To avoid numerical instability in the formation of \( B_\theta \) and \( C_\theta \) (due to the possibly large variation in \( \|A\|, \|B\|, \) and \( \|C\| \)) we would ideally like to choose \( \theta \) close to

\[
\theta_{opt} = \inf\{\theta \in \mathbb{R} : B + 2\theta A > 0, \ C + \theta B + \theta^2 A \geq 0\}.
\]

Rather than solving this optimization problem we choose \( \theta \) to be an upper bound on the modulus of \( \lambda_1 \), the right-most eigenvalue of \( Q \). With such a shift, all of the eigenvalues of \( Q_\theta \) lie in the left half plane. When \( Q \) is hyperbolic, \( Q_\theta \) is also hyperbolic with real and nonpositive eigenvalues. Thus \( B_\theta > 0 \) and \( C_\theta \geq 0 \) by Theorem 2.5. Therefore if \( B_\theta \neq 0 \) or \( C_\theta \geq 0 \) we can conclude that \( Q \) is not hyperbolic. If \( B_\theta > 0 \) and \( C_\theta \geq 0 \) we proceed to step 2.
Table 5.1
Operation count for the preprocessing step. Matrices are assumed real and of dimension $n$.

<table>
<thead>
<tr>
<th>Operations</th>
<th>Cost (flops)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Cholesky factorization of $B$ and $C$ to check definiteness.</td>
<td>$2n^3/3$ or less</td>
</tr>
<tr>
<td>Computation of $\theta$ when $B$ and/or $C$ not positive definite:</td>
<td>$n^3/3$</td>
</tr>
<tr>
<td>Cholesky factorization of $A$.</td>
<td>$4n^2$</td>
</tr>
<tr>
<td>$|A^{-1}|$ (1-norm estimation [11, section 15.3], typically 4 solves).</td>
<td>$6n^2$</td>
</tr>
<tr>
<td>Form $B_0 = B + 2\theta A$, $C_0 = C + \theta B + \theta^2 A$.</td>
<td>$2n^3/3$ or less</td>
</tr>
<tr>
<td>Cholesky factorizations of $B_0$ and $C_0$.</td>
<td></td>
</tr>
</tbody>
</table>

To construct the shift $\theta$ we use the following strategy: let

$$a = \|A\|, \quad b = \|B\|, \quad c = \|C\|,$$

where $\|\cdot\|$ is any consistent matrix norm. Then, from [18, Lemmas 3.1 and 4.1], for every eigenvalue $\lambda$ of $Q$ we have

(5.2) $|\lambda| \leq \frac{1}{2}\|A^{-1}\|\left(b + \sqrt{b^2 + 4c/\|A^{-1}\|}\right) =: \sigma_1,$

(5.3) $|\lambda| \leq (1 + \|A^{-1}\|) \max(c^{1/2}, b) =: \sigma_2.$

We take the 1-norm and set $\sigma = \min(\sigma_1, \sigma_2)$. Since $\sigma$ must greatly overestimate $|\lambda_1|$ when $|\lambda_n| \gg |\lambda_1|$, we carry on one step further and form the shifted quadratic $Q_{-\sigma/2}(\lambda) = Q(\lambda - \sigma/2)$ for which (5.2)–(5.3) give two new bounds $\tau_1$ and $\tau_2$ (and $A$ is unchanged, so $\|A^{-1}\|$ can be reused). We then take $\theta = \min(\sigma, \tau - \frac{1}{2}\sigma)$, where $\tau = \min(\tau_1, \tau_2)$.

As shown by Theorem 3.4, the speed of convergence of iteration (3.7) for overdamped $Q$ depends on the ratio $\lambda_n/\lambda_{n+1}$. An unnecessarily large shift of the spectrum to the left can make this ratio very close to 1, potentially causing slow convergence of the iteration. However, we showed in section 4 that for the generic case of weakly overdamped $Q$ with $\lambda_n = \lambda_{n+1}$ the convergence is at least linear with constant 1/2, so convergence of the iteration cannot be unduly delayed by a conservative choice of shift.

Table 5.1 details the computations and their cost. (Costs of all the operations used here are summarized in [12, Appendix C].) Preprocessing requires at most $\frac{1}{3}n^3$ flops.

5.2. Overdamping test. The following algorithm is based on Corollary 3.6. It runs the scaled iteration (3.7) until either a non-positive definite $B_k$ or a negative definite $Q(\mu_k)$ is detected, signaling that $Q$ is not overdamped or is overdamped, respectively. The algorithm terminates on one of these conditions or because of possible nonconvergence of the iteration for a non-overdamped $Q$. It is intended to be applied to $Q_0$ from the preprocessing step.

Algorithm 5.1 (overdamping test). This algorithm tests whether a quadratic $Q(\lambda) = \lambda^2A + AB + C$ with $A, B > 0$ and $0 \neq C \geq 0$ is overdamped and, if it is, computes $\mu < 0$ such that $Q(\mu) < 0$. Input parameters are the maximum number of iterations $k_{\text{max}}$ and a convergence tolerance $\epsilon > 0$.

1. Set $A_0 = A$, $B_0 = B$, $C_0 = C$.
2. Set $a_0 = \|C_0\|/\|A_0\|$, $\mu_0 = -a_0^{1/2}$, $k = 0$.
3. If $Q(\mu_0) < 0$, $Q(\lambda)$ is (hyperbolic and hence) overdamped, $\mu = \mu_0$, quit, end

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4 while $k < k_{\text{max}}$
5 \[ B_{k+1} = B_k - A_k B_k^{-1} C_k - C_k B_k^{-1} A_k \]
6 if \[ \|B_{k+1} - B_k\|_1/\|B_{k+1}\|_1 \leq \epsilon, \text{ goto line 15, end} \]
7 if $B_{k+1} \neq 0$, $Q$ is not overdamped, quit, end
8 $A_{k+1} = \alpha_k A_k B_k^{-1} A_k$
9 $C_{k+1} = \alpha_k^{-1} C_k B_k^{-1} C_k$
10 $\alpha_{k+1} = \|C_{k+1}\|_1/\|A_{k+1}\|_1$
11 $\mu_{k+1} = \mu_k \alpha_{k+1}$
12 if $Q(\mu_{k+1}) < 0$, $Q$ is overdamped, $\mu = \mu_{k+1}$, quit, end
13 $k = k + 1$
14 end
15 $Q$ is not overdamped. % See the discussion below.

Note that the crucial definiteness test on line 12 of Algorithm 5.1 is carried out on $Q$ and not on $Q_k$ in (3.5). Hence a positive test can be interpreted irrespective of rounding errors in the iteration: The only errors are in forming $Q_{k-1}$ and in computing its Cholesky factor. For a non-overdamped $Q$, it is possible that $B_k > 0$ for all $k$ (see the example at the end of section 4). However, if convergence of the $B_k$ is detected on line 6 then $Q$ is declared not overdamped because by this point an overdamped $Q$ would have been detected, while if $k_{\text{max}}$ is large enough (say, $k_{\text{max}} = 20$) and this iteration limit is reached then $Q$ can reasonably be declared not overdamped in view of the fast (quadratic) convergence of (3.7) for an overdamped $Q$.

The implementation details of Algorithm 5.1 and the cost per iteration are described in Table 5.2. The total cost for $m$ iterations is $\frac{2}{3} n^3$ flops for $m = 0$ and roughly $\frac{2}{3} m n^3$ flops for $m \geq 1$.

Guo and Lancaster’s test for overdamping is based on iteration (3.1), scaled as in (3.7). For the computation of $\tilde{S}$, $\frac{19}{3} \ell n^3$ flops are required, where $\ell$ is the number of iterations for convergence of (3.1). An extra $5n^3$ flops is needed to form the two solvents $S^{(1)}$ and $S^{(2)}$ (which are nonsymmetric in general) via (3.3). Then the smallest eigenvalue $\lambda_n$ of $S^{(1)}$ and the largest eigenvalue $\lambda_{n+1}$ of $S^{(2)}$ need to be computed and the definiteness of $Q((\lambda_n + \lambda_{n+1})/2)$ tested. The total cost is $(\frac{2}{3} \ell + \frac{20}{3}) n^3$ flops plus the cost of finding $\lambda_n$ and $\lambda_{n+1}$. Since $m \leq \ell$, Algorithm 5.1 is clearly the more efficient, possibly significantly so.

We mention two alternative ways to test hyperbolicity. Both are based on the fact that a Hermitian $Q$ with $A > 0$ is hyperbolic if and only if a certain $2n \times 2n$ pair

<p>| Table 5.2 |
|-------------------------|-------------------------|
| <strong>Operation count per complete iteration of Algorithm 5.1. Matrices are assumed real and of dimension $n$.</strong> |</p>
<table>
<thead>
<tr>
<th>Operations</th>
<th>Cost (flops)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Cholesky factorization of $B_k = L_k L_k^T$ available from previous step.</td>
<td>$n^3$</td>
</tr>
<tr>
<td>Form $V_k = L_k^{-1} A_k$.</td>
<td>$n^3$</td>
</tr>
<tr>
<td>Form $W_k = L_k^{-1} C_k$.</td>
<td>$n^3$</td>
</tr>
<tr>
<td>Compute $A_k B_k^{-1} C_k = V_k^T W_k$.</td>
<td>$2n^3$</td>
</tr>
<tr>
<td>Cholesky of $B_{k+1}$.</td>
<td>$n^3/3$</td>
</tr>
<tr>
<td>Compute $A_k B_k^{-1} A_k = V_k^T V_k$.</td>
<td>$n^3$</td>
</tr>
<tr>
<td>Compute $C_k B_k^{-1} C_k = W_k^T W_k$.</td>
<td>$n^3$</td>
</tr>
<tr>
<td>Cholesky of $-Q(\mu_{k+1})$.</td>
<td>$n^3/3$</td>
</tr>
<tr>
<td>Total</td>
<td>$20n^3/3$</td>
</tr>
</tbody>
</table>
The equivalence between $Q$ to the fact that the ill-conditioned congruence transformation is implicitly present in no warning. The poor performance of PDFIND when working with the pair is due whether the congruence is used or not, and when it gives the wrong answer there is 0.51966, 0.51969, 0.51970. The conclusion is that PDFIND is numerically unreliable gruence incorrectly detects nonoverdamping for 0.51965, 0.51967, 0.51968 and fails for 0.5177–0.5181, 0.5188, 0.5189, 0.5194. Next, we take Algorithm 5.1. However, as $\alpha$ decreases, PDFIND becomes more unreliable, due to the increasing ill conditioning of the congruence transformation with decreasing $\alpha$. To be more specific, we take $\alpha = 10^{-7}$. First, consider $\beta = 0.5157:0.0001:0.5197$. For $\beta = 0.5197$ our algorithm detects overdamping in 3 iterations, and for other values our algorithm detects nonoverdamping in at most 8 iterations. PDFIND, using the congruence, incorrectly detects nonoverdamping for $\beta = 0.5197$ and fails for $\beta = 0.5157, 0.5177–0.5181, 0.5188, 0.5189, 0.5194$. Next, we take $\beta = 0.51965:0.00001:0.51971$. Our algorithm detects overdamping in at most 5 iterations. PDFIND without the congruence incorrectly detects nonoverdamping for 0.51965, 0.51967, 0.51968 and fails for 0.51966, 0.51969, 0.51970. The conclusion is that PDFIND is numerically unreliable whether the congruence is used or not, and when it gives the wrong answer there is no warning. The poor performance of PDFIND when working with the pair is due to the fact that the ill-conditioned congruence transformation is implicitly present in the equivalence between $Q$ being hyperbolic and $(A, B)$ being definite.

For our algorithm, instability could potentially arise if $B_k$ is ill-conditioned. However, we know from [2, p. 40, line 10] that $B_k$ is well-conditioned if $B_0 = B$ is well-conditioned (which is verifiable right from the beginning) and if $\lambda_n/\lambda_{n+1}$ is not too close to 1. When $\lambda_n/\lambda_{n+1}$ is extremely close to 1, $B_k$ is known to be ill-conditioned for large $k$. However, $B_k$ appears in our algorithm only in terms like $A_k B_k^{-1} C_k$, and $A_k$ and $C_k$ converge to 0, so the ill-conditioning of $B_k$ has only a limited effect on our algorithm. Indeed, instability has not been observed in any of our tests.
5.3. Solving hyperbolic QEPs via definite linearizations. Recall that the scalar $\mu$ computed by Algorithm 5.1 applied to $Q_\nu$ is such that $Q(\mu + \theta) = Q_\theta(\mu) < 0$. Hence with $\omega = \mu + \theta$ we have

$$
\check{Q}(t) = Q(t + \omega) = t^2 A + t(B + 2\omega A) + C + \omega B + \omega^2 A
= t^2 \check{A} + t \check{B} + \check{C},
$$

with $\check{C} = Q(\omega) < 0$ and $\check{A} = A > 0$. The pencils

$$
L_1(\lambda) = \lambda \begin{bmatrix} \check{A} & 0 \\ 0 & -\check{C} \end{bmatrix} + \begin{bmatrix} \check{B} \\ \check{C} \end{bmatrix}, \quad L_2(\lambda) = \lambda \begin{bmatrix} 0 & \check{A} \\ \check{A} & 0 \end{bmatrix} + \begin{bmatrix} -\check{A} & 0 \\ 0 & \check{C} \end{bmatrix}
$$

are both Hermitian definite linearizations of $\check{Q}$ with a positive definite leading coefficient of $L_1$ and a negative definite trailing coefficient of $L_2$. They share the same eigenvalues as $Q$, and the eigenvectors of $\check{Q}$ are easy to recover from those of $L_1$ or $L_2$. The sensitivity and stability of these linearizations have recently been studied in [14], [15], [17]. It is shown therein that the scaling of Fan, Lin and Van Dooren [6] should be applied to $\check{Q}$ before linearizing. The choice between $L_1$ and $L_2$ should be guided by the fact that, in terms of conditioning and backward error, they favor large and small eigenvalues, respectively. However, if $\check{C}$ or $\check{A}$ is well-conditioned and $||B||/(||A|| ||C||)^{1/2}$ is not much bigger than 1 then $L_1$ or $L_2$, respectively, can safely be used to stably obtain all of the eigenpairs. For more details on conditioning and backward error for $L_1$ and $L_2$, see [14], [15], [17].

Using Cholesky factorizations $A = L_A L_A^T$ and $-\check{C} = L_C L_C^T$, the definite generalized eigenvalue problem $L_1(\lambda)z = 0$ or $L_2(\lambda)z = 0$ is transformed to a Hermitian (or real symmetric) standard eigenvalue problem [5]. For example, $L_1(\lambda)$ reduces to

$$
\lambda I + \begin{bmatrix} L_A^{-1} B L_A^{-T} & -L_A^{-1} L_C \\ -L_C L_A^{-T} & 0 \end{bmatrix}.
$$

As Table 5.3 explains, this phase requires about $13n^3$ flops, giving a grand total of $(\frac{20}{3}m + 13)n^3$ flops.

Guo and Lancaster’s solution algorithm has a total cost of $(\frac{12}{3}m + 25)n^3$ flops, assuming the eigenvalues of $S^{(1)}$ and $S^{(2)}$ (which are the eigenvalues of $Q$) are computed by the QR algorithm. In practice this is significantly more than the cost of our algorithm given that $m \leq \ell$ is usually small.

<table>
<thead>
<tr>
<th>Operations</th>
<th>Cost (flops)</th>
</tr>
</thead>
<tbody>
<tr>
<td>Cholesky factorizations of $A = L_A L_A^T$ and $-C = L_C L_C^T$ already available.</td>
<td>$n^3/3$</td>
</tr>
<tr>
<td>Form $R = -(L_A^{-1} L_C)^T$.</td>
<td>$3n^3/2$</td>
</tr>
<tr>
<td>Form $G = L_A^{-1} B L_A^{-T}$.</td>
<td>$&lt; 4(2n)^3/3$</td>
</tr>
<tr>
<td>Tridiagonalization of $[ -R^T 0 ]$.</td>
<td>$O(n^2)$</td>
</tr>
<tr>
<td>Eigenvalues via (e.g.) QR iteration.</td>
<td>Total $\approx 13n^3$</td>
</tr>
</tbody>
</table>
The most common way of solving the QEP is to apply the QZ algorithm or a Krylov method to a linearization $L$ of $Q$. The QZ algorithm applied to the $2n \times 2n$ $L$ costs $240n^3$ flops for the computation of the eigenvalues.

Our algorithm has two important advantages over that of Guo and Lancaster and QZ applied to a linearization, besides its more favorable operation count. First, it works entirely with symmetric matrices, which reduces the storage requirement. Second, it guarantees to produce real eigenvalues in floating point arithmetic; the other two approaches cannot do so because they invoke the QZ algorithm and the nonsymmetric QR algorithm.

6. Numerical experiment. We describe an experiment that illustrates the behavior of our algorithm for testing overdamping. More extensive testing of this algorithm, and of the preprocessing and solving procedures described in section 5, will be presented in a future publication. Our experiments were performed in MATLAB 7.4 (R2007a), for which the unit roundoff is $u = 2^{-53} \approx 1.1 \times 10^{-16}$. We took $k_{\text{max}} = 30$ and $\epsilon = u$ in Algorithm 5.1.

We first describe a useful technique for generating symmetric quadratic matrix polynomials with prescribed eigenvalues and eigenvectors and positive definite coefficient matrices.

Let $(\lambda_k, v_k), k=1:2n$, be a set of given real eigenpairs such that, with

$$
\Lambda = \text{diag}(\lambda_1, \ldots, \lambda_{2n}) =: \Lambda_1 \oplus \Lambda_2, \quad \Lambda_1, \Lambda_2 \in \mathbb{R}^{n \times n},
$$

$$
V := [v_1, \ldots, v_{2n}] =: [V_1 \quad V_2], \quad V_1, V_2 \in \mathbb{R}^{n \times n},
$$

$V_1$ and $V_2$ are nonsingular, and

$$(6.1) \quad V_1 V_1^T = V_2 V_2^T, \quad V_1 \Lambda_1 V_1^T - V_2 \Lambda_2 V_2^T =: \Gamma \text{ is nonsingular.}$$

Then the symmetric quadratic polynomial defined by the matrices

$$(6.2a) \quad A = \Gamma^{-1}, \quad B = -A(V_1 \Lambda_1^2 V_1^T - V_2 \Lambda_2^2 V_2^T)A,$$

$$(6.2b) \quad C = -A(V_1 \Lambda_1^3 V_1^T - V_2 \Lambda_2^3 V_2^T)A + B \Gamma B$$

has eigenpairs $(\lambda_k, v_k), k=1:2n$ (see [23] for example). We now show how to generate a potentially overdamped quadratic.

**Lemma 6.1.** Assume that $0 > \lambda_1 \geq \cdots \geq \lambda_n > \lambda_{n+1} \geq \cdots \geq \lambda_{2n}$. Then $\Gamma$ is nonsingular and the matrices generated by (6.2) satisfy $A > 0$, $B > 0$, and $C > 0$.

**Proof.** It follows from Weyl’s theorem [20, p. 181] that $\Gamma \succ 0$ and hence that $A > 0$. All matrices $V_2$ that satisfy the first constraint in (6.1) can be written as $V_1 U$ for some orthogonal $U$. Hence we can write

$$
B = -AV_1(\Lambda_1^2 - U \Lambda_2^2 U^T)V_1^T A = -AV_1(H_1^2 - H_2^2)V_1^T A,
$$

where $H_1 = \Lambda_1$ and $H_2 = U \Lambda_2 U^T$, and again Weyl’s theorem guarantees that $B \succ 0$.

It is known that $(V, \Lambda, P\Lambda^T)$, where $P = \text{diag}(I_n, -I_n)$, forms a self-adjoint triple for $Q(\lambda)$ [7, section 10.2]. Since $Q$ has no zero eigenvalues, $C$ is nonsingular and a formula for its inverse is easily obtained from the resolvent form of $Q(\lambda)$: For $\lambda \neq \lambda_i$,

$$
Q(\lambda)^{-1} = V(\lambda I_{2n} - \Lambda)^{-1} P \Lambda^T.
$$
Table 6.1
Minimum, average, and maximum number of iterations performed by Algorithm 5.1 and percentage of overdamped problems, for each n and matrix type.

<table>
<thead>
<tr>
<th>n</th>
<th>type 1</th>
<th>type 2</th>
<th>type 3</th>
</tr>
</thead>
<tbody>
<tr>
<td>5</td>
<td>0.0, 2.4, 6.0</td>
<td>0.0, 0.8, 3.0</td>
<td>0.0, 2.4, 5.0</td>
</tr>
<tr>
<td>10</td>
<td>0.0, 3.6, 10.0</td>
<td>0.0, 0.5, 3.0</td>
<td>2.0, 2.7, 4.0</td>
</tr>
<tr>
<td>50</td>
<td>0.0, 4.2, 11.0</td>
<td>0.0, 2.1, 4.0</td>
<td>2.0, 2.1, 3.0</td>
</tr>
<tr>
<td>100</td>
<td>3.0, 6.2, 10.0</td>
<td>0.0, 2.6, 4.0</td>
<td>2.0, 2.0, 2.0</td>
</tr>
<tr>
<td>250</td>
<td>2.0, 6.0, 11.0</td>
<td>2.0, 3.0, 4.0</td>
<td>2.0, 2.0, 2.0</td>
</tr>
<tr>
<td>500</td>
<td>3.0, 7.5, 11.0</td>
<td>2.0, 3.0, 4.0</td>
<td>2.0, 2.0, 2.0</td>
</tr>
</tbody>
</table>

Setting $\lambda = 0$ in the above expression gives
\[
C^{-1} = -V\Lambda^{-1}PV^T = -V_1(H_1^{-1} - H_2^{-1})V_1^T,
\]
and once again Weyl’s theorem guarantees that $C^{-1}$, and therefore also $C$, is positive definite.

We use the following eigenvalue distributions:

**type 1:** $\lambda_k$, $k = 1:2n$, is uniformly distributed in $[-100, -1]$.

**type 2:** $\lambda_k$ is uniformly distributed in $[-100, -6]$ for $k = n + 1:2n$ and $[-5, -1]$ for $k = 1:n$.

**type 3:** $\lambda_k$ is uniformly distributed in $[-100, 20]$.

$B$ and $C$ are then shifted as in (2.2) with $\theta = 1.1\lambda_1$ to ensure that the eigenvalues are all negative.

We took $V_1 = U_1$ and $V_2 = V_1U_2$, where $U_1$ and $U_2$ are random orthogonal matrices from the Haar distribution [11, section 28.3]. For types 1 and 2, $A$, $B$, and $C$ are all positive definite by construction; for type 3 nothing can be said about the definiteness of $A$, $B$, and $C$. Table 6.1 shows the minimum, average, and maximum number of iterations for Algorithm 5.1 over 20 quadratics for each of several values of $n$, along with the percentage of $Q$ found to be overdamped for each $n$ and matrix type. In all cases where $Q$ was deemed overdamped, the computed $\mu$ was verified to lie in $(\lambda_{n+1}, \lambda_n)$.

We make several observations:

- For all three eigenvalue distributions, Algorithm 5.1 is quick to terminate, especially for types 2 and 3, with only very occasional need for more than 10 iterations. The gap between $\lambda_n$ and $\lambda_{n+1}$ is larger for type 2 than type 1, which explains the greater number of iterations for type 1.

- With $V_1$ orthogonal the coefficients matrices $A$, $B$, and $C$ are well-conditioned, with 2-norm condition numbers of order $10^2$. If instead we take $V_1$ as a random matrix with 2-norm condition number $10^4$ (computed in MATLAB as `gallery('randsvd', n, 1e4, ...)`), the condition numbers of $A$, $B$, and $C$ are of order $10^8$ and the number of iterations of the algorithm increases, though only slightly: The maximum number of iterations over all tests is 13, and the largest average over all $n$ rises to 7.8, 3.1, and 3.2 for types 1, 2, and 3, respectively.

- After detecting overdamping an average of 6–9 more iterations are needed for convergence of the block cyclic iteration. Recall that the algorithm of Guo and Lancaster [9] needs to iterate to convergence in order to show overdamping.

**Acknowledgments.** This work was started while the first author visited MIMS in the School of Mathematics at the University of Manchester in 2005; he thanks the School for its hospitality. We thank Qiang Ye for helpful comments concerning the algorithm of Crawford and Moon.
REFERENCES