A fast algorithm for the computation of an upper bound on the \( \mu \)-norm

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Abstract

A fast algorithm for the computation of the optimally frequency-dependent scaled \( \mathcal{H}_\infty \)-norm of a finite-dimensional LTI system is presented. It is well known that this quantity is an upper bound to the “\( \mu \)-norm”; furthermore, it was recently shown to play a special role in the context of slowly time-varying uncertainty. Numerical experimentation suggests that the algorithm generally converges quadratically. © 2000 Elsevier Science Ltd. All rights reserved.

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1. Introduction

In the context of robust control analysis and synthesis a quantity of great interest is the structured singular-value norm, or \( \mu \)-norm, of the system. Consider a feedback connection of a continuous-time system with real coefficients as in Fig. 1. Let \( P(s) = C(sI - A)^{-1}B \) be an \( m \times m \) stable transfer function matrix and let \( \Delta(s) \) be a structured perturbation constrained to lie in the set

\[ \mathcal{R}(\Lambda) \triangleq \{ \Lambda \in \mathcal{R}_S : \Delta(s_0) \in \Lambda \quad \forall s_0 \in \mathbb{C}_+ \} \]

where \( \mathcal{R}_S \) is the set of all real-rational, proper, stable, \( m \times m \) transfer matrices and the uncertainty set

\[ \Delta \triangleq \{ \text{diag}[\Delta_1, \ldots, \Delta_F] : \Delta_j \in \mathbb{C}_+^{m_j \times m_j} \} \]

is defined for integers \( m_1, \ldots, m_F \). Let \( \sigma_1(\cdot) \) denote the maximum singular value of its matrix argument and define the complex structured singular value for a constant matrix \( M \in \mathbb{C}_+^{m \times m} \) as (see Zhou, Doyle & Glover, 1995 for a complete discussion of the structured singular value)

\[ \mu_\Lambda(M) \triangleq \frac{1}{\min[\sigma_1(\Delta): \Delta \in \Lambda, \det(I - M\Delta) = 0]} \]

unless \( \det(I - M\Delta) \neq 0 \) for all \( \Delta \in \Lambda \), in which case \( \mu_\Lambda(M) \neq 0 \). Finally, define the “\( \mu \)-norm”\(^3\) as

\[ ||P||_\Lambda \triangleq \sup \mu_\Lambda(P(j\omega)). \]

It has been shown that the computation of \( \mu_\Lambda(M) \) is NP-hard (see Toker & Özbay, 1995), in consequence no efficient algorithms are likely to exist for its computation.

\(^3\) Actually, it is not a norm.
In practice, a standard upper bound is used in its place. This upper bound is computed as follows. Define
\[ \mathcal{D} = \{ \text{diag}[d_1 I_{m_1}, \ldots, d_{F-1} I_{m_{F-1}}, I_{m_F}] : 0 < d_j \in \mathbb{R} \}. \]
Then, for a constant complex matrix \( M \), an upper bound for \( \mu_\Delta(M) \) is
\[ \mu_\Delta(M) \triangleq \inf_{D \in \mathcal{D}} \sigma_1(DMD^{-1}). \]
Substituting this into the expression for \( \|P\|_\Delta \), an upper bound on the \( \mu \)-norm of the system \( P(s) \) is obtained as the optimally frequency-dependent scaled \( \mathcal{H}_\infty \)-norm:
\[ \|P\|_\Delta \triangleq \sup_{\omega \in \mathcal{F}} \mu_\Delta(P(j\omega)). \]
It is well-known that if \( P(s) \) is stable, then \( \|P\|_\Delta < 1 \) is necessary and sufficient (thus \( \|P\|_\Delta < 1 \) is sufficient) for uniform robust stability of the \( P - \Delta \) loop (Fig. 1) for any linear time-invariant structured \( \Delta(s) \) of \( L_2 \)-gain no greater than one (see, e.g., Corollary 3 in Tits & Balakrishnan, 1998). Further, it has recently been shown that \( \|P\|_\Delta < 1 \) is necessary and sufficient for uniform robust stability of the \( P - \Delta \) loop for any linear, arbitrarily slowly time-varying structured \( \Delta \) of \( L_2 \)-gain no greater than one (Poolla & Tikku, 1995).

Algorithms for the efficient computation of \( \mu_\Delta(P(j\omega)) \), for given \( \omega \), have long been available. In fact, given \( \omega \),
\[ \mu_\Delta(P(j\omega))^2 = \inf_{\alpha, D \in \mathcal{D}} \{ \mathbf{z} : P(j\omega)D P(j\omega)^* - \alpha D < 0 \}, \]
which is a linear matrix inequality (LMI) problem. Efficient algorithms exist for obtaining global solutions to such problems, e.g., Boyd, El Ghaoui, Feron and Balakrishnan (1994). Note that minimizers (or approximate minimizers) \( D \) for \( \sigma_1(DMD^{-1}) \) are related to the minimizers (or approximate minimizers) \( \tilde{D} \) for (1) by \( D = \tilde{D}^{1/2} \). \( \|P\|_\Delta \) is usually computed via a “frequency sweep”, i.e. choose a set of frequencies \( \{ \omega_1, \ldots, \omega_N \} \) and use the approximation
\[ \|P\|_\Delta \approx \max \{ \mu_\Delta(P(j\omega_1)), \ldots, \mu_\Delta(P(j\omega_N)) \}. \]
The drawbacks of this approach are obvious. First, a large number of \( \mu_\Delta(\cdot) \) computations are required.

Second, an upper bound on \( \|P\|_\Delta \) is not necessarily obtained. Finally, the result can be arbitrarily bad, i.e. it is difficult to bound the error.

Another approach that has been used to compute an upper bound to the \( \mu \)-norm is based on the Main Loop Theorem (see, e.g., Packard & Doyle, 1993) and the extension to \( \mu \) of the Maximum Modulus Theorem (see, e.g., top of p. 1201 in Tits & Fan, 1995). In the discrete-time case, given \( P(z) = C(zI - A)^{-1}B + D \), with \( A \) stable, and a positive scalar \( \gamma \), it can be seen that the \( \mu \)-norm of \( P \) is less than \( \gamma \) if and only if \( \mu_\Delta(M(\gamma)) < 1 \), where
\[ M(\gamma) = \begin{bmatrix} A & B \\ \frac{\gamma C}{\delta} & \frac{\gamma D}{} \end{bmatrix} \]
and \( \Delta_p \) is an “augmented” uncertainty structure given by
\[ \Delta_p \triangleq \{ \text{diag}[\delta I_n, \Delta] : \delta \in \mathbb{C}, \Delta \in \Delta \}. \]
The original idea is due to Doyle and Packard (1987). The continuous-time case can be reduced to the discrete-time case by means of a bilinear transformation (see Section 10.2 in Zhou et al., 1995 for details). Repeated evaluation of the upper bound \( \mu_\Delta(\cdot) \) according to a bisection search over \( \gamma \) yields an upper bound on the \( \mu \)-norm of \( P \). On the down side, note that this upper bound is generally less tight than that obtained by gridding, because the augmented uncertainty structure involves an additional block which, to make things worse, is of the “repeated” type. The bisection search can be done away with, as shown by Ferreres and Fromion (1997), by invoking the “skewed \( \mu \)” proposed and studied by Fan and Tits (1992).\(^5\) Computation then entails the solution of a single LMI constrained quasi-convex generalized eigenvalue minimization problem (GEVP; see Boyd et al., 1994).\(^6\) The approach proposed in this paper improves on the scheme of Ferreres and Fromion (1997) in that it computes the same (tighter) upper bound as the gridding approach. As far as computational cost is concerned, the tradeoff is that of a sequence of LMIs (\( \mu \) upper bound computations with respect to the original, non-augmented uncertainty structure) versus a single GEVP with a larger number of variables and constraints. Due to the “repeated” block, the complexity of the latter is strongly affected by the dimension of the state space.

Recently, an efficient algorithm has been proposed (Boyd & Balakrishnan, 1990; Bruinsma & Steinbuch, 1991).\(^5\) In Ferreres and Fromion (1997) the continuous-time case is dealt with directly, by limiting the frequency range to a compact interval and using a “repeated real scalar” block in the augmentation. This idea is due to Sideris (1992).\(^6\) To make the optimization problem of Fan and Tits (1992) into a standard GEVP, a symmetric \( n \times n \) matrix of slack variables must be introduced; see p. 8–33 in Gahinet et al. (1995).
This algorithm makes use of the well-known fact that a given scalar \( \xi > 0 \) is a singular value of \( (C(j\omega) - A)^{-1}B \) if and only if \( j\omega \) is an eigenvalue of the related Hamiltonian matrix

\[
H(\xi, A, B, C) = \begin{bmatrix}
A & \xi^{-1}BB^* \\
-\xi^{-1}C^*C & A^*
\end{bmatrix}
\]

Given any \( \xi \in (0, \|P\|_\infty) \), a set of frequency intervals may thus be computed where maximizers for (2) are known to lie. At step \( k \), \( \xi \) is selected as \( \sigma_1(P(j\omega_{k-1})) \) and \( \omega_k \) is chosen as the mid-point of the largest among the intervals thus determined. A quadratic rate of convergence ensues.

In this paper, using the idea just outlined as a stepping stone, an algorithm is constructed for the fast computation of \( \|P\|_\infty \). (A similar algorithm can be used for the computation of the real eigenvalues; see Sreedhar, Van Dooren & Tits, 1996.)

2. Key ideas

Let

\[
\tilde{\mu}(\omega) = \inf_{D \in \mathcal{D}} \sigma_1(\mathcal{L}(j\omega)D^{-1}).
\]

The goal of the algorithm is to maximize \( \tilde{\mu}(\omega) \) over \( \omega \in \mathbb{R} \).

Define

\[
\tilde{\mu}^* = \sup_{\omega \in \mathbb{R}} \tilde{\mu}(\omega)
\]

and

\[
\omega^* = \arg \max_{\omega \in \mathbb{R}} \tilde{\mu}(\omega)
\]

assuming for now that such a unique maximizer exists.

Finally, for a fixed \( D \in \mathcal{D} \), define the curve

\[
v_D(\omega) = \sigma_1(\mathcal{L}(j\omega)D^{-1}).
\]

While reading the following, it may be helpful to refer to Fig. 2. The main idea of the algorithm is as follows (more details are given in Section 3). At iteration \( k \), suppose \( \tilde{\omega}_k \) is the best-known lower bound to \( \tilde{\mu}^* \) thus far, and let \( \omega_k \) be the current trial frequency. Suppose further that \( \omega^* \) is known to lie in a certain open set \( \Omega_k \).

Compute

\[
D_{k+1} = \arg \min_{D \in \mathcal{D}} \sigma_1(\mathcal{L}(j\omega_k)D^{-1})
\]

(assume for simplicity that the minimum is achieved), and let

\[
\tilde{\omega}_{k+1} = \tilde{\mu}(\omega_k) = v_{D_{k+1}}(\omega_k).
\]

If \( \tilde{\omega}_{k+1} > \tilde{\omega}_k \), take \( \tilde{\omega}_{k+1} = \tilde{\omega}_k \) as the new estimate of \( \tilde{\mu}^* \) (as is the case in Fig. 2), otherwise keep the old estimate (i.e. \( \tilde{\omega}_{k+1} = \tilde{\omega}_k \)). Note that

\[
\tilde{\mu}(\omega) \leq v_{D_{k+1}}(\omega) \quad \forall \omega \in \mathbb{R},
\]

with equality at \( \omega = \omega_k \), since by definition \( \tilde{\mu}(\omega) \) is the lower envelope of the family of curves \( \{v_D(\omega) : D \in \mathcal{D}\} \).

Now note that \( v_{D_{k+1}}(\omega) \) is the maximum singular value curve for the transfer function \( D_{k+1}P(s)D_{k+1}^{-1} \) at \( s = j\omega \). Thus, as suggested in Section 1, the ideas developed for the computation of the \( \mathcal{H}_\infty \)-norm may be applied to locate the open set of frequencies, call it \( \Omega_{k+1} \), in which \( v_{D_{k+1}}(\omega) > \tilde{\omega}_{k+1} \) (in Fig. 2, \( \Omega_{k+1} = (1.5, \infty) \)). In view of (3) and since \( \tilde{\omega}_{k+1} \) is a lower bound to \( \tilde{\mu}^* \), it is clear that (i) if \( \Omega_{k+1} = 0 \), then \( \tilde{\mu}^* = \tilde{\omega}_{k+1} \), and (ii) if \( \Omega_{k+1} \neq 0 \), then \( \omega^* \in \Omega_{k+1} \).

In the latter case it follows that \( \omega^* \in \Omega_{k+1} \). \( \Omega_{k+1} \), \( \Omega_{k+1} \) (in Fig. 2, it is assumed that \( \Omega_{k+1} = \Omega_{k+1} \), \( \Omega_{k+1} = \Omega_{k+1} \)). Now, the next frequency \( \omega_{k+1} \) must be selected. Several possibilities come to mind. For example, \( \omega_{k+1} \) could simply be chosen as the mid-point of the largest interval contained in \( \Omega_{k+1} \).

Since, for all \( k \), \( \Omega_{k+1} \not\supset \omega_k \), the size of the largest interval in \( \Omega_{k+1} \) would go to zero as \( k \to \infty \) and convergence of \( \tilde{\omega}_k \) to \( \tilde{\mu}^* \) would ensue (see Section 4 below). This, however, results in a relatively slow rate of convergence. More sophisticated alternatives are considered in the next section.

A rough outline of the algorithm described above is as follows.

Algorithm

Data: \( P(s) = C(sI - A)^{-1}B \), \( A \), \( \varepsilon > 0 \).

Initialization: \( k = 0 \), pick \( \omega_0 \geq 0 \), \( \tilde{\omega}_0 = 0 \), \( \Omega_0 = (0, \infty) \).

Footnotes:

1. Below, \( \Omega_k \) is restricted to be a subset of \((0, \infty)\) for all \( k \). Thus, in the case that \( \omega^* = 0 \), clearly \( \omega^* \) will not be in \( \Omega_k \). It is readily checked that this will not affect any of the results.
Step 1: Compute $D_{k+1}, \xi_{k+1}$.
Step 2: Compute $\Omega_{k+1}$.
Step 3: Compute $\omega_k + 1$.
Step 4: $k \leftarrow k + 1$. If an appropriate stopping criterion is satisfied, stop. Otherwise, go back to Step 1.

3. Details of the algorithm

The computation of $\Omega_k$ is now considered in greater detail. The pure imaginary eigenvalues of the Hamiltonian matrix

$$H_k = H(\xi_{k+1}, A, B D_{k+1}^{-1}, D_{k+1} C),$$

tell us the frequencies at which one of the singular values of $D_{k+1} P(\omega) D_{k+1}^{-1}$ takes the value $\xi_{k+1}$. That is, if $j\omega$ is an imaginary eigenvalue of $H_k$, then

$$\sigma_l(D_{k+1} P(j\omega) D_{k+1}^{-1}) = \xi_{k+1}$$

for some $r \in \{1, \ldots, m\}$ (see Fig. 3). Of course, the frequencies of interest are those for which $r = 1$. Suppose the Schur decomposition of the Hamiltonian matrix $H_k$ has been performed and the $\omega_l$'s have been identified. Let $\omega_1, \ldots, \omega_q$ be the nonnegative $\omega_l$'s, sorted in increasing order, and let $\omega_l = j\omega_l, \ell = 1, \ldots, q$. In Boyd and Balakrishnan (1990), the authors suggest determining those frequencies corresponding to the maximum singular value by eliminating the frequencies for which

$$\xi_{k+1}^2 I - D_{k+1}^{-1} P\omega D_{k+1}^2 P(\omega) D_{k+1}^{-1}$$

is not positive semi-definite. This is a rather costly computation, though. An alternative scheme is as follows. The pure imaginary simple eigenvalues of the Hamiltonian may be differentiated with respect to the scalar parameter $\xi$ via the expression

$$\frac{\partial \lambda_l}{\partial \xi}(\xi_{k+1}) = \frac{u_l^T (iH/e\xi) \xi_{k+1} A, B D_{k+1}^{-1} D_{k+1} C) v_l}{u_l^T v_l},$$

\(\ell = 1, \ldots, q,\)

(see, e.g., Theorem IV.2.3 in Stewart and Sun (1990)) where $v_l$ and $u_l^\#$ are respectively right and left eigenvectors of $H_k$ corresponding to the imaginary eigenvalue $\omega_l$; these few vectors can be computed from the Schur decomposition at little extra cost, e.g., using “inverse iteration”. Since $\omega_l = j\omega_l$,

$$\frac{\partial \xi}{\partial \omega}(\omega_l) = \left( - \frac{\partial \lambda_l}{\partial \xi}(\xi_{k+1}) \right)^{-1},$$

which is the slope of the corresponding singular-value curve for the matrix $D_{k+1} P(j\omega) D_{k+1}^{-1}$ at the frequency $\omega = \omega_l$ and is a real number for every simple eigenvalue on the imaginary axis. This is all the information needed to compute $\Omega_{k+1}$.

Indeed, define

$$s_l = \frac{\partial \xi}{\partial \omega}(\omega_l), \quad \ell = 1, \ldots, q.$$

For the example in Fig. 3, the sign of $s_l$ is marked at each intersection point. As $P(s)$ is strictly proper, $v_{Dk+1}(\omega) \to 0$ as $\omega \to \infty$. Thus, as long as $\xi_{k+1} > 0$, $\omega_l$ must correspond to the largest singular value, i.e. $v_{Dk+1}(\omega_l) = \xi_{k+1}$, and $\Omega_{k+1}$ will be bounded. Let $c_l, \ell = 1, \ldots, q, c_l \in \{1, \ldots, m\}$, denote the number of the singular-value curve corresponding to $\omega_l$. That is,

$$\sigma_{c_l}(D_{k+1} P(j\omega_l) D_{k+1}^{-1}) = \xi_{k+1}.$$ 

Thus, $c_q = 1$. The rest of the $c_l$'s may be assigned by sweeping to the left from $\omega_{q-1}$ and using the following rules:

- If $s_{r+1} < 0$ and $s_r < 0$ then $c_r = c_{r+1} + 1.$
- If $s_{r+1} < 0$ and $s_r > 0$ then $c_r = c_{r+1}.$
- If $s_{r+1} > 0$ and $s_r < 0$ then $c_r = c_{r+1}.$
- If $s_{r+1} > 0$ and $s_r > 0$ then $c_r = c_{r+1} - 1.$

Now define

$$\{\omega_1, \ldots, \omega_q\} \triangleq \{\omega_l; c_l = 1\},$$

ordered so that $\omega_1 < \cdots < \omega_q$. Clearly, $r \leq q$. In the obvious manner define $s_l, \ell = 1, \ldots, q$ as the slopes of the maximum singular value curve computed at $\omega_l$. If $s_1 < 0$,

8Certainly true if $\|P(\omega_0)\| > 0$, which will hold in all but a few pathological cases.
9The algorithm as stated does not handle the case when $s_l = 0$ or the eigenvalue curve is non-differentiable. This happens when multiple eigenvalues are found on the imaginary axis. A simple way to get around the difficulty in such cases is to increase the current level $\xi_k$ slightly.
then
\[ \Omega_{k+1} = (\hat{0}_1, \hat{0}_2, \hat{0}_3) \cup \cdots \cup (\hat{0}_{r-1}, \hat{0}_r). \]

If \( \hat{s}_1 > 0 \), then
\[ \Omega_{k+1} = (\hat{0}_1, \hat{0}_2) \cup \cdots \cup (\hat{0}_{r-1}, \hat{0}_r). \]

To help clarify the ideas, consider the example in Fig. 3 where \( q = 6 \). Letting \( c_6 = 1 \), we begin the left sweep. As \( s_6 < 0 \) and \( s_3 < 0 \), \( c_6 = c_6 + 1 = 2 \). Then, since \( s_3 < 0 \) and \( s_4 > 0 \), \( c_5 = c_5 = 2 \). Next, since \( s_4 > 0 \) and \( s_3 > 0 \), \( c_4 = c_4 - 1 = 1 \). Proceeding in this manner, we find \( \{c_1, \ldots, c_6\} = \{1, 1, 1, 2, 2, 1\} \). Thus, \( r = 4 \) and \( \{\hat{0}_1, \ldots, \hat{0}_4\} = \{\hat{0}_1, \hat{0}_2, \hat{0}_3, \hat{0}_6\} \). Finally, since \( \hat{s}_1 = s_1 > 0 \), \( \Omega_{k+1} = (\hat{0}_1, \hat{0}_2) \cup \cdots \cup (\hat{0}_4, \hat{0}_6) \).

As mentioned in Section 2, one way to select the next trial point \( \omega_{k+1} \) is to pick the mid-point of the largest interval in \( \Omega_{k+1} \). A better approach, yielding faster convergence, is to use some form of interpolation scheme in the final iterations. In particular, suppose \( (\omega_{k_1}, \hat{c}_{k_1} + 1) \) and \( (\omega_{k_2}, \hat{c}_{k_2} + 1) \) are the two most recent trial points for which \( \hat{\omega}(\hat{0}_k) = \hat{c}_{k_1} + 1, \hat{c}_{k_2} + 1 \), \( i = 1, 2 \). Then, for \( k = k_1, i = 1, 2 \), \( \hat{c}_{k_1} + 1 = \hat{c}_{k_2} + 1 \). Thus \( \omega_{k_1} \) is one of the \( \hat{0}_k \)s for the \( (k + 1) \)st iteration, say \( \omega_{k_1} = \hat{0}_{i_1} \), and the corresponding \( c_{i_1} \) is 1, i.e.,
\[ \hat{\mu}(\omega_{k_1}) = \sigma_1(D_{k_1} + 1)P(j\omega_{k_1})D_{k_1}^{-1} = \hat{c}_{k_1} + 1. \]

Since the \( \hat{\mu}(\cdot) \) curve lies entirely below the \( \sigma_1(D_{k_1} + 1)P(j\omega_{k_1})D_{k_1}^{-1} \) curve, for all such \( k \), whenever both curves are differentiable at \( \omega_{k_1} \),
\[ \frac{\partial \hat{\mu}}{\partial \omega}(\omega_{k_1}) = \frac{\partial}{\partial \omega} \sigma_1(D_{k_1} + 1)P(j\omega_{k_1})D_{k_1}^{-1} = \frac{\partial \hat{c}}{\partial \omega}(\hat{0}_{i_1}) = s_{i_1}. \]

Using the available values of \( \hat{\mu} \) and of its derivative at \( \omega_{k_1} \) and \( \omega_{k_2} \), a cubic or rational (specifically, quadratic over linear) function (Barzilai & Ben-Tal, 1982) may be passed through the two points and \( \omega_{k_1} \) and \( \omega_{k_2} \) may be taken as the maximum of the interpolating function, subject to the constraint \( \omega_{k_1} \in \Omega_{k+1} \) (closure of \( \Omega_{k+1} \)).

Finally, a stopping criterion which guarantees that the algorithm terminates with an estimate that is within prescribed \( \epsilon > 0 \) of \( \|P\|_\mu \) is as follows. At iteration \( k \), given \( \hat{c}_k, D_k \), and \( \Omega_k \), define the set
\[ S_k \triangleq \{\omega: v_D(\omega) > \hat{c}_k + \epsilon\}. \]

Like \( \Omega_{k+1} \), \( S_k \) may be computed using an appropriate Hamiltonian eigen-decomposition. Since \( \omega^* \in \Omega_k \), if \( \Omega_k \cap S_k = \emptyset \), then \( \omega^* \notin \Omega_k \), hence
\[ \hat{c}_k < \|P\|_\mu \leq \hat{c}_k + \epsilon. \]

4. Convergence
Theorems 1 and 2 below are proved for the case when the mid-point rule is used throughout to select the next trial point \( \omega_{k+1} \).

Assumption 1. The choice of \( \omega_0 \) is such that
\[ \hat{\mu}(\omega_0) > 0. \]

Assumption 1 is necessary in order for the algorithm to be well-defined. Typically, \( \omega_0 = 0 \) will satisfy this assumption; if not, a random search will provide a suitable \( \omega_0 \) as long as \( P(s) \) is not identically 0. Under this assumption, \( \Omega_1 \) is bounded and, in view of the fact that \( \omega_0 \notin \Omega_{k+1} \), a simple induction argument shows that the length of the largest interval in \( \Omega_k \) tends to zero as \( k \to \infty \).

Theorem 1. The sequence \( \{\hat{c}_k\}_{k \in \mathbb{N}} \) converges to \( \hat{\mu}^* \) as \( k \to \infty \).

Proof. By contradiction. Suppose \( \hat{c}_k \to \hat{\mu}^* \). Since the sequence \( \{\hat{c}_k\}_{k \in \mathbb{N}} \) is monotone non-decreasing, and since \( \hat{c}_k \leq \hat{\mu}^* \) for all \( k \), \( \hat{c}_k \to \hat{\mu}^* \) for some \( \hat{\mu}^* < \hat{\mu}^* \). Thus, \( E \triangleq \{\omega: \hat{\mu}(\omega) > \hat{\mu}^* \} \) is not empty and, clearly, \( E \subset \Omega_k \) for all \( k \). Since \( \hat{\mu}(\cdot) \) is continuous (see Bercovici, Foias & Tannenbaum, 1990), \( E \) contains a non-trivial interval \( S \). This contradicts the fact that the size of the largest interval in \( \Omega_k \) goes to 0.

For this preliminary convergence analysis, consider now the problem of maximizing \( \hat{\mu}(\omega) \), where
\[ \hat{\mu}(\omega) \triangleq \min_{\mu} \{\sigma_1(D_P(j\omega)D_P^{-1}): D \in \mathcal{D}, bI \leq D^2 \leq bI\} \quad (4) \]
with \( b \geq b > 0 \) arbitrary prescribed numbers. Note that \( \hat{\mu}(\omega) \) is an upper bound to \( \hat{\mu}(\omega) \), that it can be efficiently computed, and that it can be made arbitrarily close to \( \hat{\mu}(\omega) \), uniformly over \( \omega \in \mathbb{R} \), by selecting \( b > 0 \) small enough and \( b \) large enough (Lee & Tits, 1993). Let
\[ \hat{\mu}^* \triangleq \sup_{\omega \in \mathbb{R}} \hat{\mu}(\omega). \]

It is not difficult to show that Theorem 1 still holds when \( \hat{\mu}(\omega) \) is replaced with \( \hat{\mu}(\omega) \) and \( \hat{\mu}^* \) is replaced with \( \hat{\mu}^* \).

Theorem 2. Assume that \( \omega^* \) is the unique global maximizer of \( \hat{\mu}(\omega) \). Then \( \omega_k \to \omega^* \) as \( k \to \infty \).

Proof. By contradiction. Let \( \mathcal{X} \) be an infinite index set and suppose that \( \omega_k \to \hat{\omega} \) on \( \mathcal{X} \), with \( \hat{\omega} \neq \omega^* \) (since \( \Omega_1 \) is bounded, such \( \mathcal{X} \) exists for some \( \hat{\omega} \) in the closure of \( \Omega_1 \)). By continuity of \( \hat{\mu}(\cdot) \) and uniqueness of the global maximizer \( \omega^* \), there exists \( \delta > 0 \) such that, for \( k \) large enough, \( k \in \mathcal{X} \),
\[ \hat{\mu}(\omega_k) < \hat{\mu}^* - \delta. \]

Let \( D_k \) be the minimizer of \( (4) \) at \( \omega_k \). Since \( \{D_k\} \) lies in a compact set, there exists an infinite index set \( \mathcal{X}' \subset \mathcal{X} \) such that \( D_k \to D^* \) on \( \mathcal{X}' \). By uniform continuity on compact sets,
\[ v_{D_k}(\omega) < \hat{\mu}^* - \delta, \]
for all \( k \in \mathcal{K} \), in some fixed neighborhood \( V \) of \( \hat{\omega} \). On the other hand, in view of Theorem 1, \( \bar{\zeta}_k > \bar{\mu}^* - \delta \) for \( k \) large enough. Thus, for \( k \in \mathcal{K} \) large enough, \( V \) is taken out of \( \Omega_k \) in Step 2 of the algorithm, which contradicts the fact that \( \omega_k \to \hat{\omega} \) on \( \mathcal{K} \).

5. Numerical experiments

The algorithm was implemented\(^{10}\) in MATLAB\(^{\text{TM}}\). In the implementation, a mid-point rule is used until \( \Omega_k \) is reduced to just one interval and enough information has been accumulated to compute an interpolating function. All numerical experiments were run on a Sun UltraSparc 10, 300 MHz machine with 128 Meg of RAM, running Solaris 2.5 operating system.

The performance of the new algorithm was first compared to that of the “skewed \( \mu \)” approach mentioned in the introduction. For the former, function \( m_k \) of the \( \mu \)-Analysis and Synthesis Toolbox (Balas, Doyle, Glover, Packard & Smith, 1993) was used for computing \( \hat{\zeta}_{k+1} \) and \( D_{k+1} \) in Step 1. The GEVP in the latter was solved using the LMI Control Toolbox Version 1.0.4 (Gahinet, Nemirovski, Laub & Chilali 1995). Tolerance parameters were set so as to obtain roughly the same accuracy with both approaches; in particular, option \( C \) was used in \( m_k \). On 100 randomly generated\(^{11}\) stable systems with 10 variables, 5 inputs, 5 outputs and 4 complex uncertainty block problems, there were 5 instances where the new approach yielded a \( \mu \)-norm upper bound 1% or more lower than the upper bound obtained with the skewed-\( \mu \) approach; in one of these instances, the difference was about 33%. The computation time was always less with the new algorithm, often by more than an order of magnitude.

To illustrate the behavior of the new algorithm, we then applied it to an example taken from Balas et al. (1993). The system is a model of an experimental aircraft (the HIMAT vehicle) with an \( \mathcal{H}_\infty \) sub-optimal pitch axis controller. Uncertainty and performance weightings are included in the model, for a total of 16 state variables. The resultant system is not strictly proper, hence a state space realization will have a non-zero feed-through matrix \( E \). Redefine the Hamiltonian matrix as

\[
H(\xi, A, B, C, E) = \begin{bmatrix}
A & 0 \\
0 & -A^*
\end{bmatrix} + \begin{bmatrix}
B & 0 \\
0 & -C^*
\end{bmatrix} \\
\begin{bmatrix}
-E & \xi I \\
\xi I & -E^*
\end{bmatrix}^{-1}
\begin{bmatrix}
C & 0 \\
0 & B^*
\end{bmatrix}
\]

and make sure that the initial frequency \( \omega_0 \) is such that \( \hat{\mu}(\omega_0) > \sigma_1(D) \). Then the algorithm can handle such problems with no other changes. The uncertainty structure \( \Delta \) consists of two \( 2 \times 2 \) full blocks (robust performance problem). Fig. 4 shows a plot of \( \hat{\mu}_m(P(j\omega)) \) (in this case, equal to \( \mu_m(P(j\omega)) \)) as a function of frequency. The algorithm may now be applied to compute \( \| P \|_\mu \).

In Table 1 the progress of the algorithm from iteration to iteration is shown. All quantities are as defined earlier. An asterisk in the column labeled I indicates that rational interpolation was used to generate \( \omega_k \). If \( \Omega_k \) contains multiple intervals, or if the algorithm has not yet accumulated enough information, the mid-point rule is used to generate \( \omega_k \). It appears from the table that the algorithm exhibits quadratic convergence. Note that, in the process of computing \( \hat{\mu}^* \), \( \hat{\mu} \) has been evaluated at only a few points, which demonstrates great computational savings in comparison with a grid search. The total computation time was about 2.5 s. In comparison, on the same problem, the “skewed \( \hat{\mu} \)” approach took over 50 s. Interestingly, it yielded the same value for \( \hat{\mu}^* \).

6. Extension to mixed-\( \mu \)

The algorithm is readily extended to compute an upper bound on the mixed-\( \mu \) norm, that is, the \( \mu \)-norm for systems with mixed dynamic and real parametric uncertainty. (See Feron (1997) for a sophisticated grid-based approach to this computation, accounting for possible discontinuities.) Two expressions for the widely used “D-G” upper bound \( \hat{\mu}_m(M) \) are as follows (see, e.g., Chapter 18 in Zhou & Doyle, 1998)

\[
\hat{\mu}_m(M) = \inf_{D \in \mathcal{D}, G \in \mathcal{G}} \left\{ \xi: \| \left( \left( DMD^{-1} \right)^{-1} - jG \right) \times (I + G^2)^{-1/2} \|_\infty \leq 1 \right\},
\]

\( \hat{\mu} \)-plot for the closed-loop system.
\( \tilde{\mu}_k(M) = \inf_{D \in \mathcal{D}, G \in \mathcal{G}} \{ \xi^* : M^*DM + j(GM - M^*G) \\
- \xi^2D = 0 \} \), \hspace{1cm} (7)

where \( \mathcal{D} \) and \( \mathcal{G} \) are determined by the structure. Expression (7) is an efficiently solvable LMI problem (see Boyd et al., 1994) and, given \( x, D, G \) satisfying the constraint in (7), the values \( \tilde{\xi}_k = x, \; D = \tilde{D}^{1/2}, \) and \( G = (1/2)\tilde{D}^{-1/2}G\tilde{D}^{-1/2} \) satisfy the constraint in (6) for the same \( M \) (Zhou & Doyle, 1998). Moreover, given \( \tilde{\xi}, D, \) and \( G, \) it follows from (6), with \( M = P(j\omega) \), that \( \tilde{\mu}_x(P(j\omega)) \leq \tilde{\xi} \) for all \( \omega \) such that

\[ \bar{\sigma}(F(j\omega)) \leq 1, \hspace{1cm} (8) \]

with

\[ F(j\omega) = \left( \frac{DP(j\omega)D^{-1}}{\tilde{\xi}} - jG \right)(I + G^2)^{-1/2}. \hspace{1cm} (9) \]

Note that, given a realization \( (A, B, C, E) \) for \( P(j\omega) \), \( F(j\omega) \) can be expressed as \( C_F(j\omega I - A)^{-1}B_F + E_F \), where \( B_F = \tilde{\xi}^{-1/2}BD^{-1}(I + G^2)^{-1/2}, \; C_F = \tilde{\xi}^{-1/2}DC, \) and \( E_F = (\tilde{\xi}^{-1}DED^{-1} - jG)(I + G^2)^{-1/2} \) are obtained by inspection from (9). Thus, the end-points of intervals where (8) holds correspond to imaginary eigenvalues of \( H = H(1, A, B_F, C_F, E_F) \), where \( H \) is as in (5). When \( E = 0 \) (no feedthrough term in \( P \)), this reduces to

\[ H = \begin{bmatrix}
A & 0 \\
0 & -A^*
\end{bmatrix}
+ \frac{1}{\tilde{\xi}} \begin{bmatrix}
-jBD^{-1}GDC & BD^{-2}B^* \\
-jC^*D(I + G^2)DC & jC^*DGD^{-1}B^*
\end{bmatrix}. \]

With this in hand, the algorithm discussed in the previous sections is readily extended to the mixed-\( \mu \) case.

However, the \( \hat{\mu}(\cdot) \) curve does not enjoy the same regularity properties as in the purely complex case. As a result, the interpolation rules for updating \( \phi_k \) may not be appropriate and quadratic convergence may be lost. In the extreme case where the supremum of \( \hat{\mu}(\omega) \) is strictly larger than the supremum of its lower envelope (function whose epigraph is the closure of that of \( \hat{\mu}(\omega) \)), convergence of \( \xi_k \) to \( \hat{\mu}^* \) (Theorem 1) may even be lost unless frequency points where \( \hat{\mu} \) is discontinuous (including possibly \( \omega = \infty \)) are checked separately, and the initial value of \( \xi_k \) is set accordingly. This can happen even in the trivial case of “real \( \mu \)” (real uncertainty) of a scalar transfer function such as \( p(s) = 1/(s^3 + (3/2)s^2 + s + 1) \), since \( \hat{\mu}(p(0)) = 1, \; \mu(\pm j) = 2, \) and \( \hat{\mu}(p(j\omega)) = 0 \) for \( \omega \in [0, \pm 1] \). Indeed, the real \( \mu \) of a scalar is equal to its upper bound \( \hat{\mu} \) and is the absolute value of the scalar if the scalar is real, and zero otherwise.

7. Concluding remarks

The algorithm just outlined in this note can be refined in various ways. For instance, in the process of minimizing \( \nu_G(\phi_k) \) to evaluate \( \tilde{\mu}(\phi_k) \), if it is found that the minimum value is less than \( \tilde{\xi}_k \), then it is unnecessary to compute it (or \( D_k \)) with great accuracy.

The new algorithm can be extended to discrete-time systems. The Hamiltonian eigenvalue problem \( H(\xi, A, B, C) \) is replaced by a simplectic eigenvalue problem \( S(\xi, A, B, C) \) which will have an eigenvalue \( e^{j\omega} \) on the unit circle if and only if \( \xi \) is a singular value of \( C(e^{j\omega}I - A)^{-1}B \). We refer the reader to Hinrichsen and Son (1989), for the appropriate formulas. Alternatively, a corresponding continuous-time quadruple \( (A, B, C, D) \) can be obtained by means of a bilinear transformation.

\begin{table}[h]
\centering
\caption{Results for example system}
\begin{tabular}{|c|c|c|c|c|}
\hline
\( k \) & \( \phi_{k-1} \) & \( \phi_k \) & \( |\phi_{k-1} - \phi_k| \) & \( |\phi_k - \phi_k^*| \) & I \\
\hline
1 & 0 & 0.7959766 & 8.0 \times 10^{-1} & 8.9 \times 10^{-1} \\
2 & 3790.8315 & 0.7959766 & 3.7 \times 10^{-4} & 8.9 \times 10^{-1} \\
3 & 367.3152 & 0.7959766 & 2.9 \times 10^{-2} & 8.9 \times 10^{-1} \\
4 & 129.91042 & 1.3439461 & 5.0 \times 10^{-1} & 3.4 \times 10^{-1} \\
5 & 64.955208 & 1.6041191 & 1.6 \times 10^{-1} & 7.9 \times 10^{-2} \\
6 & 76.939264 & 1.6791253 & 3.5 \times 10^{0} & 3.7 \times 10^{-3} & * \\
7 & 80.420400 & 1.6827812 & 2.1 \times 10^{-1} & 1.3 \times 10^{-5} & * \\
8 & 80.450844 & 1.6827940 & 5.6 \times 10^{-4} & 9.5 \times 10^{-11} & * \\
\hline
\end{tabular}
\end{table}

References


