

OVERESTIMATION IN LINEAR INTERVAL EQUATIONS*

A. NEUMAIER

Abstract. Bounds are given for the overestimation of the solution set of a system of linear interval equations by various inclusion intervals for this set. In particular, a theorem of Gay [this Journal, 19 (1982), pp. 858-870] on the quadratic approximation property of the preconditioned fix-point inverse is strengthened.

Key words. linear interval equations, overestimation, quadratic approximation

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1. Introduction. A well-studied problem in interval analysis is the construction of enclosures for the set of solutions of the equations

$$(1) \quad \tilde{A}\tilde{x} = \tilde{b} \quad (\tilde{A} \in A, \tilde{b} \in b),$$

where A is a regular interval matrix (i.e. all $\tilde{A} \in A$ are nonsingular) and b is an interval vector. While preconditioned Gauss elimination [Hansen [6]] and preconditioned interval iteration [Krawczyk [7]] turned out to be practically useful methods, the quality of the computed enclosures was paid little attention, apart from the special case of M -matrices. There, Barth and Nuding [2] showed that the limit $A^{\infty}b$ of the Jacobi iteration for (1) agrees with the hull $A^{\infty}b$ of the solution set of (1), and Barth and Nuding [2] and Beck [3] showed the same for the result $A^G b$ of Gauss elimination provided that the right-hand side b satisfies $0 \leq b, 0 \leq b$, or $0 \in b$.

In the general case an asymptotic investigation of the overestimation of preconditioned Gauss elimination has been made by Müller [10] (cf. the exposition in Alefeld and Herzberger [1, Chap. 16]). Moreover, exponential worst case overestimation of Gauss elimination without preconditioning was demonstrated by Wongwises [15] and Schätzle [14] under natural simplifying assumptions. For the overestimation in case of preconditioned interval iteration a neat theorem of Gay [5] is available; however, his proof contains a gap. The present paper partly owes its existence to an attempt to repair the proof; in fact we prove a strengthened version of Gay's theorem (Theorem 3). We also prove a new optimality result for the midpoint inverse as preconditioning matrix (Theorem 4) and derive another enclosure for the hull $A^{\infty}b$ with small overestimation, based on an approximate inverse C and an approximate solution \tilde{x} (Theorem 1). Using the concept of a skew vector, we are even able to approximate the solution set of (1) itself by a skew vector with small overestimation (Theorem 2).

The terminology follows Neumaier [12]; it is assumed that the reader is acquainted with §§ 1-4 and 7 of [12]. However, for convenience of the reader we repeat some basic notation. We denote by \mathbb{R}, \mathbb{R}^n , and $\mathbb{R}^{n \times n}$ the set of real numbers, n -dimensional column vectors and $n \times n$ -matrices, and by \mathbb{R}, \mathbb{R}^n , and $\mathbb{R}^{n \times n}$ the set of real closed intervals, n -dimensional interval vectors and $n \times n$ interval matrices, respectively. If $x = [x, \bar{x}] \in \mathbb{R}^n$ then we write $\text{mid } x = \tilde{x} := (x + \bar{x})/2$ for the midpoint, $\text{rad } x = \rho(x) := (x - \bar{x})/2$ for the radius, $|x| := \sup \{|\tilde{x}| \mid \tilde{x} \in x\}$ for the absolute value, and $\tilde{x} := x - \tilde{x} = [-\rho(x), \rho(x)]$ for the radial part of x ; similar definitions apply to interval

* Received by the editors June 12, 1985; accepted for publication (in revised form) March 31, 1986. † Institut für Angewandte Mathematik, Universität Freiburg, D-7800 Freiburg, West Germany.

¹ To emphasize the nature of C as a fixed real matrix, we do not write \tilde{C} which would be the "generic" notation.

A. NEUMAIER

matrices. The Ostrowski operator (\cdot) associates with a square matrix $A \in \mathbb{R}^{n \times n}$ the matrix $A' = (A)$ with entries $A'_{ik} := \min \{|\tilde{a}| \mid \tilde{a} \in A_{ik}\}$, $A'_{ik} := -|A_{ik}|$ for $i \neq k$. In terms of the Jacobi splitting $A = D - E$ of A , defined by $D_{ii} := A_{ii}$, $E_{ii} := 0$ and $D_{ik} := 0$, $E_{ik} := -A_{ik}$ for $i \neq k$ we have

$$(A) = (D) - |E|, \quad |A| = |D| + |E| \cong (D) + |E|.$$

Sum and difference of two bounded subsets $\Sigma_1, \Sigma_2 \subseteq \mathbb{R}^n$ are defined as $\Sigma_1 \pm \Sigma_2 := \{\tilde{x} \pm \tilde{y} \mid \tilde{x} \in \Sigma_1, \tilde{y} \in \Sigma_2\}$, and the hull of a bounded subset $\Sigma \subseteq \mathbb{R}^n$ is defined as $\text{Hull } \Sigma := [\inf \Sigma, \sup \Sigma]$, the smallest interval vector containing Σ . The (vector valued) distance of two interval vectors $x, y \in \mathbb{R}^n$ is defined as

$$(2) \quad q(x, y) := \inf \{q \in \mathbb{R}^n \mid q \geq 0, x \subseteq y + [-q, q], y \subseteq x + [-q, q]\}.$$

One immediately establishes the following properties of the distance.

LEMMA 1. Let $x, y, z \in \mathbb{R}^n$. Then

$$\begin{aligned} q(x, y) = 0 & \text{ iff } x = y, \\ q(y, x) & = q(x, y), \\ q(x, z) & \leq q(x, y) + q(y, z). \end{aligned}$$

LEMMA 2. Let $x, y \in \mathbb{R}^n$. Then

$$\begin{aligned} (3) \quad x \subseteq y & \text{ iff } |x - y| \leq \rho(y) - \rho(x), \\ (4) \quad q(x, y) & = |\tilde{x} - \tilde{y}| + |\rho(x) - \rho(y)|. \end{aligned}$$

Proof. Equation (3) is well known, and the equivalence of (2) and (4) is proved in Lemma 1 of Neumaier [13]. \square

Note that often (4) is taken as a definition of the distance of intervals.

LEMMA 3. If $f(\xi_1, \dots, \xi_m)$ is an arithmetic expression in m variables such that the variables ξ_1, \dots, ξ_m (where $n \leq m$) occur only once then for $x_1, \dots, x_m \in \mathbb{R}$, $\rho(x_i) = 0$ for $i > n$,

$$f(x_1, \dots, x_m) = \{f(\xi_1, \dots, \xi_m) \mid \xi_i \in x_i (i = 1, \dots, m)\}.$$

In particular, if $a, b \in \mathbb{R}^n$ then

$$a + b = \{\tilde{a} + \tilde{b} \mid \tilde{a} \in a, \tilde{b} \in b\},$$

and if $A \in \mathbb{R}^{n \times n}$, $\tilde{x} \in \mathbb{R}^n$ then

$$A\tilde{x} = \{\tilde{A}\tilde{x} \mid \tilde{A} \in A\}.$$

Proof. Well known; see [1] or [11]. \square

We define a skew vector as a set of the form

$$(5) \quad \{\tilde{x} + C\tilde{r} \mid \tilde{r} \in r\},$$

where $\tilde{x} \in \mathbb{R}^n$, $C \in \mathbb{R}^{n \times n}$ and $r \in \mathbb{R}^n$. Geometrically, a skew vector is (possibly degenerate) parallelepiped with axes parallel to the columns of C . Skew vectors occur in the literature on interval differential equations, hidden behind the concept of a moving coordinate frame (see e.g. Moore [11], Eijgenraam [4]). The hull of the skew vector (5) is easily seen to be the interval $\tilde{x} + Cr$. In case of large off-diagonal entries in C the distance between a skew vector and its hull may be quite large; cf. Fig. 1, which shows a skew vector (5) and its hull, where

$$\tilde{x} = \begin{pmatrix} 2 \\ 2 \end{pmatrix}, \quad C = \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \quad r = \begin{pmatrix} [-1, 1] \\ [-1, 1] \end{pmatrix}.$$

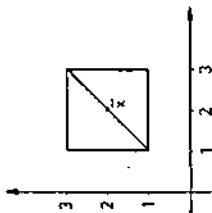


FIG. 1

Therefore it can be expected that, in general, the overestimation of a set Σ by an enclosing interval (e.g. its hull) can be considerably reduced by enclosing Σ in a suitable skew vector.

2. The solution set of linear interval equations. Let $A \in \mathbb{R}^{n \times n}$, $b \in \mathbb{R}^n$, and suppose that A is regular. We are interested in narrow enclosures of the set

$$\Sigma(A, b) := \{\tilde{A}^{-1}\tilde{b} \mid \tilde{A} \in A, \tilde{b} \in b\}$$

of solutions of the linear system of interval equations

$$\tilde{A}\tilde{x} = \tilde{b} \quad (\tilde{A} \in A, \tilde{b} \in b).$$

The optimal interval enclosure is

$$A''b := \Omega[\tilde{A}^{-1}\tilde{b} \mid \tilde{A} \in A, \tilde{b} \in b];$$

but in practice often only cruder enclosures are available. In this section we consider a particular enclosure based on the residual $r := b - A\tilde{x}$ of an approximation \tilde{x} to the solution.

THEOREM 1. Suppose that $C \in \mathbb{R}^{n \times n}$ is such that CA is regular. Let $\tilde{x} \in \mathbb{R}^n$, $r := b - A\tilde{x}$, and suppose that $d \in \mathbb{R}^n$ satisfies

$$A''((I - AC)r) \subseteq d \quad \text{or} \quad (CA)''((I - CA)(Cr)) \subseteq d.$$

Then $x := \tilde{x} + Cr + d$ satisfies

$$(1) \quad A''b \subseteq x \subseteq A''b + (d - d);$$

in particular,

$$(2) \quad q(A''b, x) \subseteq 2\rho(d),$$

$$(3) \quad 0 \subseteq \text{rad } x - \text{rad } A''b \subseteq 2\rho(d).$$

Proof. By assumption, C and A are regular. Let $\tilde{A} \in A$, $\tilde{b} \in b$ and put

$$(4) \quad \tilde{r} := \tilde{b} - \tilde{A}\tilde{x}, \quad \tilde{y} := \tilde{x} + Cr, \quad \tilde{z} := \tilde{A}^{-1}\tilde{b}.$$

Then $\tilde{z} - \tilde{y} = \tilde{A}^{-1}\tilde{b} - \tilde{x} - C(\tilde{b} - \tilde{A}\tilde{x}) = \tilde{A}^{-1}(I - \tilde{A}C)(\tilde{b} - \tilde{A}\tilde{x})$, hence

$$(5) \quad \tilde{z} - \tilde{y} = \tilde{A}^{-1}(I - \tilde{A}C)\tilde{r} = (C\tilde{A})^{-1}(I - C\tilde{A})C\tilde{r} \in d.$$

Since $\tilde{y} \in \tilde{x} + Cr$, this implies (1). To prove (2) and (3), we use the fact that Lemma 3 implies

$$r = \{\tilde{b} - \tilde{A}\tilde{x} \mid \tilde{A} \in A, \tilde{b} \in b\};$$

hence every $\tilde{r} \in r$ and every \tilde{y} in

$$(6) \quad \tilde{\Sigma} := \{\tilde{x} + Cr \mid \tilde{r} \in r\}$$

can be written in the form (4). Therefore $\Sigma(A, B) \subseteq \tilde{\Sigma} + d \subseteq \tilde{x} + Cr + d = x$, and $\tilde{\Sigma} \subseteq \Sigma(A, B) - d \subseteq A''b - d$, whence $\tilde{x} + Cr \subseteq A''b - d$, and $x = \tilde{x} + Cr + d \subseteq A''b - d + d$. This implies (1), and (2). (3) are immediate consequences. \square

For practical application, C should be an approximation of the inverse of some $\tilde{A} \in A$ so that $CA = I$, and \tilde{x} should be an approximation of $\tilde{A}^{-1}\tilde{b}$ for some $\tilde{b} \in b$ so that $A\tilde{x} = \tilde{b}$. If the residuals $I - CA$ and $b - A\tilde{x}$ are of order $O(\epsilon)$ then d can be taken of order $O(\epsilon^2)$, so that (1) overestimates the hull by only $O(\epsilon^2)$. Note that the amount of overestimation can be computed a posteriori from (3).

The proof of Theorem 1 suggests that $\Sigma(A, b)$ itself may be enclosed with little overestimation by a skew vector. This is indeed the case.

THEOREM 2. Suppose that $C \in \mathbb{R}^{n \times n}$ is such that AC is regular. Let $\tilde{x} \in \mathbb{R}^n$, $r := b - A\tilde{x}$ and suppose that

$$(AC)''((I - AC)r) \subseteq e \in \mathbb{R}^n.$$

Then $\Sigma_0 := \{\tilde{x} + C\tilde{r} \mid \tilde{r} \in r + e\}$ satisfies

$$(7) \quad \Sigma(A, b) \subseteq \Sigma_0 \subseteq \Sigma(A, b) + |C|(e - e).$$

Proof. As before,

$$\tilde{e} := C^{-1}(\tilde{z} - \tilde{y}) = (\tilde{A}C)^{-1}(I - \tilde{A}C)\tilde{r} \in e,$$

so that $\tilde{z} = \tilde{y} + C\tilde{e} = \tilde{x} + C(\tilde{r} + \tilde{e})$, which implies $\Sigma(A, b) \subseteq \Sigma_0$. Again as before, $\tilde{\Sigma} \subseteq \Sigma(A, b) - \{C\tilde{e} \mid \tilde{e} \in e\}$ so that $\Sigma_0 \subseteq \tilde{\Sigma} + \{C\tilde{e} \mid \tilde{e} \in e\} \subseteq \Sigma(A, b) + \{C(\tilde{e} - \tilde{e}) \mid \tilde{e} \in e\} \subseteq \Sigma(A, b) + |C|(e - e)$. \square

Again, if the residuals $I - AC$ and $b - A\tilde{x}$ are of order $O(\epsilon)$ then e can be chosen of order $O(\epsilon^2)$, so that (6) overestimates the solution set by only $O(\epsilon^2)$. Thus (6) can be expected to be a very accurate enclosure for the solution set. Note that the amount of overestimation can be calculated a posteriori from (7).

3. The preconditioned fixpoint inverse. $A \in \mathbb{R}^{n \times n}$ is called an H -matrix if $(A)_{ii} > 0$ for some $u > 0$; cf. Neumaier [12, § 2.5]. In [12], the fixpoint inverse of an H -matrix $A \in \mathbb{R}^{n \times n}$ was defined as the sublinear map $A'' : \mathbb{R}^n \rightarrow \mathbb{R}^n$ which maps $b \in \mathbb{R}^n$ to the unique solution $x = A''b$ of the fixpoint equation

$$x = L^r(b + Ex),$$

where $A = L - E$ is a direct splitting of A . By [12, Prop. 11], the vector $x = A''b$ is an enclosure of the hull $A''b$. Here we consider the preconditioned equations

$$CA\tilde{x} = C\tilde{b} \quad (\tilde{A} \in A, \tilde{b} \in b),$$

and, assuming that CA is an H -matrix, the corresponding enclosure $x = (CA)''(C\tilde{b})$ for $(CA)''(C\tilde{b})$. The norm used in the following is the maximum norm.

THEOREM 3. Suppose that $C \in \mathbb{R}^{n \times n}$ is such that CA is an H -matrix. Then

$$(1) \quad z = A''b \subseteq (CA)''(C\tilde{b}) \subseteq (CA)''(Cb) \subseteq \tilde{z} + (CA)''|CA|\tilde{z}.$$

Moreover, if CA is strictly diagonally dominant, i.e.,

$$\beta := \max_i \frac{|CA|_{ii}}{\sum_{k \neq i} |CA|_{ki}} < 1$$

then

$$(2) \quad \|\text{rad } A''b\| \leq \|\text{rad } (CA)''(C\tilde{b})\| \leq \|\text{rad } (CA)''(Cb)\| \leq \frac{1+\beta}{1-\beta} \|\text{rad } A''b\|.$$

Proof. Let $CA = D - E$ be the Jacobi splitting of CA . Then $x := (CA)^F(Cb)$ satisfies the equation

$$(3) \quad x = D^F(Cb + Ex).$$

Fix $i \in \{1, \dots, n\}$. The splitting equation implies

$$D_{ii} = \sum_j C_{ij}A_{jk}, \quad E_{ik} = -\sum_j C_{ij}A_{jk} \quad (k \neq i),$$

and (3) gives

$$x_i = \sum_k \frac{C_{ik}b_k + E_{ik}x_k}{D_{ii}}.$$

In these expressions, each variable apart from the $C_{ij} \in \mathbb{R}$ occurs only once (note that i is fixed); hence these expressions equal the range (Lemma 3), and for each $\xi \in x_i$ we can find $\tilde{x} \in x$, $\tilde{b} \in b$, $\tilde{D} \in D$, $\tilde{E} \in E$, $\tilde{A} \in A$ such that $CA = \tilde{D} - \tilde{E}$ and $\tilde{y} := \tilde{D}^{-1}(Cb + \tilde{E}\tilde{x})$ satisfies $\tilde{y}_i = \xi$.

Now $\tilde{z} := \tilde{A}^{-1}\tilde{b} \in A^{-1}b = z$ and $\tilde{y} = \tilde{D}^{-1}(C\tilde{A}\tilde{z} + \tilde{E}\tilde{x}) = \tilde{D}^{-1}((\tilde{D} - \tilde{E})\tilde{z} + \tilde{E}\tilde{x}) = \tilde{z} + \tilde{D}^{-1}\tilde{E}(\tilde{x} - \tilde{z}) \in z + D^F E(x - z)$; therefore $\xi \in z + D^F E(x - z)$, and since $\xi \in x_i$ and i were arbitrary we find

$$x \subseteq z + D^F E(x - z).$$

This implies

$$\begin{aligned} |x - \tilde{z}| &\subseteq |z + D^F E(x - z)| \subseteq \rho(z) + \langle D \rangle^{-1} |E| |x - z|, \\ \langle D \rangle |x - \tilde{z}| &\subseteq \langle D \rangle \rho(z) + |E| (|x - \tilde{z}| + \rho(z)). \end{aligned}$$

Therefore

$$(4) \quad |x - \tilde{z}| \subseteq (\langle D \rangle - |E|)^{-1} (\langle D \rangle + |E|) \rho(z).$$

In particular,

$$|x - \tilde{z}| \subseteq \langle CA \rangle^{-1} |CA| \rho(z), \quad x \in \tilde{z} + \langle CA \rangle^{-1} |CA| \xi$$

and (1) follows. Finally if CA is strictly diagonally dominant then $B := \langle D \rangle^{-1} |E|$ satisfies $\|B\| \leq \beta < 1$ whence

$$\begin{aligned} \| \text{rad } x \| &\leq \| x - \tilde{z} \| \leq \| (\langle D \rangle - |E|)^{-1} (\langle D \rangle + |E|) \rho(z) \| \\ &= \| (I - B)^{-1} (I + B) \rho(z) \| \leq \frac{1 + \beta}{1 - \beta} \| \rho(z) \| \end{aligned}$$

which implies (2). \square

Remarks. 1) If $\rho(A) = O(\epsilon)$ and $C = \tilde{A}^{-1}$ then $CA = I + \tilde{A}^{-1}A = I + O(\epsilon)$ whence $\beta = O(\epsilon)$. From (2) we therefore get

$$\frac{\text{rad}((CA)^F(Cb))}{\text{rad } A^{-1}b} = 1 + O(\epsilon),$$

i.e., very little overestimation.

2) For a computable bound on the overestimation we simply replace $A^{-1}b$ in (2) by the slightly bigger $(CA)^F(Cb)$.

3) Theorem 3 implies that preconditioning with a good approximation C of \tilde{A}^{-1} (resulting in a small β , say $\beta \leq \frac{1}{2}$) increases the radius of $A^{-1}b$ only by a small factor $\frac{1}{1 - \beta}$ (usually gives excellent enclosures. An asymptotic result of this type is already known for some time: Miller [10] proved

$$\text{rad}((CA)^F(Cb)) = \text{rad } A^{-1}b + O(\epsilon^2)$$

if $C = \tilde{A}^{-1}$ and $\rho(A)$, $\rho(b) = O(\epsilon)$; cf. the exposition in Alefeld and Herzberger [1, Chap. 16].

4) Since $q(x, z) = |\tilde{x} - \tilde{z}| + |\rho(x) - \rho(z)| = |x - \tilde{z}| - \rho(z)$ (note $\rho(x) \geq \rho(z)$ since $z \subseteq x$), we get from (4) the relation

$$q(x, z) \leq 2(\langle D \rangle - |E|)^{-1} |E| \rho(z),$$

which in the strictly diagonally dominant case implies

$$\|q((CA)^F(Cb), A^{-1}b)\| \leq \frac{2\beta}{1 - \beta} \|\text{rad } A^{-1}b\|.$$

5) If $\|I - CA\| \leq \beta_0 < 1$ then the iteration $x^{i+1} = Cb + Ex^i$, where $E := I - CA$, converges to a limit x^* satisfying $x^* = Cb + Ex^*$ and

$$z \subseteq (CA)^F(Cb) \subseteq x^*;$$

see Mayer [9], Gay [5] and Neumaier [12]. Proceeding as in the proof of Theorem 3 but with the splitting $CA = I - E$ one obtains the relation

$$z \subseteq x^* \subseteq \tilde{z} + (I - |E|)^{-1} (I + |E|) \tilde{z},$$

and as in Remark 4 one gets

$$(5) \quad \|q(x^*, A^{-1}b)\| \leq \frac{2\beta_0}{1 - \beta_0} \|\text{rad } A^{-1}b\|.$$

6) Relation (5) improves Theorem 4.1 of Gay [5], who states almost the same inequality for $\beta_0 < \frac{1}{2}$, but with the denominator $1 - 3\beta_0$. However his proof seems to have a gap: In [5, p. 868, line 4] one can conclude $w(x^i) = w(d) + w(\langle E \rangle d)$ but it is not clear whether this implies the required relation $w(x^i) = w(d) + \sup(\langle \langle E \rangle d \rangle)$.

7) It would be interesting to have a result similar to Theorem 3 for preconditioned Gauss elimination.

4. On the choice of the preconditioning matrix. In Neumaier [12], the optimal choice of the preconditioning matrix C was discussed. It was shown that if $\|I - CA\| < 1$ for some $C \in \mathbb{R}^{n \times n}$ then $\|I - CA\|$ takes its minimum for $C = \tilde{A}^{-1}$. Here we present another optimality criterion based on an upper bound for the preconditioned fixpoint inverse or Gauss inverse. We remind the reader that $A^{-1}b$ denotes the result of interval Gauss elimination applied without pivoting to the matrix A and the right-hand side b ; cf. [12, Chap. 5].

THEOREM 4. Let $A \in \mathbb{R}^{n \times n}$, $b \in \mathbb{R}^n$, and $C \in \mathbb{R}^{n \times n}$, and suppose that CA is an H -matrix. Then

$$(1) \quad |(CA)^F(Cb)| \subseteq (CA)^{-1}|Cb|,$$

$$(2) \quad |(CA)^G(Cb)| \subseteq (CA)^{-1}|Cb|,$$

and if $\tilde{b} = 0$ then equality holds in (1). Moreover, $\tilde{A}^{-1}A$ is an H -matrix and

$$(3) \quad \langle \tilde{A}^{-1}A \rangle^{-1} | \tilde{A}^{-1}b | \subseteq \langle CA \rangle^{-1} | Cb |.$$

Proof. By Theorem 9 of [12], $|(CA)^F| = (CA)^{-1}$ so that $|(CA)^F(Cb)| \cong |(CA)^F| |Cb| = (CA)^{-1} |Cb|$, and (1) holds. If $\beta = 0$ then $a = Cb$ satisfies $\hat{a} = 0$, and since $(CA)^F$ is normal [12, Chap. 4], (R8) in [12] implies $(CA)^F(Cb) = |(CA)^F| |(Cb)| = (CA)^{-1} |Cb|$, so that (1) holds with equality. If, instead, Theorems 3 and 4 of [12] are used one obtains in the same way relation (2).

We now observe that CA is an H -matrix, so that by [12, Thm. 5], $\hat{A}^{-1}A$ is an H -matrix. To prove (3), let $u > 0$ be such that $\langle CA \rangle u > 0$. Then $0 \notin \langle CA \rangle_i$ for $i = 1, \dots, n$ so that by [12, § 2.2 (O4) and § 1.3, (B9) and (B14) (the last rule contains a misprint and should read $\rho(BA) = \rho(B)|A|$), we have

$$\langle CA \rangle = \langle \hat{CA} \rangle - |C| \rho(A).$$

Therefore the matrix

$$P := \langle \hat{CA} \rangle^{-1} \langle CA \rangle = I - \langle \hat{CA} \rangle^{-1} |C| \rho(A)$$

satisfies $P \preceq I$, $Pu = \langle \hat{CA} \rangle^{-1} \langle CA \rangle u > 0$. Hence P is an M -matrix; in particular, $P^{-1} \cong 0$. Since

$$\begin{aligned} \langle \hat{A}^{-1}A \rangle &= \langle I + \hat{A}^{-1}\hat{A} \rangle \cong I - |\hat{A}^{-1}| \rho(A) \\ &= I - |\langle \hat{CA} \rangle^{-1} |C| \rho(A) \cong I - \langle \hat{CA} \rangle^{-1} |C| \rho(A) = P, \end{aligned}$$

we get $P \langle \hat{A}^{-1}A \rangle^{-1} \preceq I$; hence $\langle \hat{A}^{-1}A \rangle^{-1} \preceq P^{-1}$, and finally

$$\begin{aligned} \langle \hat{A}^{-1}A \rangle^{-1} |\hat{A}^{-1}b| &\preceq P^{-1} |\hat{A}^{-1}b| = P^{-1} |\langle \hat{CA} \rangle^{-1} Cb| \\ &\preceq P^{-1} \langle \hat{CA} \rangle^{-1} |Cb| = \langle CA \rangle^{-1} |Cb|. \end{aligned} \quad \square$$

Remarks. 1) Inequality (3) describes a minimality property of the upper bound $\langle CA \rangle^{-1} |Cb|$. It would be interesting to know an optimal C for which $|\langle CA \rangle^{-1} |Cb|$ and $|\langle CA \rangle^G(Cb)|$ themselves take their minimum. It seems that at least for H -matrices A and $\langle CA \rangle^G(Cb)$, the optimal choice is not $C = \hat{A}^{-1}$ but $C = \hat{A}^{-1}$ where $\hat{A} \in A$ satisfies

$$\begin{aligned} |\hat{A}_i| &= \max(|\hat{A}_i|, |\hat{A}_i|), \\ |\hat{A}_i| &= \min(|\hat{A}_i|, |\hat{A}_i|) \quad \text{for } i \neq k. \end{aligned}$$

2) A result similar to Theorem 4 is implicitly in Krawczyk [8]: If the spectral radius of $|I - CA|$ is less than one then the expression $(I - |I - CA|)^{-1} |Cb|$ (which is also an upper bound for $\langle CA \rangle^G(Cb)$) takes its minimum for $C = \hat{A}^{-1}$; indeed this follows by multiplying the third inequality in the lemma in [8] with $(e - \hat{f})^{-1}$ which in the present notation is $(I - |I - \hat{A}^{-1}A|)^{-1}$.

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