Numerical Analysis - Part II

Anders C. Hansen

Lecture 21

Iterative methods for linear algebraic systems

The conjugate gradient method

Here it is.

- (A) For any initial vector $\mathbf{x}^{(0)}$, set $\mathbf{d}^{(0)} = \mathbf{r}^{(0)} = \mathbf{b} A\mathbf{x}^{(0)}$;
- (B) For $k \geq 0$, calculate $\mathbf{x}^{(k+1)} = \mathbf{x}^{(k)} + \alpha_k \mathbf{d}^{(k)}$ and the residual

$$\mathbf{r}^{(k+1)} = \mathbf{r}^{(k)} - \alpha_k A \mathbf{d}^{(k)}, \quad \text{with}$$

$$\alpha_k := \{ \mathbf{r}^{(k+1)} \perp \mathbf{d}^{(k)} \} = \frac{(\mathbf{r}^{(k)}, \mathbf{d}^{(k)})}{(A \mathbf{d}^{(k)}, \mathbf{d}^{(k)})}, \quad k \ge 0.$$
(1)

(C) For the same k, the next conjugate direction is the vector

$$\mathbf{d}^{(k+1)} = \mathbf{r}^{(k+1)} + \beta_k \mathbf{d}^{(k)}, \text{ with}$$

$$\beta_k := \{ \mathbf{d}^{(k+1)} \perp A \mathbf{d}^{(k)} \} = -\frac{(\mathbf{r}^{(k+1)}, A \mathbf{d}^{(k)})}{(\mathbf{d}^{(k)}, A \mathbf{d}^{(k)})}, \quad k \ge 0.$$
(2)

Theorem 1 (Properties of CGM)

For every $m \ge 0$, the conjugate gradient method has the following properties.

(1) The linear space spanned by the residuals $\{\mathbf{r}^{(i)}\}$ is the same as the linear space spanned by the conjugate directions $\{\mathbf{d}^{(i)}\}$ and it coincides with the space spanned by $\{A^i\mathbf{r}^{(0)}\}$:

$$\mathrm{span}\{\bm{r}^{(i)}\}_{i=0}^m = \mathrm{span}\{\bm{d}^{(i)}\}_{i=0}^m = \mathrm{span}\{A^i\bm{r}^{(0)}\}_{i=0}^m \,.$$

- (2) The residuals satisfy the orthogonality conditions: $(\mathbf{r}^{(m)}, \mathbf{r}^{(i)}) = (\mathbf{r}^{(m)}, \mathbf{d}^{(i)}) = 0$ for i < m.
- (3) The directions are conjugate (A-orthogonal): $(\mathbf{d}^{(m)}, \mathbf{d}^{(i)})_A = (\mathbf{d}^{(m)}, A\mathbf{d}^{(i)}) = 0$ for i < m.

Proof. We use induction on $m \ge 0$, the assertions being trivial for m = 0, since $\mathbf{d}^{(0)} = \mathbf{r}^{(0)}$ and (2)-(3) are void. Therefore, assuming that the assertions are true for some m = k, we ask if they remain true when m = k + 1.

(1) Formula (2)
$$\mathbf{d}^{(k+1)} = \mathbf{r}^{(k+1)} + \beta_k \mathbf{d}^{(k)}$$

readily implies that equivalence of the spaces spanned by $(\mathbf{r}^{(i)})_0^k$ and $(\mathbf{d}^{(i)})_0^k$, is preserved when k is increased to k+1. Similarly, from $\mathbf{r}^{(k+1)} = \mathbf{r}^{(k)} - \alpha_k A \mathbf{d}^{(k)}$ in (1), and from the inductive assumption $\mathbf{r}^{(k)}, \mathbf{d}^{(k)} \in \operatorname{span}\{A^i\mathbf{r}^{(0)}\}_{i=0}^k$, it follows that $\mathbf{r}^{(k+1)} \in \operatorname{span}\{A^i\mathbf{r}^{(0)}\}_{i=0}^{k+1}$. To see that $A^{k+1}\mathbf{r}^{(0)} \in \operatorname{span}\{\mathbf{r}^{(i)}\}_{i=0}^{k+1}$, since $\alpha_k \neq 0$, the claim follows by (5) if $\mathbf{d}^{(k)}$ has a non-zero component from $A^k\mathbf{r}^{(0)}$, and if not the claim follows from the induction hypothesis.

Proof. Cont. (2) Turning to assertion (2), we need $\mathbf{r}^{(k+1)} \perp \mathbf{r}^{(i)}$ for $i \leq k$, which by (1) is equivalent to

$$\mathbf{r}^{(k+1)} \perp \mathbf{d}^{(i)}$$
 for $i \leq k$.

We have $\mathbf{r}^{(k+1)} \perp \mathbf{d}^{(k)}$ by the definition of α_k in (1), so we need

$$\mathbf{r}^{(k+1)} \stackrel{(1)}{=} \mathbf{r}^{(k)} - \alpha_k A \mathbf{d}^{(k)} \perp \mathbf{d}^{(i)}$$
 for $i < k$,

and this follow from the inductive assumptions $\mathbf{r}^{(k)} \perp \mathbf{d}^{(i)}$ and $A\mathbf{d}^{(k)} \perp \mathbf{d}^{(i)}$.

Proof. Cont. (3) It remains to justify (3), namely that $\mathbf{d}^{(k+1)}$ defined in (2) satisfies

$$\mathbf{d}^{(k+1)} \perp A\mathbf{d}^{(i)}$$
 for $i \leq k$.

The value of β_k in (2) is defined to give $\mathbf{d}^{(k+1)} \perp A\mathbf{d}^{(k)}$, so we need

$$\mathbf{d}^{(k+1)} \stackrel{(2)}{=} \mathbf{r}^{(k+1)} + \beta_k \mathbf{d}^{(k)} \perp A \mathbf{d}^{(i)} \quad \text{for} \quad i < k.$$

By the inductive hypothesis $\mathbf{d}^{(k)} \perp A\mathbf{d}^{(i)}$, hence it remains to establish that $\mathbf{r}^{(k+1)} \perp A\mathbf{d}^{(i)}$ for i < k. Now, the formula (1) yields $A\mathbf{d}^{(i)} = (\mathbf{r}^{(i)} - \mathbf{r}^{(i+1)})/\alpha_i$, therefore we require the conditions $\mathbf{r}^{(k+1)} \perp (\mathbf{r}^{(i)} - \mathbf{r}^{(i+1)})$ for i < k, and they are a consequence of the assertion (2) for m = k+1 obtained previously.

Termination property

Corollary 2 (A termination property)

If the conjugate gradient method is applied in exact arithmetic, then, for any $\mathbf{x}^{(0)} \in \mathbb{R}^n$, termination occurs after at most n iterations. More precisely, termination occurs after at most s iterations, where $s = \dim \operatorname{span}\{A^i\mathbf{r}_0\}_{i=0}^{n-1}$ (which can be smaller than n).

Termination property

Proof. Assertion (2) of Theorem 1 states that residuals $(\mathbf{r}^{(k)})_{k\geq 0}$ form a sequence of mutually orthogonal vectors in \mathbb{R}^n , therefore at most n of them can be nonzero. Since they also belong to the space $\operatorname{span}\{A^i\mathbf{r}_0\}_{i=0}^{n-1}$, their number is bounded by the dimension of that space.

The Krylov subspaces

Definition 3 (The Krylov subspaces)

Let A be an $n \times n$ matrix, $\mathbf{v} \in \mathbb{R}^n$ nonzero, and $m \in \mathbb{N}$. The linear space $K_m(A, \mathbf{v}) := \operatorname{span}\{A^i\mathbf{v}\}_{i=0}^{m-1}$ is called the m-th Krylov subspace of \mathbb{R}^n .

Theorem 4 (Number of iterations in CGM)

Let A > 0, and let s be the number of its distinct eigenvalues. Then, for any \mathbf{v} ,

$$\dim K_m(A, \mathbf{v}) \le s \quad \forall m. \tag{3}$$

Hence, for any A > 0, the number of iterations of the CGM for solving $A\mathbf{x} = \mathbf{b}$ is bounded by the number of distinct eigenvalues of A.

The Krylov subspaces

Proof. Inequality (3) is true not just for positive definite A > 0, but for any A with n linearly independent eigenvectors (\mathbf{u}_i) . Indeed, in that case one can expand $\mathbf{v} = \sum_{i=1}^n a_i \mathbf{u}_i$, and then group together eigenvectors with the same eigenvalues: for each λ_{ν} we set $\mathbf{w}_{\nu} = \sum_{k=1}^{m_{\nu}} a_{i_k} \mathbf{u}_{i_k}$ if $A\mathbf{u}_{i_k} = \lambda_{\nu} \mathbf{u}_{i_k}$. Then

$$\mathbf{v} = \sum_{\nu=1}^{s} c_{\nu} \mathbf{w}_{\nu}, \qquad c_{\nu} \in \{0,1\},$$

hence $A^i \mathbf{v} = \sum_{\nu=1}^s c_\nu \lambda_\nu^i \mathbf{w}_\nu$, thus for any m we get $K_m(A,\mathbf{v}) \subseteq \operatorname{span}\{\mathbf{w}_1,\mathbf{w}_2,\ldots,\mathbf{w}_s\}$, and that proves (3). By Corollary 2, the number of iteration in CGM is bounded by $\dim K_m(A,\mathbf{r}^{(0)})$, hence the final conclusion.

The Krylov subspaces

Remark 5

Theorem 4 shows that, unlike other iterative schemes, the conjugate gradient method is both iterative and direct: each iteration produces a reasonable approximation to the exact solution, and the exact solution itself will be recovered after n iterations at most.

Simplifying the CGM-algorithm

We now simplify and reformulate the CGM-algorithm.

Firstly, we rewrite expressions for the parameters α_k and β_k in (1)-(2) as follows:

$$\alpha_{k} = \frac{(\mathbf{r}^{(k)}, \mathbf{d}^{(k)})}{(\mathbf{d}^{(k)}, A\mathbf{d}^{(k)})} \stackrel{(c)}{=} \frac{\|\mathbf{r}^{(k)}\|^{2}}{(\mathbf{d}^{(k)}, A\mathbf{d}^{(k)})} > 0,$$

$$\beta_{k} = -\frac{(\mathbf{r}^{(k+1)}, A\mathbf{d}^{(k)})}{(\mathbf{d}^{(k)}, A\mathbf{d}^{(k)})} \stackrel{(a)}{=} -\frac{(\mathbf{r}^{(k+1)}, \mathbf{r}^{(k+1)} - \mathbf{r}^{(k)})}{(\mathbf{d}^{(k)}, \mathbf{r}^{(k+1)} - \mathbf{r}^{(k)})} \stackrel{(b)}{=} \frac{\|\mathbf{r}^{(k+1)}\|^{2}}{(\mathbf{d}^{(k)}, \mathbf{r}^{(k)})} \stackrel{(c)}{=} \frac{\|\mathbf{r}^{(k+1)}\|^{2}}{\|\mathbf{r}^{(k)}\|^{2}} > 0.$$

Here, for β , we used in (a) the fact that $A\mathbf{d}^{(k)}$ is a multiple of $\mathbf{r}^{(k+1)} - \mathbf{r}^{(k)}$ by (1), and in (b) orthogonality of $\mathbf{r}^{(k+1)}$ to both $\mathbf{r}^{(k)}, \mathbf{d}^{(k)}$ proved in Theorem 1(2). Then, for both β and α , we used in (c) the property $(\mathbf{d}^{(k)}, \mathbf{r}^{(k)}) = \|\mathbf{r}^{(k)}\|^2$ which follows from (2) with index k+1, taking in account orthogonality $\mathbf{r}^{(k+1)} \perp \mathbf{d}^{(k)}$. Secondly, we let $\mathbf{x}^{(0)}$ be the zero vector.

Standard form of the conjugate gradient method

Here it is.

- (1) Set k = 0, $\mathbf{x}^{(0)} = 0$, $\mathbf{r}^{(0)} = \mathbf{b}$, and $\mathbf{d}^{(0)} = \mathbf{r}^{(0)}$;
- (2) Calculate the matrix-vector product $\mathbf{v}^{(k)} = A\mathbf{d}^{(k)}$ and $\alpha_k = \|\mathbf{r}^{(k)}\|^2/(\mathbf{d}^{(k)}, \mathbf{v}^{(k)}) > 0$;
- (3) Apply the formulae $\mathbf{x}^{(k+1)} = \mathbf{x}^{(k)} + \alpha_k \mathbf{d}^{(k)}$ and $\mathbf{r}^{(k+1)} = \mathbf{r}^{(k)} \alpha_k \mathbf{v}^{(k)}$;
- (4) Stop if $\|\mathbf{r}^{(k+1)}\|$ is acceptably small;
- (5) Set $\mathbf{d}^{(k+1)} = \mathbf{r}^{(k+1)} + \beta_k \mathbf{d}^{(k)}$, where

$$\beta_k = \|\mathbf{r}^{(k+1)}\|^2 / \|\mathbf{r}^{(k)}\|^2 > 0;$$

(6) Increase $k \to k+1$ and go back to (2).

Standard form of the conjugate gradient method

The total work is dominated by the number of iterations, multiplied by the time it takes to compute $\mathbf{v}^{(k)} = A\mathbf{d}^{(k)}$. Thus the conjugate gradient algorithm is highly suitable when most of the elements of A are zero, i.e. when A is *sparse*.

Preconditioning

In $A\mathbf{x} = \mathbf{b}$, we change variables, $\mathbf{x} = P^T \widehat{\mathbf{x}}$, where P is a nonsingular $n \times n$ matrix, and multiply both sides with P. Thus, instead of $A\mathbf{x} = \mathbf{b}$, we are solving the linear system

$$PAP^{T}\widehat{\mathbf{x}} = P\mathbf{b} \Leftrightarrow \widehat{A}\widehat{\mathbf{x}} = \widehat{\mathbf{b}}.$$
 (4)

Note that symmetry and positive definiteness of A imply that $\widehat{A} = PAP^T$ is also symmetric and positive definite since $(\widehat{A}\mathbf{y},\mathbf{y}) = (PAP^T\mathbf{y},\mathbf{y}) = (AP^T\mathbf{y},P^T\mathbf{y}) > 0$. Therefore, we can apply conjugate gradients to the new system. This results in the solution $\widehat{\mathbf{x}}$, hence $\mathbf{x} = P^T\widehat{\mathbf{x}}$. This procedure is called the preconditioned conjugate gradient method and the matrix P is called the preconditioner.

Condition number and convergence rate of CGM

The condition number of a matrix A is the value $\kappa(A) := \|A\| \cdot \|A^{-1}\|$, so for a symmetric positive definite matrix A it is the ratio between its largest and smallest eigenvalues,

$$\kappa(A) = rac{\lambda_{\sf max}(A)}{\lambda_{\sf min}(A)} \geq 1$$
 .

The closer this number is to 1, the faster the convergence is of CGM. More precisely, for the rate of convergnce of CGM, we have the upper estimate

$$\|\mathbf{e}^{(k)}\|_{A} \le 2\rho^{k} \|\mathbf{e}^{(0)}\|_{A}, \qquad \rho = \rho_{A} = \frac{\sqrt{\kappa(A) - 1}}{\sqrt{\kappa(A) + 1}} < 1.$$
 (5)

The main idea of preconditioning is to pick P in (4) so that $\kappa(\widehat{A})$ is much smaller than $\kappa(A)$, thus accelerating convergence.

Preconditioning – Choosing *P*

To this end, we note that the similarity transform $B \to C^{-1}BC$ preserves spectrum, hence

$$\kappa(\widehat{A}) = \kappa(PAP^T) = \kappa(P^{-1}[PAP^T]P) = \kappa(AP^TP),$$

and if we set

$$S^{-1} := P^T P =: (QQ^T)^{-1},$$

then it is suggestive to choose S as an approximation to A which is easy to Cholesky-factorize, i.e., $S = QQ^T$ (or already in this form), and then take $P = Q^{-1}$. Then $AP^TP = AS^{-1}$ is close to identity, hence

$$\kappa(\widehat{A}) = \kappa(AP^TP) \approx \kappa(I) = 1 \quad \Rightarrow \quad \kappa(\widehat{A}) \ll \kappa(A),$$

and the preconditioned system (4) will be solved much faster because of (5).

Preconditioning – Extra cost

Each step in the CGM for solving $A\mathbf{x} = \mathbf{b}$ requires one matrix-vector product $A\mathbf{y}$, so with $P = Q^{-1}$, additional expense in each step of the CGM for the preconditioned system (4) while computing $\widehat{A}\mathbf{y} = PAP^T\mathbf{y}$ is two additional computations

$$\mathbf{u} = P^T \mathbf{y} = Q^{-T} \mathbf{y}, \qquad \mathbf{v} = P \mathbf{z} = Q^{-1} \mathbf{z},$$

for some $\mathbf{y}, \mathbf{z} \in \mathbb{R}^n$, but note that computing $Q^{-1}\mathbf{z}$ is the same as solving the linear system $Q\mathbf{v} = \mathbf{z}$, which is cheap (via forward substitution) as Q is a lower triangular matrix.

Preconditioning – Examples

Example 6

- 1) The simplest choice of S is D = diag A, then $P = D^{-1/2}$ in (4).
- 2) Another possibility is to choose S as a band matrix with small bandwidth. For example, solving the Poisson equation with the five-point formula, we may take S to be the tridiagonal part of A.
- 3) One can also take $P=L^{-1}$, where L is the lower triangular part of A (maybe imposing some changes). For example, for the Poisson equation, with m=20 hence dealing with 400×400 system, we take P^{-1} as the lower triangular part of A, but change the diagonal elements from 4 to $\frac{5}{2}$. Then we get a computer precision after just 30 iterations.

Preconditioning – Examples

For the tridiagonal system $A\mathbf{x} = \mathbf{b}$ below, we choose the preconditioner as follows.

$$A = \begin{bmatrix} 2 - 1 \\ -1 & 2 & \ddots \\ & \ddots & \ddots & -1 \\ & -1 & 2 \end{bmatrix}, \qquad Q = \begin{bmatrix} 1 \\ -1 & 1 \\ & \ddots & \ddots \\ & & -1 & 1 \end{bmatrix},$$

$$S = QQ^{T} = \begin{bmatrix} 1 - 1 \\ -1 & 2 & \ddots \\ & \ddots & \ddots & -1 \\ & & -1 & 2 \end{bmatrix}.$$

Preconditioning – **Examples**

The matrix S coincides with A except at the (1,1)-entry. The matrix

$$\widehat{A} = Q^{-1}AQ^{-T}$$

for the preconditioned CGM has just two distinct eigenvalues, and we recover the exact solution just in two steps. To see the latter, note that \widehat{A} is similar to

$$Q^{-T}Q^{-1}A = S^{-1}A,$$

hence it has the same spectrum. Since $A = S + \mathbf{e}_1 \mathbf{e}_1^T$, we have

$$S^{-1}A = I + \mathbf{ue}_1^T,$$

a rank-1 perturbation of the identity matrix, with all eigenvalues but one equal 1 (the remaining one equal $1 + u_1$).

Theorem 7

Consider the CGM. We then have the upper estimate

$$\|\mathbf{e}^{(k)}\|_{A} \leq 2\rho^{k} \|\mathbf{e}^{(0)}\|_{A}, \qquad \rho = \rho_{A} = \frac{\sqrt{\kappa(A) - 1}}{\sqrt{\kappa(A) + 1}} < 1,$$

where $\mathbf{e}^{(k)} = x^* - x^{(k)}$ and $x^{(k)}$ is the k-th output of the CGM.

Conjugate gradient - Global approach

Theorem 8 (Non-examinable)

Given $A \in \mathbb{R}^{n \times n}$, A > 0, let $\{\mathbf{d}^{(k)}\}_{k=0}^{m-1}$ be a set of the conjugate directions, i.e., $(A\mathbf{d}^{(k)}, \mathbf{d}^{(i)}) = 0$ for i < k, and consider

$$F(\mathbf{x}^{(k)}) := \|\mathbf{x}^* - \mathbf{x}^{(k)}\|_A^2 = \|\mathbf{e}^{(k)}\|_A^2$$
.

Then the value of $F(\mathbf{x}^{(m+1)})$ obtained through the CGM coincides with the minimum of $F(\mathbf{y})$ taken over all $\mathbf{y} = \mathbf{x}^{(0)} + \sum_{k=0}^{m} c_k \mathbf{d}^{(k)}$ simultaneously, namely

$$\arg\min_{c_0,...,c_m} F(\mathbf{y}) = \mathbf{x}^{(m+1)} = \mathbf{x}^{(0)} + \sum_{k=0}^m \alpha_k \mathbf{d}^{(k)}.$$

Proof of Theorem 7. As we have seen, every direction $\mathbf{d}^{(i)}$ in CGM is a linear combination of the vectors $(A^s\mathbf{r}^{(0)})_{s=0}^i$, therefore, any vector of the form $\widehat{\mathbf{x}}^{(k)} = \mathbf{x}^{(0)} + \sum_{i=0}^{k-1} a_i \mathbf{d}^{(i)}$ can be represented as

$$\hat{\mathbf{x}}^{(k)} = \mathbf{x}^{(0)} + \sum_{i=0}^{k-1} c_i A^i \mathbf{r}^{(0)}$$
 (6)

Subtracting both parts of (6) from the exact solution \mathbf{x}^* we obtain $\widehat{\mathbf{e}}^{(k)} = \mathbf{e}^{(0)} - \sum_{i=0}^{k-1} c_i A^i \mathbf{r}^{(0)}$, and since $\mathbf{r}^{(0)} = A \mathbf{e}^{(0)}$, we can express the error $\widehat{\mathbf{e}}^{(k)} = \mathbf{x}^* - \widehat{\mathbf{x}}^{(k)}$ as

$$\widehat{\mathbf{e}}^{(k)} = (I - \sum_{i=1}^{k} c_i A^i) \, \mathbf{e}^{(0)} = P_k(A) \, \mathbf{e}^{(0)}, \tag{7}$$

where P_k is a polynomial of degree $\leq k$, which satisfies $P_k(0) = 1$.

Proof. Cont. Now recall from Theorem 8 that, at the k-th stage, the CGM produces the vector $\mathbf{x}^{(k)}$ that minimizes the functional

$$F(\widehat{\mathbf{x}}^{(k)}) = \|\widehat{\mathbf{e}}^{(k)}\|_A^2 = (A\widehat{\mathbf{e}}^{(k)}, \widehat{\mathbf{e}}^{(k)})$$

over all vectors $\widehat{\mathbf{x}}^{(k)}$ of the form $\widehat{\mathbf{x}}^{(k)} = \mathbf{x}^{(0)} + \sum_{i=0}^{k-1} a_i \mathbf{d}^{(i)}$, hence over all $\widehat{\mathbf{e}}^{(k)}$ of the form (7). Expressing $\mathbf{e}^{(0)}$ as $\mathbf{e}^{(0)} = \sum \gamma_i \mathbf{w}_i$, where (\mathbf{w}_i) are orthonormal eigenvectors of A, we find from (7) that $\widehat{\mathbf{e}}^{(k)} = \sum_i \gamma_i P_k(\lambda_i) \mathbf{w}_i$, and $A\widehat{\mathbf{e}}^{(k)} = \sum_i \gamma_i P_k(\lambda_i) \lambda_i \mathbf{w}_i$, and respectively

$$\|\widehat{\mathbf{e}}^{(k)}\|_{\mathcal{A}}^2 = \sum_i [P_k(\lambda_i)]^2 \lambda_i \gamma_i^2 \leq \max_{\lambda \in \sigma(\mathcal{A})} [P_k(\lambda)]^2 \|\mathbf{e}^{(0)}\|_{\mathcal{A}}^2.$$

Hence, because of the minimization property of CGM,

$$\|\mathbf{e}^{(k)}\|_{A} = \min_{P_{k}} \|\widehat{\mathbf{e}}^{(k)}\|_{A} \le \min_{P_{k}} \max_{\lambda \in \sigma(A)} |P_{k}(\lambda)| \|\mathbf{e}^{(0)}\|_{A}.$$

Proof. Cont. Now, assume that, for the spectrum $\sigma(A)$, we know the largest and the smallest eigenvalues, or some lower and upper bounds, say, $0 < m \le \lambda \le M$. Then the following minimization problem, on the class of polynomials of degree k, arises:

$$P_k(0) = 1$$
, $\max_{x \in [m,M]} |P_k(x)| \to \min$.

This problem has a classical solution $P_k^* = T_k^*$, where T_k^* is the Chebyshev polynomial on the interval [m, M], which is obtained by dilation and translation of the standard Chebyshev polynomial T_k given on the interval [-1,1]:

$$T_k(x) = \cos k\theta, \qquad x = \cos \theta, \qquad \theta \in [0, \pi].$$

One can show that $|T_k^*(x)| \le 2\rho^k$ on the interval [m, M], hence the rate of convergence of CGM admits the following estimate:

$$\|\mathbf{e}^{(k)}\|_{A} \le 2\rho^{k} \|\mathbf{e}^{(0)}\|_{A}, \qquad \rho = \frac{\sqrt{M} - \sqrt{m}}{\sqrt{M} + \sqrt{m}} < 1, \qquad \sigma(A) \in [m, M].$$