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# NEARLY OPTIMAL CONVERGENCE RESULT FOR MULTIGRID <br> WITH AGGRESSIVE COARSENING AND POLYNOMIAL SMOOTHING 

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#### Abstract

We analyze a general multigrid method with aggressive coarsening and polynomial smoothing. We use a special polynomial smoother that originates in the context of the smoothed aggregation method. Assuming the degree of the smoothing polynomial is, on each level $k$, at least $C h_{k+1} / h_{k}$, we prove a convergence result independent of $h_{k+1} / h_{k}$. The suggested smoother is cheaper than the overlapping Schwarz method that allows to prove the same result. Moreover, unlike in the case of the overlapping Schwarz method, analysis of our smoother is completely algebraic and independent of geometry of the problem and prolongators (the geometry of coarse spaces).


Keywords: multigrid, aggressive coarsening, optimal convergence result
MSC 2010: 65F10, 65M55

## 1. Introduction

This paper is concerned with convergence of a multigrid method featuring aggressive coarsening. We analyze a general (abstract) multigrid algorithm with a special polynomial smoother that allows to prove a convergence bound independent of the relative size of the spaces on subsequent levels.

Assuming that the resolution on level $k$ can be characterized by a meshsize $h_{k}$, and employing a carefully designed polynomial smoother as a multigrid relaxation

[^0]process, we prove a convergence result independent of the ratio $h_{k+1} / h_{k}$, provided that the degree of our smoother is greater than or equal to $C h_{k+1} / h_{k}$, where $C$ is a positive constant influencing convergence, and $h_{k}$ and $h_{k+1}$ are the characteristic resolutions of finer and coarser level, respectively (throughout the paper $l$ denotes the coarsest level, and 0 the finest level). Thus, we allow the coarse space to be dramatically smaller than the preceding fine space and still obtain a multilevel convergence result not influenced by the ratio of their sizes. Here the aggressive coarsening is compensated for by a more powerful multigrid relaxation that consists of a sequence of Richardson type sweeps whose number is at least $C h_{k+1} / h_{k}, C>0$. We note that the assumption of existence of characteristic meshsize on each level is a sufficient condition, and the abstract convergence result presented in Theorem 4.1 does not depend on this assumption. The assumption will, however, allow us to verify the prerequisites of the abstract convergence result for the model problem considered. We stress that the abstract theory (Theorem 4.1) is not restricted to the quasiuniform case.

The smoother we use was originally derived from a prolongator smoother in the context of the smoothed aggregation method [9], [7], [4], and the previously proved theory [9], [8] was also limited to that context. Recent improvements of the convergence theory in [4], establishing the same convergence result as presented here, were also restricted to the smoothed aggregation method.

The regularity-free theory of [2] is known to derive no theoretical benefit from the use of more than $O(1)$ smoothing steps. Thus, until recently, the authors believed that the current result was possible only within the framework of smoothed aggregation, that is, smoothing the prolongator was deemed essential to establishing the result. The earlier works on this topic [6], [9], [7], depend crucially on this argument. In this paper, we prove a nearly optimal multilevel convergence result for this smoother used in general multigrid, and show that with a special choice of iteration parameters, $C h_{k+1} / h_{k}$ smoothing steps suffice to prove convergence independent of how aggressive the coarsening is. The near optimality of the convergence estimate is understood in the sense of the regularity-free theory [2], i.e., the convergence bound has a linear dependence on the number of levels.

Note that a similar convergence result can be proved for the overlapping Schwarz smoother. However, the relevant analysis requires verification of geometry-dependent assumptions on the overlapping subdomains which are tied to the geometrical properties of the coarse-level basis. In contrast, all assumptions on our smoother are strictly algebraic, and the analysis of the smoother is therefore independent of the geometry of the problem or the particular choice of prolongation operators (geometrical properties of coarse-levels). For our smoother, we only need the assumption that its degree is sufficiently large, which in the quasiuniform case means, greater than or
equal to $C h_{k+1} / h_{k}, C>0$. Moreover, our smoother is cheaper than the overlapping Schwarz method, while its polynomial nature allows for easy and efficient parallel implementation whenever a highly tuned parallel matrix-vector multiply subroutine is available.

The paper is organized as follows: Section 2 presents a convergence result for the multigrid method in an abstract setting. As usual, we require that the multigrid relaxation process satisfies a smoothing condition, and that the hierarchy of coarse spaces (and the associated prolongators) satisfies a weak approximation property. By the very nature of aggressive coarsening, the smoothing procedure needs to do some of the work done under normal circumstances by the coarse-grid correction process. Therefore we have a weaker approximation condition and correspondingly a stronger smoothing condition. We see the introduction of this relaxed weak approximation condition, (2.14), as the main contribution of this paper. The corresponding stronger smoothing condition, (2.16), is shown to be satisfied by our choice of the polynomial smoother.

Section 3 is devoted to the analysis of an appropriate polynomial smoother sufficient to satisfy the smoothing condition (2.16), needed to establish the convergence bound. The smoother analysis presented here essentially follows [4] and presents a minor generalization.

Section 4 summarizes the results presented in Sections 2 and 3 in the form of a final abstract convergence estimate.

We conclude the paper by considering a model example in Section 5. Here we demonstrate how the abstract result can be applied to obtain a convergence result independent of the coarsening aggressivity in the case of a model example of a geometric multigrid method with aggressive coarsening for a simple $H_{0}^{1}$-equivalent model problem discretized over a quasiuniform mesh.

## 2. Multigrid algorithm and abstract estimates

We are solving a problem

$$
A \mathbf{x}=\mathbf{f}
$$

with a symmetric positive definite (s.p.d.) matrix $A$ of order $n$. We set $A_{0}=A$ and $n_{0}=n$. We assume injective prolongators

$$
P_{k+1}^{k}: \mathbb{R}^{n_{k+1}} \rightarrow \mathbb{R}^{n_{k}}, \quad n_{k+1}<n_{k}, k=0, \ldots, l-1
$$

where $l$ is the number of levels, are given. Define a composite prolongator

$$
P_{k}^{0}=P_{1}^{0} \ldots P_{k}^{k-1}
$$

and assume that the coarse-level matrices are defined by the usual variational (Galerkin) formula

$$
\begin{equation*}
A_{k+1}=\left(P_{k+1}^{k}\right)^{T} A_{k} P_{k+1}^{k}=\left(P_{k+1}^{0}\right)^{T} A P_{k+1}^{0} \tag{2.1}
\end{equation*}
$$

To define a standard $V$-cycle multigrid, in addition to the hierarchy of matrices $\left\{A_{k}\right\}$ and prolongators $\left\{P_{k+1}^{k}\right\}$, we also need a multigrid relaxation, defined here on level $k$ as an iterative process with an error propagation operator $I-M_{k}^{-1} A_{k}$. We assume that the smoothing matrices $M_{k}$ are such that the relaxation process is an $A_{k}$-convergent iterative method, which is equivalent to $M_{k}^{T}+M_{k}-A_{k}$ being a positive definite matrix. We, in fact, assume that there is a constant $\alpha>0$, uniform with respect to $k \geqslant 0$, such that,

$$
\begin{equation*}
\mathbf{v}_{k}^{T}\left(M_{k}^{T}+M_{k}-A_{k}\right) \mathbf{v}_{k} \geqslant \alpha \mathbf{v}_{k}^{T} A_{k} \mathbf{v}_{k} \quad \text { for all } \mathbf{v}_{k} \in \mathbb{R}^{n_{k}} \tag{2.2}
\end{equation*}
$$

We denote by $\bar{M}_{k}$ the symmetrized smoother

$$
\begin{equation*}
\bar{M}_{k}=M_{k}\left(M_{k}^{T}+M_{k}-A_{k}\right)^{-1} M_{k}^{T} . \tag{2.3}
\end{equation*}
$$

It can be defined implicitly from the relation

$$
\begin{equation*}
I-\bar{M}_{k}^{-1} A_{k}=\left(I-M_{k}^{-T} A_{k}\right)\left(I-M_{k}^{-1} A_{k}\right) . \tag{2.4}
\end{equation*}
$$

Based on a given choice of $P_{k+1}^{k}, M_{k}$ (that is $A_{k}$-convergent) for $0 \leqslant k \leqslant l-1$, and $A_{k}$ obtained variationally from $A_{k-1}$ for $1 \leqslant k \leqslant l$, starting with $B_{l}=A_{l}$, for $k=l-1, \ldots, 1,0$, we recursively define a $V$-cycle preconditioner (a s.p.d. matrix) $B_{k}$ in the following standard way:

$$
I-B_{k}^{-1} A_{k}=\left(I-M_{k}^{-T} A_{k}\right)\left(I-P_{k+1}^{k} B_{k+1}^{-1}\left(P_{k+1}^{k}\right)^{T} A_{k}\right)\left(I-M_{k}^{-1} A_{k}\right) .
$$

Letting $B=B_{0}$, we are concerned in what occurs with the (upper) bound $K_{*}$ in the estimate

$$
\begin{equation*}
\mathbf{v}^{T} A \mathbf{v} \leqslant \mathbf{v}^{T} B \mathbf{v} \leqslant K_{\star} \mathbf{v}^{T} A \mathbf{v} \tag{2.5}
\end{equation*}
$$

(the lower bound holds because our algorithm is a variational multigrid). In what follows, $\|\cdot\|$ and $\langle\cdot, \cdot\rangle$ denote the Euclidean norm and the inner product in the relevant vector space. Further, for a symmetric positive definite matrix $B$, we define $\langle\cdot, \cdot\rangle_{B}=\langle B \cdot, \cdot\rangle$ and $\|\cdot\|_{B}=\langle\cdot, \cdot\rangle_{B}^{1 / 2}$.

Our analysis is based on the XZ-identity ([11]), formulated here in its matrixvector form suitable for our purposes as follows. Given multigrid smoothers defined
by $M_{j}$ such that $M_{j}^{T}+M_{j}-A_{j}$ are the s.p.d., interpolation matrices $P_{j+1}^{j}$, and the coarse matrices defined as $A_{j+1}=\left(P_{j+1}^{j}\right)^{T} A_{j} P_{j+1}^{j}$, the following XZ-identity holds (cf. [10]):

$$
\begin{align*}
& \mathbf{v}^{T} A \mathbf{v} \leqslant \mathbf{v}^{T} B \mathbf{v}  \tag{2.6}\\
&=\inf _{\left\{\mathbf{v}_{k}\right\}}\left\{\left\|\mathbf{v}_{l}\right\|_{A_{l}}^{2}+\sum_{j=0}^{l-1}\left\|M_{j}^{T} \mathbf{v}_{j}^{f}+A_{j} P_{j+1}^{j} \mathbf{v}_{j+1}\right\|_{\left(M_{j}^{T}+M_{j}-A_{j}\right)^{-1}}^{2}\right\}, \\
& \mathbf{v}_{0}=\mathbf{v}, \quad \mathbf{v}_{k}^{f} \equiv \mathbf{v}_{k}-P_{k+1}^{k} \mathbf{v}_{k+1} .
\end{align*}
$$

The infimum here is taken over the components $\left\{\mathbf{v}_{k}\right\}$ of all possible decompositions of $\mathbf{v}$ obtained as follows: Starting with $\mathbf{v}_{0}=\mathbf{v}$, for $k \geqslant 0, \mathbf{v}_{k}=\mathbf{v}_{k}^{f}+P_{k+1}^{k} \mathbf{v}_{k+1}$, i.e., choosing $\mathbf{v}_{k+1} \in \mathbb{R}^{n_{k+1}}$ arbitrary, we then let $\mathbf{v}_{k}^{f}=\mathbf{v}_{k}-P_{k+1}^{k} \mathbf{v}_{k+1}$.

We observe that applying the triangle inequality together with the trivial inequality $(a+b)^{2} \leqslant 2\left(a^{2}+b^{2}\right), a, b \in \mathbb{R}$ and using the property (2.2), results in the estimate

$$
\begin{align*}
& \sum_{j=0}^{l-1}\left\|M_{j}^{T} \mathbf{v}_{j}^{f}+A_{j} P_{j+1}^{j} \mathbf{v}_{j+1}\right\|_{\left(M_{j}^{T}+M_{j}-A_{j}\right)^{-1}}^{2}  \tag{2.7}\\
& \quad \leqslant \sum_{j=0}^{l-1} 2\left(\left\|M_{j} \mathbf{v}_{j}^{f}\right\|_{\left(M_{j}^{T}+M_{j}-A_{j}\right)^{-1}}^{2}+\left\|A_{j} P_{j+1}^{j} \mathbf{v}_{j+1}\right\|_{\left.\left(M_{j}^{T}+M_{j}-A_{j}\right)^{-1}\right)}^{2}\right. \\
& \quad=2 \sum_{j=0}^{l-1}\left\|M_{j} \mathbf{v}_{j}^{f}\right\|_{\left(M_{j}^{T}+M_{j}-A_{j}\right)^{-1}}^{2}+2 \sum_{j=0}^{l-1}\left\|A_{j} P_{j+1}^{j} \mathbf{v}_{j+1}\right\|_{\left(M_{j}^{T}+M_{j}-A_{j}\right)^{-1}}^{2} \\
& \quad \leqslant 2 \sum_{j=0}^{l-1}\left\|\mathbf{v}_{j}^{f}\right\|_{M_{j}}^{2}+\frac{2}{\alpha} \sum_{j=0}^{l-1}\left\|A_{j} P_{j+1}^{j} \mathbf{v}_{j+1}\right\|_{A_{j}^{-1}}^{2} \\
& \quad=2 \sum_{j=0}^{l-1}\left\|\mathbf{v}_{j}^{f}\right\|_{\bar{M}_{j}}^{2}+\frac{2}{\alpha} \sum_{j=0}^{l-1}\left\|P_{j+1}^{j} \mathbf{v}_{j+1}\right\|_{A_{j}}^{2} \\
& \quad=2 \sum_{j=0}^{l-1}\left\|\mathbf{v}_{j}^{f}\right\|_{M_{j}}^{2}+\frac{2}{\alpha} \sum_{j=0}^{l-1}\left\|\mathbf{v}_{j+1}\right\|_{A_{j+1}}^{2} \\
& \quad=2 \sum_{j=0}^{l-1}\left\|\mathbf{v}_{j}^{f}\right\|_{M_{j}}^{2}+\frac{2}{\alpha} \sum_{j=1}^{l}\left\|\mathbf{v}_{j}\right\|_{A_{j}}^{2} .
\end{align*}
$$

From here, we see that in order to bound the relative condition number of the $V$-cycle preconditioner $B$ with respect to $A$ based on estimates (2.6) and (2.7), it is sufficient
to bound the expressions below in terms of $\|\mathbf{v}\|_{A}^{2}$ for some particular choice of $\left\{\mathbf{v}_{k}\right\}$ :

$$
\begin{align*}
& \sum_{k=0}^{l-1}\left\|\mathbf{v}_{k}^{f}\right\|_{M_{k}}^{2} \leqslant C_{1}\|\mathbf{v}\|_{A}^{2},  \tag{2.8}\\
& \sum_{k=1}^{l}\left\|\mathbf{v}_{k}\right\|_{A_{k}}^{2} \leqslant C_{2}\|\mathbf{v}\|_{A}^{2},
\end{align*}
$$

and

$$
\begin{equation*}
\left\|\mathbf{v}_{l}\right\|_{A_{l}}^{2} \leqslant C_{3}\|\mathbf{v}\|_{A}^{2} . \tag{2.10}
\end{equation*}
$$

Indeed, the estimates (2.6) and (2.7) give

$$
\begin{align*}
\mathbf{v}^{T} A \mathbf{v} \leqslant \mathbf{v}^{T} B \mathbf{v} & \leqslant\left\|\mathbf{v}_{l}\right\|_{A_{l}}^{2}+2 \sum_{j=0}^{l-1}\left\|\mathbf{v}_{j}^{f}\right\|_{M_{j}}^{2}+\frac{2}{\alpha} \sum_{j=1}^{l}\left\|\mathbf{v}_{j}\right\|_{A_{j}}^{2}  \tag{2.11}\\
& \leqslant\left(C_{3}+2 C_{1}+\frac{2}{\alpha} C_{2}\right) \mathbf{v}^{T} A \mathbf{v}
\end{align*}
$$

Note that (2.10) follows from (2.9) with $C_{3}=C_{2}$.
We define coarse-spaces $V_{k}$ and associated norms $\|\cdot\|_{k}$ by

$$
\begin{align*}
V_{k} & =\operatorname{Rng}\left(P_{k}^{0}\right),  \tag{2.12}\\
\|\cdot\|_{k}: \quad P_{k}^{0} \mathbf{x} \mapsto\|\mathbf{x}\| & \equiv \sqrt{\mathbf{x}^{T} \mathbf{x}}, \quad k=0, \ldots, l\left(P_{0}^{0}=I\right) .
\end{align*}
$$

Further, we define

$$
\begin{equation*}
\lambda_{k, j}=\sup _{\mathbf{x} \in \mathbb{R}^{n_{k}} \backslash\{\mathbf{0}\}} \frac{\left\langle A P_{k}^{0} \mathbf{x}, P_{k}^{0} \mathbf{x}\right\rangle}{\left\|P_{k}^{0} \mathbf{x}\right\|_{j}^{2}}, \quad k=0, \ldots, l, 0 \leqslant j \leqslant k . \tag{2.13}
\end{equation*}
$$

Note that $\lambda_{k, j} \leqslant \varrho\left(A_{j}\right)$ and $\lambda_{k, k}=\varrho\left(A_{k}\right)$.
Remark 2.1. Definition (2.13) allows the following interpretation: The spectral bound

$$
\varrho\left(A_{k}\right)=\sup _{\mathbf{x} \in \mathbb{R}^{n} k \backslash\{\mathbf{0}\}} \frac{\left\langle A_{k} \mathbf{x}, \mathbf{x}\right\rangle}{\|\mathbf{x}\|^{2}}=\sup _{\mathbf{x} \in \mathbb{R}^{n}{ }^{n} \backslash\{\mathbf{0}\}} \frac{\left\langle A P_{k}^{0} \mathbf{x}, P_{k}^{0} \mathbf{x}\right\rangle}{\left\|P_{k}^{0} \mathbf{x}\right\|_{k}^{2}}
$$

indicates the smoothness of the space $V_{k}$ with respect to the norm $\|\cdot\|_{k}$. The quantity

$$
\lambda_{k, j}=\sup _{\mathbf{x} \in \mathbb{R}^{n_{k}} \backslash\{\mathbf{0}\}} \frac{\left\langle A P_{k}^{0} \mathbf{x}, P_{k}^{0} \mathbf{x}\right\rangle}{\left\|P_{k}^{0} \mathbf{x}\right\|_{j}^{2}}=\sup _{\mathbf{x} \in \mathbb{R}^{n_{k}} \backslash\{\mathbf{0}\}} \frac{\left\langle A P_{k}^{0} \mathbf{x}, P_{k}^{0} \mathbf{x}\right\rangle}{\left\|P_{j}^{0} P_{k}^{j} \mathbf{x}\right\|_{j}^{2}}=\sup _{\mathbf{x} \in \mathbb{R}^{n_{k}} \backslash\{\mathbf{0}\}} \frac{\left\langle A P_{k}^{0} \mathbf{x}, P_{k}^{0} \mathbf{x}\right\rangle}{\left\|P_{k}^{j} \mathbf{x}\right\|^{2}},
$$

$j<k$, indicates the smoothness of the space $V_{k}$ with respect to the finer space norm $\|\cdot\|_{j}$.

We now formulate our abstract convergence estimate,

Theorem 2.1. Let $\bar{\lambda}_{k+1, k} \geqslant \lambda_{k+1, k}, k=0, \ldots, l-1$ be upper bounds. We assume the existence of linear mappings $Q_{k}: V_{0} \rightarrow V_{k}, Q_{0}=I$, satisfying

$$
\begin{equation*}
\left\|\left(Q_{k}-Q_{k+1}\right) \mathbf{v}\right\|_{k} \leqslant \frac{C_{a}}{\sqrt{\bar{\lambda}_{k+1, k}}}\|\mathbf{v}\|_{A} \quad \forall \mathbf{v} \in V_{0}, k=0, \ldots, l-1 \tag{2.14}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|Q_{k}\right\|_{A} \leqslant C_{s}, \quad k=0, \ldots, l . \tag{2.15}
\end{equation*}
$$

Further, we assume that our smoothers, $M_{k}$, satisfy (2.2) and the symmetrized smoothers $\bar{M}_{k}$ satisfy

$$
\begin{equation*}
\|\mathbf{v}\|_{M_{k}}^{2} \leqslant \beta\left(\bar{\lambda}_{k+1, k}\|\mathbf{v}\|^{2}+\|\mathbf{v}\|_{A_{k}}^{2}\right) \quad \forall \mathbf{v} \in \mathbb{R}^{n_{k}}, k=0, \ldots, l-1 . \tag{2.16}
\end{equation*}
$$

Then the resulting multigrid operator $B$ is nearly spectrally equivalent to $A$, more precisely,

$$
\begin{equation*}
\mathbf{v}^{T} A \mathbf{v} \leqslant \mathbf{v}^{T} B \mathbf{v} \leqslant\left[C_{s}^{2}+2 l\left(\beta\left(C_{a}^{2}+4 C_{s}^{2}\right)+\frac{1}{\alpha} C_{s}^{2}\right)\right] \mathbf{v}^{T} A \mathbf{v} \quad \forall \mathbf{v} \in \mathbb{R}^{n_{0}} \tag{2.17}
\end{equation*}
$$

Remark 2.2. The difference from the results previously obtained based on the theory in [2] is in our use of the weak approximation condition (2.14). The original theory relied instead on the condition

$$
\begin{equation*}
\left\|\left(Q_{k}-Q_{k+1}\right) \mathbf{v}\right\|_{k} \leqslant \frac{C_{a}}{\sqrt{\varrho\left(A_{k}\right)}}\|\mathbf{v}\|_{A} \tag{2.18}
\end{equation*}
$$

and the approximation properties of the space $V_{k+1}$ were thus measured against the smoothness of the space $V_{k}$ (because of $\varrho\left(A_{k}\right)$ ). In typical applications, the approximation on the left-hand side of (2.18) is guided by $h_{k+1}$, while the spectral bound of $A_{k}$ and the scaling of the $\|\cdot\|_{k}$-norm are guided by $h_{k}$. To prove (2.18), the ratio $h_{k+1} / h_{k}$ has to be bounded, and the resolutions of spaces $V_{k}$ and $V_{k+1}$ have to be comparable.

In our case, the approximation properties of the space $V_{k+1}$ are measured against (the upper bound of) $\lambda_{k+1, k} \equiv \sup _{\mathbf{x} \in \operatorname{Rng}\left(P_{k+1}^{k}\right) \backslash\{\mathbf{0}\}}\left\langle A_{k} \mathbf{x}, \mathbf{x}\right\rangle /\|\mathbf{x}\|^{2} \leqslant \varrho\left(A_{k}\right)$, that is, against the smoothness of the space $V_{k+1}$ (measured with respect to the norm $\|\cdot\|_{k}$ used on the left-hand side of (2.14)), and therefore the resolutions of the spaces $V_{k}$ and $V_{k+1}$ do not have to be comparable. The current estimate thus allows us to prove a convergence result independent of the coarsening ratio. The cost of the uniform convergence result, when the coarsening ratio becomes large $\left(\lambda_{k+1, k} \ll \varrho\left(A_{k}\right)\right)$, is in the increasing demand on the smoother that arises through the smoothing condition (2.16).

Proof. We write linear mappings $Q_{k}: V_{0} \rightarrow V_{k} \equiv \operatorname{Rng}\left(P_{k}^{0}\right), k=0, \ldots, l$, in the form

$$
Q_{k}=P_{k}^{0} \widetilde{Q}_{k}, \quad \widetilde{Q}_{k}: V_{0}=\mathbb{R}^{n_{0}} \rightarrow \mathbb{R}^{n_{k}}
$$

In the XZ-identity (2.6), we choose

$$
\mathbf{v}_{k}=\widetilde{Q}_{k} \mathbf{v}, \quad k=1, \ldots, l .
$$

Therefore (see (2.6)),

$$
\mathbf{v}_{k}^{f}=\mathbf{v}_{k}-P_{k+1}^{k} \mathbf{v}_{k+1}=\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}
$$

Thus, to prove our theorem means to verify the inequalities (2.8), (2.9), and (2.10) for the above particular decomposition of $\mathbf{v}$.

To prove (2.8), we estimate using the assumptions (2.16), (2.14) and (2.15), the definition (2.12) of $\|\cdot\|_{k}, Q_{k}=P_{k}^{0} \widetilde{Q}_{k}$ for $k=0, \ldots, l$ and the triangle inequality:

$$
\begin{align*}
\left\|\mathbf{v}_{k}^{f}\right\|_{\bar{M}_{k}}^{2} & =\left\|\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{\bar{M}_{k}}^{2}  \tag{2.19}\\
& \leqslant \beta\left(\bar{\lambda}_{k+1, k}\left\|\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|^{2}+\left\|\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{A_{k}}^{2}\right) \\
& =\beta\left(\bar{\lambda}_{k+1, k}\left\|P_{k}^{0}\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{k}^{2}+\left\|P_{k}^{0}\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{A}^{2}\right) \\
& =\beta\left(\bar{\lambda}_{k+1, k}\left\|\left(Q_{k}-Q_{k+1}\right) \mathbf{v}\right\|_{k}^{2}+\left\|\left(Q_{k}-Q_{k+1}\right) \mathbf{v}\right\|_{A}^{2}\right) \\
& \leqslant \beta\left(C_{a}^{2}\|\mathbf{v}\|_{A}^{2}+2\left(\left\|Q_{k} \mathbf{v}\right\|_{A}^{2}+\left\|Q_{k+1} \mathbf{v}\right\|_{A}^{2}\right)\right) \\
& \leqslant \beta\left(C_{a}^{2}+4 C_{s}^{2}\right)\|\mathbf{v}\|_{A}^{2} .
\end{align*}
$$

Thus,

$$
\sum_{k=0}^{l-1}\left\|\mathbf{v}_{k}^{f}\right\|_{M_{k}}^{2} \leqslant l \beta\left(C_{a}^{2}+4 C_{s}^{2}\right)\|\mathbf{v}\|_{A}^{2},
$$

proving (2.8) with a constant

$$
C_{1}=l \beta\left(C_{a}^{2}+4 C_{s}^{2}\right) .
$$

To prove (2.9) and (2.10), we realize that

$$
\left\|\mathbf{v}_{k}\right\|_{A_{k}}^{2}=\left\|\widetilde{Q}_{k} \mathbf{v}\right\|_{A_{k}}^{2}=\left\|Q_{k} \mathbf{v}\right\|_{A}^{2} \leqslant C_{s}^{2}\|\mathbf{v}\|_{A}^{2},
$$

hence (2.9) immediately follows with a constant

$$
C_{2}=C_{s}^{2} l
$$

and (2.10) with a constant

$$
C_{3}=C_{s}^{2}
$$

The estimate (2.17) now follows by (2.11).

## 3. Polynomial smoother

In this section, we investigate a polynomial smoother with the error propagation operator

$$
\begin{equation*}
I-M_{k}^{-T} A_{k}=I-M_{k}^{-1} A_{k}=S_{k}^{\gamma}\left(I-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2} A_{k}\right) \tag{3.1}
\end{equation*}
$$

where $S_{k}$ is a polynomial in $A_{k}$ such that $\varrho\left(S_{k}\right) \leqslant 1, \bar{\lambda}_{S_{k}^{2} A_{k}} \geqslant \varrho\left(S_{k}^{2} A_{k}\right)$ and $\gamma$ is a positive integer. The particular cases of interest are $\gamma=1$ and $\gamma=2$.

From (2.4) and the fact that the error propagation operator corresponding to the symmetrized smoother $\bar{M}_{k}$ is

$$
I-\bar{M}_{k}^{-1} A_{k}=\left(I-M_{k}^{-T} A_{k}\right)\left(I-M_{k}^{-1} A_{k}\right)=S_{k}^{2 \gamma}\left(I-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2} A_{k}\right)^{2}
$$

it follows that the corresponding symmetrized smoother $\bar{M}_{k}$ is given by

$$
\begin{equation*}
\bar{M}_{k}^{-1}=A_{k}^{-1}\left[I-\left(I-\frac{1}{\overline{\lambda_{S}^{2} A_{k}}} S_{k}^{2} A_{k}\right)^{2} S_{k}^{2 \gamma}\right] \tag{3.2}
\end{equation*}
$$

Lemma 3.1. We assume that $S_{k}$ is a polynomial in $A_{k}$ such that $\varrho\left(S_{k}\right) \leqslant 1$, $\bar{\lambda}_{S_{k}^{2} A_{k}} \geqslant \varrho\left(S_{k}^{2} A_{k}\right)$, and $\gamma$ is a positive integer. Let $\left\{v_{i}\right\}$ be the eigenvectors of $A_{k}$ and $\lambda_{i}\left(S_{k}\right)$ the corresponding eigenvalues of $S_{k}$. For a given parameter, $q \in(0,1)$, define

$$
U_{1}=\left\{\operatorname{span}\left\{v_{i}\right\}:\left|\lambda_{i}\left(S_{k}\right)\right| \leqslant q\right\} \quad \text { and } \quad U_{2}=\left\{\operatorname{span}\left\{v_{i}\right\}:\left|\lambda_{i}\left(S_{k}\right)\right|>q\right\}
$$

Then the symmetrized smoother $\bar{M}_{k}$ in (3.2) is positive definite and satisfies

$$
\|\mathbf{x}\|_{\bar{M}_{k}}^{2} \leqslant \frac{1}{1-q^{2 \gamma}}\|\mathbf{x}\|_{A_{k}}^{2} \forall \mathbf{x} \in U_{1}, \quad\|\mathbf{x}\|_{\bar{M}_{k}}^{2} \leqslant \frac{\bar{\lambda}_{S_{k}^{2} A_{k}}}{q^{2}}\|\mathbf{x}\|^{2} \forall \mathbf{x} \in U_{2}
$$

and

$$
\begin{equation*}
\|\mathbf{x}\|_{\bar{M}_{k}}^{2} \leqslant \frac{1}{1-q^{2 \gamma}}\|\mathbf{x}\|_{A_{k}}^{2}+\frac{\bar{\lambda}_{S_{k}^{2} A_{k}}}{q^{2}}\|\mathbf{x}\|^{2} \quad \forall \mathbf{x} \in \mathbb{R}^{n_{k}}, q \in(0,1) \tag{3.3}
\end{equation*}
$$

Remark 3.1. Our goal is to satisfy the smoothing condition (2.16). Therefore, in view of (3.3), the property the smoother (3.1) needs to satisfy is

$$
\begin{equation*}
\bar{\lambda}_{S_{k}^{2} A_{k}} \leqslant C \bar{\lambda}_{k+1, k} \tag{3.4}
\end{equation*}
$$

(with $\bar{\lambda}_{k+1, k} \ll \varrho\left(A_{k}\right)$ for aggressive coarsening). Indeed, from (3.4) and (3.3), it follows that

$$
\begin{aligned}
\|\mathbf{x}\|_{\bar{M}_{k}}^{2} & \leqslant \max \left\{\frac{1}{1-q^{2 \gamma}}, \frac{1}{q^{2}}\right\}\left(\|\mathbf{x}\|_{A_{k}}^{2}+\bar{\lambda}_{S_{k}^{2} A_{k}}\|\mathbf{x}\|^{2}\right) \\
& \leqslant \max \left\{\frac{1}{1-q^{2 \gamma}}, \frac{1}{q^{2}}\right\} \cdot \max \{1, C\}\left(\|\mathbf{x}\|_{A_{k}}^{2}+\bar{\lambda}_{k+1, k}\|\mathbf{x}\|^{2}\right) \quad \forall \mathbf{x} \in \mathbb{R}^{n_{k}}
\end{aligned}
$$

Here $q \in(0,1)$ is a parameter we choose. Thus, (2.16) follows from (3.4) and (3.3) with

$$
\begin{equation*}
\beta=\min _{q \in(0,1)} \max \left\{\frac{1}{1-q^{2 \gamma}}, \frac{1}{q^{2}}\right\} \cdot \max \{1, C\} . \tag{3.5}
\end{equation*}
$$

The role of the smoothing polynomial, $S_{k}=p\left(A_{k}\right)$, is therefore to minimize $\varrho\left(S_{k}^{2} A_{k}\right)$ (to attain the same order of magnitude as $\bar{\lambda}_{k+1, k}$ ), subject to the constraint that $S_{k}$ is an error propagation operator of an $A_{k}$-non-divergent smoother, that is, $p(0)=1$ and $\varrho\left(S_{k}\right) \leqslant 1$. Let $\bar{\lambda}_{k} \geqslant \varrho\left(A_{k}\right)$ be an available upper bound. The polynomial $p$ of a given degree, $N_{k}$, satisfying the above constraints and minimizing the right-hand side of the inequality

$$
\varrho\left(S_{k}^{2} A_{k}\right)=\varrho\left(p^{2}\left(A_{k}\right) A_{k}\right)=\max _{t \in \sigma\left(A_{k}\right)} p^{2}(t) t \leqslant \max _{t \in\left[0, \lambda_{k}\right]} p^{2}(t) t
$$

will be given in Lemma 3.2.
Remark 3.2. For $\gamma=1$, using the minimizer $\hat{q}=1 / \sqrt{2}$, we have (see (3.5))

$$
\min _{q \in(0,1)} \max \left\{\frac{1}{1-q^{2 \gamma}}, \frac{1}{q^{2}}\right\}=2
$$

Similarly, for $\gamma=2$, using the minimizer $\hat{q}=\sqrt{\frac{1}{2}(-1+\sqrt{5})}$, we get

$$
\min _{q \in(0,1)} \max \left\{\frac{1}{1-q^{2 \gamma}}, \frac{1}{q^{2}}\right\}=\frac{2}{-1+\sqrt{5}} \doteq 1.618034
$$

Proof. The proof given here is a generalization of the one given in [4].
Recall that both $S_{k}$ and $I-\bar{\lambda}_{S_{k}^{2} A_{k}}^{-1} S_{k}^{2} A_{k}$ are polynomials in $A_{k}$, hence all these matrices have common eigenvectors, mutually commute and $U_{1}$ and $U_{2}$ are their common invariant subspaces. Further, $\varrho\left(S_{k}\right) \leqslant 1$ and $\varrho\left(I-\bar{\lambda}_{S_{k}^{2} A_{k}}^{-1} S_{k}^{2} A_{k}\right) \leqslant 1$.

To prove

$$
\begin{equation*}
\left\langle\bar{M}_{k} \mathbf{x}, \mathbf{x}\right\rangle \leqslant \frac{1}{1-q^{2 \gamma}}\left\langle A_{k} \mathbf{x}, \mathbf{x}\right\rangle \text { on } U_{1}, \tag{3.6}
\end{equation*}
$$

$\left(\langle\cdot, \cdot\rangle\right.$ denotes the Euclidean inner product in $\left.\mathbb{R}^{n_{k}}\right)$, we use (3.2) and estimate for $\mathrm{x} \in U_{1}$ :

$$
\begin{aligned}
\left\langle\bar{M}_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle & =\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle-\left\langle A_{k}^{-1}\left(I-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2} A_{k}\right)^{2} S_{k}^{2 \gamma} \mathbf{x}, \mathbf{x}\right\rangle \\
& \geqslant\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle-\left\langle A_{k}^{-1} S_{k}^{2 \gamma} \mathbf{x}, \mathbf{x}\right\rangle \\
& \geqslant\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle-q^{2 \gamma}\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle=\left(1-q^{2 \gamma}\right)\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle
\end{aligned}
$$

Since $U_{1}$ is an invariant subspace of both $\bar{M}_{k}$ and $A_{k}$, both $\bar{M}_{k}^{-1}$ and $\bar{M}_{k}$ are symmetric, positive definite on $U_{1}$ and the statement (3.6) follows.

To prove

$$
\begin{equation*}
\left\langle\bar{M}_{k} \mathbf{x}, \mathbf{x}\right\rangle \leqslant \frac{\bar{\lambda}_{S_{k}^{2} A_{k}}}{q^{2}}\|\mathbf{x}\|^{2} \quad \text { on } U_{2} \tag{3.7}
\end{equation*}
$$

we estimate for $\mathbf{x} \in U_{2}$ :

$$
\begin{aligned}
\left\langle\bar{M}_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle & =\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle-\left\langle A_{k}^{-1}\left(I-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2} A_{k}\right)^{2} S_{k}^{2 \gamma} \mathbf{x}, \mathbf{x}\right\rangle \\
& \geqslant\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle-\left\langle A_{k}^{-1}\left(I-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2} A_{k}\right) \mathbf{x}, \mathbf{x}\right\rangle \\
& =\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}}\left\langle S_{k}^{2} \mathbf{x}, \mathbf{x}\right\rangle \\
& >\frac{q^{2}}{\bar{\lambda}_{S_{k}^{2} A_{k}}}\|\mathbf{x}\|^{2}
\end{aligned}
$$

Since $U_{2}$ is an invariant subspace of $\bar{M}_{k}$, both $\bar{M}_{k}^{-1}$ and $\bar{M}_{k}$ are symmetric, positive definite on $U_{2}$ and the statement (3.7) follows.

Let us consider the decomposition of $\mathbf{x} \in \mathbb{R}^{n_{k}} \backslash\{\mathbf{0}\}$,

$$
\mathbf{x}=\mathbf{x}_{1}+\mathbf{x}_{2}, \quad \mathbf{x}_{1} \in U_{1}, \mathbf{x}_{2} \in U_{2}
$$

From the definition of the spaces $U_{1}, U_{2}$ it follows that the spaces $U_{1}$ and $U_{2}$ are orthogonal, that is,

$$
\left\langle\mathbf{x}_{1}, \mathbf{x}_{2}\right\rangle=0
$$

Since $U_{1}$ and $U_{2}$ are invariant subspaces of both $A_{k}$ and $\bar{M}_{k}$, it also follows that

$$
\left\langle A_{k} \mathbf{x}_{1}, \mathbf{x}_{2}\right\rangle=\left\langle\bar{M}_{k} \mathbf{x}_{1}, \mathbf{x}_{2}\right\rangle=0
$$

Therefore, since $\bar{M}_{k}$ is symmetric, positive definite on both $U_{1}$ and $U_{2}$, it follows that

$$
\left\langle\bar{M}_{k} \mathbf{x}, \mathbf{x}\right\rangle=\left\langle\bar{M}_{k} \mathbf{x}_{1}, \mathbf{x}_{1}\right\rangle+\left\langle\bar{M}_{k} \mathbf{x}_{2}, \mathbf{x}_{2}\right\rangle>0
$$

hence $\bar{M}_{k}$ is symmetric, positive definite on $\mathbb{R}^{n_{k}}$. Thus, the spaces $U_{1}$ and $U_{2}$ form a decomposition of $\mathbb{R}^{n_{k}}$ that is orthogonal with respect to the norms $\|\cdot\|,\|\cdot\|_{A_{k}}$ and $\|\cdot\|_{\bar{M}_{k}}$. Then (3.6) and (3.7) give

$$
\begin{aligned}
\|\mathbf{x}\|_{\bar{M}_{k}}^{2} & =\left\|\mathbf{x}_{1}\right\|_{\bar{M}_{k}}^{2}+\left\|\mathbf{x}_{2}\right\|_{\bar{M}_{k}}^{2} \\
& \leqslant \frac{1}{1-q^{2 \gamma}}\left\|\mathbf{x}_{1}\right\|_{A_{k}}^{2}+\frac{\bar{\lambda}_{S_{k}^{2} A_{k}}}{q^{2}}\left\|\mathbf{x}_{2}\right\|^{2} \leqslant \frac{1}{1-q^{2 \gamma}}\|\mathbf{x}\|_{A_{k}}^{2}+\frac{\bar{\lambda}_{S_{k}^{2} A_{k}}}{q^{2}}\|\mathbf{x}\|^{2}
\end{aligned}
$$

proving (3.3).
While the validity of property (2.16) is addressed by Lemma 3.1, we still need to verify that inequality (2.2) is satisfied for our choice of the smoother. To this end, the smoother $S_{k}$ is introduced in the next lemma.

Lemma 3.2. For any $\lambda>0$ and integer $N>0$ there is a unique polynomial $p_{\lambda, N}$ of degree $N$ such that

$$
\max _{0 \leqslant t \leqslant \lambda} p_{\lambda, N}^{2}(t) t
$$

is minimal under the constraint $p_{\lambda, N}(0)=1$. The polynomial $p$ is given by

$$
\begin{equation*}
p_{\lambda, N}(t)=\left(1-\frac{t}{r_{1}}\right) \ldots\left(1-\frac{t}{r_{N}}\right), \quad r_{k}=\frac{\lambda}{2}\left(1-\cos \left(\frac{2 k \pi}{2 N+1}\right)\right), \tag{3.8}
\end{equation*}
$$

$k=1, \ldots, N$. The polynomial $p_{\lambda, N}$ satisfies

$$
\begin{equation*}
\max _{0 \leqslant t \leqslant \lambda} p_{\lambda, N}^{2}(t) t=\frac{\lambda}{(2 N+1)^{2}} \tag{3.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\max _{0 \leqslant t \leqslant \lambda}\left|p_{\lambda, N}(t)\right|=1 \tag{3.10}
\end{equation*}
$$

The polynomial $p_{\lambda, N}$ is the transformed Chebyshev polynomial

$$
p_{\lambda, N}(t)=(-1)^{N} \frac{1}{2 N+1} \frac{\sqrt{\lambda}}{\sqrt{t}} T_{2 N+1}\left(\frac{\sqrt{t}}{\sqrt{\lambda}}\right),
$$

where $T_{k}$ is a Chebyshev polynomial of degree $k$, that is $T_{0}(t)=1, T_{1}(t)=t$, and $T_{k+1}(t)=2 t T_{k}(t)-T_{k-1}(t)$ for $k \geqslant 1$.

Proof. Proof of the lemma in this form can be found in [3]. The analysis of Chebyshev polynomials can be found in [1], see also [10].

Let $\bar{\lambda}_{k}$ be an available upper bound of $\varrho\left(A_{k}\right)$ and let the integer $N_{k}$ be a given degree of the smoothing polynomial. We choose

$$
\begin{equation*}
S_{k}=p_{\bar{\lambda}_{k}, N_{k}}\left(A_{k}\right) \tag{3.11}
\end{equation*}
$$

where $p_{\lambda, N}$ is given by (3.8). Further, we set

$$
\bar{\lambda}_{S_{k}^{2} A_{k}}=\frac{\bar{\lambda}_{k}}{\left(2 N_{k}+1\right)^{2}}
$$

Then, by Lemma 3.2 and the spectral mapping theorem, we have

$$
\begin{align*}
\varrho\left(S_{k}^{2} A_{k}\right) & =\max _{t \in \sigma\left(A_{k}\right)} p_{\lambda_{k}, N_{k}}^{2}(t) t \leqslant \max _{t \in\left[0, \bar{\lambda}_{k}\right]} p_{\lambda_{k}, N_{k}}^{2}(t) t  \tag{3.12}\\
& =\bar{\lambda}_{S_{k}^{2} A_{k}} \equiv \frac{\bar{\lambda}_{k}}{\left(2 \operatorname{deg}\left(S_{k}\right)+1\right)^{2}}, \quad \varrho\left(S_{k}\right) \leqslant 1
\end{align*}
$$

Lemma 3.3. For the smoother (3.1), with $\gamma=1$ and $S_{k}$ given by (3.11) and (3.8), the inequality (2.2) holds with

$$
\alpha=\frac{\delta_{0}}{2-\delta_{0}}, \quad \delta_{0}=1-\frac{2}{3 \sqrt{3}} \in(0,1) .
$$

Further, for $\gamma>0$ that is even, and $S_{k}$ being a polynomial in $A_{k}$ satisfying $\varrho\left(S_{k}\right) \leqslant 1$, the inequality (2.2) holds with $\alpha=1$. (That is, for even $\gamma>0$, we do not have to assume that $S_{k}$ is given by (3.11) and (3.8), we only need $S_{k}$ to be a polynomial in $A_{k}$ such that $\varrho\left(S_{k}\right) \leqslant 1$.)

Proof. For the proof in the case of $\gamma=1$ and $S_{k}$ given by (3.11) and (3.8), see [4], Lemma 6.2 and Proposition 7.3.

For an even $\gamma$ and $S_{k}$ being a polynomial in $A_{k}$ satisfying $\varrho\left(S_{k}\right) \leqslant 1$, we have

$$
\left\langle M_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle=\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle-\left\langle A_{k}^{-1}\left(I-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2} A_{k}\right) S_{k}^{\gamma} \mathbf{x}, \mathbf{x}\right\rangle \leqslant\left\langle A_{k}^{-1} \mathbf{x}, \mathbf{x}\right\rangle
$$

Hence, $M_{k} \geqslant A_{k}$, and therefore assumption (2.2) on the smoother $M_{k}$ holds trivially with $\alpha=1$.

Remark 3.3. The most natural way to implement the action of (3.1) for a given vector $\mathbf{x}$ is the following: To perform the iteration with the linear part $S_{k}$ given by (3.11) and (3.8), we do for $i=1, \ldots, N_{k}=\operatorname{deg}\left(S_{k}\right)$,

$$
\mathbf{x} \leftarrow\left(I-\alpha_{i} A_{k}\right) \mathbf{x}+\alpha_{i} \mathbf{f}, \quad \alpha_{i}=\left(\frac{\bar{\lambda}_{k}}{2}\left(1-\cos \left(\frac{2 i \pi}{2 N_{k}+1}\right)\right)\right)^{-1}, \quad \bar{\lambda}_{k} \geqslant \varrho\left(A_{k}\right)
$$

To perform the iteration with the error propagation operator $I-\bar{\lambda}_{S_{k}^{2} A_{k}}^{-1} S_{k}^{2} A_{k}$, we do

$$
\mathbf{x} \leftarrow \mathbf{x}-\frac{1}{\bar{\lambda}_{S_{k}^{2} A_{k}}} S_{k}^{2}\left(A_{k} \mathbf{x}-\mathbf{f}\right),
$$

where the action of $S_{k}$ is evaluated as the product

$$
S_{k} \mathbf{x}=\left(I-\alpha_{1} A_{k}\right) \ldots\left(I-\alpha_{N_{k}} A_{k}\right) \mathbf{x} .
$$

## 4. The final abstract result

In this section, we summarize the results proved in Sections 2 and 3 in the form of a theorem.

Theorem 4.1. Let $\bar{\lambda}_{k+1, k} \geqslant \lambda_{k+1, k}(k=0, \ldots, l-1)$ and $\bar{\lambda}_{k} \geqslant \varrho\left(A_{k}\right)(k=$ $0, \ldots, l)$ be upper bounds. We assume the existence of linear mappings (see (2.12)) $Q_{k}: V_{0} \rightarrow V_{k}, k=0, \ldots, l, Q_{0}=I$, satisfying (2.14) and (2.15) with positive constants $C_{a}$ and $C_{s}$, independent of the level. Further, we assume that the linear part of both the pre- and post-smoother is given by (3.1) with $S_{k}=p_{\bar{\lambda}_{k}, N_{k}}\left(A_{k}\right)$, where the polynomial $p_{\lambda, N}$ is given by (3.8) and its degree, $N_{k}$, satisfies

$$
\begin{equation*}
N_{k} \geqslant C_{\operatorname{deg}} \sqrt{\frac{\bar{\lambda}_{k}}{\overline{\lambda_{k+1, k}}}}, \quad k=0, \ldots, l-1 \tag{4.1}
\end{equation*}
$$

with a constant $C_{\text {deg }}>0$ independent of the level. We assume that $\gamma$ in (3.1) is either even, or $\gamma=1$. Then (2.17) is satisfied; that is,

$$
\mathbf{v}^{T} A \mathbf{v} \leqslant \mathbf{v}^{T} B \mathbf{v} \leqslant\left[C_{s}^{2}+2 l\left(\beta\left(C_{a}^{2}+4 C_{s}^{2}\right)+\frac{1}{\alpha} C_{s}^{2}\right)\right] \mathbf{v}^{T} A \mathbf{v} \quad \forall \mathbf{v} \in \mathbb{R}^{n_{0}}
$$

holds with

$$
\alpha= \begin{cases}1 & \text { for } \gamma \text { even, }  \tag{4.2}\\ \frac{\delta_{0}}{2-\delta_{0}}, \delta_{0}=1-\frac{2}{3 \sqrt{3}} & \text { for } \gamma=1\end{cases}
$$

and

$$
\begin{equation*}
\beta=\min _{q \in(0,1)} \max \left\{\frac{1}{1-q^{2 \gamma}}, \frac{1}{q^{2}}\right\} \cdot \max \left\{1, \frac{1}{4 C_{\mathrm{deg}}^{2}}\right\} . \tag{4.3}
\end{equation*}
$$

Proof. Statement (2.17) follows from Theorem 2.1 under assumptions (2.2) and (2.16) (inequalities (2.14) and (2.15) are assumptions of this theorem).

Assumption (2.2), with $\alpha$ given by (4.2), has been verified by Lemma 3.3.
According to Remark 3.1, (2.16) holds under assumption (3.4). Using Lemma 3.2, we estimate

$$
\varrho\left(S_{k}^{2} A_{k}\right)=\max _{t \in \sigma\left(A_{k}\right)} p_{\lambda_{k}, N_{k}}^{2}(t) t \leqslant \bar{\lambda}_{S_{k}^{2} A_{k}} \equiv \max _{t \in\left[0, \lambda_{k}\right]} p_{\bar{\lambda}_{k}, N_{k}}^{2}(t) t=\frac{\bar{\lambda}_{k}}{\left(2 \operatorname{deg}\left(S_{k}\right)+1\right)^{2}}
$$

Based on assumption (4.1), we further estimate

$$
\bar{\lambda}_{S_{k}^{2} A_{k}} \equiv \frac{\bar{\lambda}_{k}}{\left(2 \operatorname{deg}\left(S_{k}\right)+1\right)^{2}} \leqslant \frac{\bar{\lambda}_{k}}{4 \operatorname{deg}^{2}\left(S_{k}\right)} \leqslant \frac{\bar{\lambda}_{k}}{4 C_{\operatorname{deg}}^{2} \frac{\bar{\lambda}_{k}}{\lambda_{k+1, k}}} \leqslant \frac{1}{4 C_{\operatorname{deg}}^{2}} \bar{\lambda}_{k+1, k}
$$

thus proving (3.4) with the constant $C=1 /\left(4 C_{\text {deg }}^{2}\right)$. Hence, inequality (2.16), with $\beta$ given by (4.3), follows by Remark 3.1. Estimate (2.17), with $\alpha$ given by (4.2) and $\beta$ given by (4.3), now follows by Theorem 2.1.

## 5. Model example

We consider a model elliptic problem with $H_{0}^{1}$-equivalent form on a bounded polygonal or polyhedral domain $\Omega \subset \mathbb{R}^{d}, d=2$ or $d=3$, that is,

$$
\begin{equation*}
\text { find } u \in H_{0}^{1}(\Omega) \text { such that } a(u, v)=(f, v)_{L_{2}(\Omega)} \forall v \in H_{0}^{1}(\Omega) \tag{5.1}
\end{equation*}
$$

where $f \in L_{2}(\Omega)$ and

$$
\begin{equation*}
c|u|_{H^{1}(\Omega)}^{2} \leqslant a(u, u) \leqslant C|u|_{H^{1}(\Omega)}^{2} \quad \forall u \in H_{0}^{1}(\Omega) \tag{5.2}
\end{equation*}
$$

Further, we consider a system of nested quasiuniform triangulations $\left\{\tau_{h_{k}}\right\}_{k=0}^{l}$ of $\Omega$ ( $\tau_{h_{k}}$ being the refinement of $\tau_{h_{k+1}}$ ) and the corresponding piecewise linear ( $P 1$ ) finite element spaces

$$
H_{0}^{1}(\Omega) \supset V_{h_{0}} \supset V_{h_{1}} \supset \ldots \supset V_{h_{l}}
$$

Here, $h_{k}$ denotes a characteristic meshsize on level $k$. Note that the case of interest is $h_{k} \ll h_{k+1}$. We denote the standard $P 1$ finite element basis of $V_{h_{k}}$ by $\left\{\varphi_{i}^{k}, i=\right.$ $\left.1, \ldots, n_{k}\right\}$, and define standard finite element interpolators in the usual way:

$$
\Pi_{h_{k}}: \mathbf{x} \in \mathbb{R}^{n_{k}} \mapsto \sum_{i=1}^{n_{k}} x_{i} \varphi_{i}^{k}, \quad k=0, \ldots, l .
$$

We assume that the matrix $A=A_{0}$ was obtained by the standard finite element discretization of (5.1) using the finite element basis $\left\{\varphi_{i}^{0}\right\}_{i=1}^{n_{0}}$, that is,

$$
a\left(\Pi_{h_{0}} \mathbf{x}, \Pi_{h_{0}} \mathbf{x}\right)=\left\langle A_{0} \mathbf{x}, \mathbf{x}\right\rangle, \quad \mathbf{x} \in \mathbb{R}^{n_{0}}
$$

The multigrid prolongators are given by

$$
\begin{equation*}
P_{k+1}^{k}=\Pi_{h_{k}}^{-1} \Pi_{h_{k+1}} . \tag{5.3}
\end{equation*}
$$

Note that $P_{k+1}^{k}$ is an $n_{k} \times n_{k+1}$ matrix whose $j$-th column is the basis function $\varphi_{j}^{k+1}$ represented in terms of the basis $\left\{\varphi_{i}^{k}\right\}$ of the immediately finer level. The coarse-level matrices are defined by (2.1), that is,

$$
\begin{aligned}
A_{k} & =\left(P_{k}^{k-1}\right)^{T} A_{k-1} P_{k}^{k-1}=\left(P_{k}^{0}\right)^{T} A P_{k}^{0} \\
\left\langle A_{k} \mathbf{x}, \mathbf{x}\right\rangle & =a\left(\Pi_{h_{k}} \mathbf{x}, \Pi_{k_{k}} \mathbf{x}\right), \quad \forall \mathbf{x} \in \mathbb{R}^{n_{k}}, k=1, \ldots, l
\end{aligned}
$$

Let $Q_{h_{k}}: H_{0}^{1}(\Omega) \rightarrow V_{h_{k}}$ be an $L_{2}(\Omega)$-orthogonal projection. We define $\widetilde{Q}_{k}: \mathbb{R}^{n_{0}} \rightarrow$ $\mathbb{R}^{n_{k}}$ by

$$
\Pi_{h_{k}} \widetilde{Q}_{k}=Q_{h_{k}} \Pi_{h_{0}}, \quad k=0, \ldots, l
$$

For $k=0, \ldots, l$, we set

$$
Q_{k}=P_{k}^{0} \widetilde{Q}_{k}
$$

We will verify the assumptions of Theorem 4.1 for the above linear mappings $Q_{k}$. Namely, we need to verify assumptions (2.14) and (2.15) for our linear mappings $Q_{k}$ and satisfy the assumption (4.1) for smoothers $M_{k}$ whose error propagation operator is given by the polynomial (3.1), where $S_{k}$ is chosen as in (3.11) and (3.8). We will show that our method converges uniformly with respect to the coarsening ratio if the polynomial $S_{k}=p_{\bar{\lambda}_{k}, N_{k}}\left(A_{k}\right)$ in (3.11) has a degree

$$
N_{k}=\operatorname{deg}\left(S_{k}\right) \geqslant C \frac{h_{k+1}}{h_{k}}, \quad C>0
$$

Note that the assumption $c h_{k+1} / h_{k} \leqslant \operatorname{deg}\left(S_{k}\right) \leqslant C h_{k+1} / h_{k}$ is equivalent to

$$
\begin{equation*}
c \frac{h_{k+1}}{h_{k}} \leqslant \operatorname{deg}\left(I-M_{k}^{-1} A_{k}\right)=(2+\gamma) \operatorname{deg}\left(S_{k}\right)+1 \leqslant C \frac{h_{k+1}}{h_{k}} \tag{5.4}
\end{equation*}
$$

(with different constants $c, C>0$ ). Again, we recall that the cases of practical interest are $\gamma=1$ and $\gamma=2$. In any case, we consider $\gamma$ bounded. Thus, in what follows, we assume (5.4).

We will use the following well-known properties of the finite element functions ([5]):

$$
\begin{align*}
\left\|\left(I-Q_{h_{k}}\right) u\right\|_{L_{2}(\Omega)} & \leqslant C h_{k}|u|_{H^{1}(\Omega)} \quad \forall u \in H_{0}^{1}(\Omega),  \tag{5.5}\\
\left|Q_{h_{k}} u\right|_{H^{1}(\Omega)} & \leqslant C|u|_{H^{1}(\Omega)} \quad \forall u \in H_{0}^{1}(\Omega),  \tag{5.6}\\
c\left\|\Pi_{h_{k}} \mathbf{x}\right\|_{L_{2}(\Omega)}^{2} & \leqslant h_{k}^{d}\|\mathbf{x}\| \leqslant C\left\|\Pi_{h_{k}} \mathbf{x}\right\|_{L_{2}(\Omega)}^{2} \quad \forall \mathbf{x} \in \mathbb{R}^{n_{k}},  \tag{5.7}\\
\varrho\left(A_{k}\right) & \leqslant C \max _{i=1, \ldots, n_{k}}\left|\varphi_{k}^{i}\right|_{H^{1}(\Omega)}^{2} \leqslant C h_{k}^{d-2} . \tag{5.8}
\end{align*}
$$

In the estimates to follow, $C, c$ denote generic constants that will depend on the constants in (5.2), (5.5), (5.6), (5.7), (5.8), and (5.4).

First we estimate the value of $\lambda_{k+1, k}$ in (2.14):

$$
\begin{align*}
\lambda_{k+1, k} & =\sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}} \frac{\left\langle A_{k+1} \mathbf{x}, \mathbf{x}\right\rangle}{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k}^{2}}  \tag{5.9}\\
& =\sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}}\left(\frac{\left\langle A_{k+1} \mathbf{x}, \mathbf{x}\right\rangle}{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k+1}^{2}} \cdot \frac{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k+1}^{2}}{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k}^{2}}\right) \\
& \leqslant \sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}} \frac{\left\langle A_{k+1} \mathbf{x}, \mathbf{x}\right\rangle}{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k+1}^{2}} \cdot \sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1} \backslash\{\mathbf{0}\}}} \frac{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k+1}^{2}}{\left\|P_{k+1}^{0} \mathbf{x}\right\|_{k}^{2}} \\
& =\sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}} \frac{\left\langle A_{k+1} \mathbf{x}, \mathbf{x}\right\rangle}{\|\mathbf{x}\|^{2}} \cdot \sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}} \frac{\|\mathbf{x}\|^{2}}{\left\|P_{k}^{0} P_{k+1}^{k} \mathbf{x}\right\|_{k}^{2}} \\
& =\varrho\left(A_{k+1}\right) \sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1} \backslash\{\mathbf{0}\}}} \frac{\|\mathbf{x}\|^{2}}{\left\|P_{k+1}^{k} \mathbf{x}\right\|^{2}} .
\end{align*}
$$

Employing the equivalence (5.7) between the $L_{2}$-norm and the Euclidean norm, together with the definition $P_{k+1}^{k}=\Pi_{h_{k}}^{-1} \Pi_{h_{k+1}}$, we obtain

$$
\left\|\Pi_{h_{k+1}} \mathbf{x}\right\|_{L_{2}(\Omega)}^{2}=\left\|\Pi_{h_{k}} P_{k+1}^{k} \mathbf{x}\right\|_{L_{2}(\Omega)}^{2} \approx h_{k}^{d}\left\|P_{k+1}^{k} \mathbf{x}\right\|^{2}
$$

From here and from (5.7), we have

$$
\left\|P_{k+1}^{k} \mathbf{x}\right\|^{2} \approx h_{k}^{-d}\left\|\Pi_{h_{k+1}} \mathbf{x}\right\|^{2}, \text { and }\|\mathbf{x}\|^{2} \approx h_{k+1}^{-d}\left\|\Pi_{h_{k+1}} \mathbf{x}\right\|^{2}
$$

The last two equivalences, together with (5.9) and (5.8), yield

$$
\begin{align*}
\lambda_{k+1, k} & \leqslant C h_{k+1}^{d-2} \sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}} \frac{\|\mathbf{x}\|^{2}}{\left\|P_{k+1}^{k} \mathbf{x}\right\|^{2}}  \tag{5.10}\\
& \leqslant C h_{k+1}^{d-2} \sup _{\mathbf{x} \in \mathbb{R}^{n_{k+1}} \backslash\{\mathbf{0}\}} \frac{h_{k+1}^{-d}\left\|\Pi_{h_{k+1}} \mathbf{x}\right\|_{L_{2}(\Omega)}^{2}}{h_{k}^{-d}\left\|\Pi_{h_{k+1}} \mathbf{x}\right\|_{L_{2}(\Omega)}^{2}} \\
& \leqslant \bar{\lambda}_{k+1, k} \equiv C \frac{h_{k}^{d}}{h_{k+1}^{2}} .
\end{align*}
$$

(We take the final estimate as an upper bound $\bar{\lambda}_{k+1, k} \geqslant \lambda_{k+1, k}$, see Theorem 2.1.) To verify (2.14), we further estimate using (5.5), (5.7), (2.12), $P_{k+1}^{k}=\Pi_{h_{k}}^{-1} \Pi_{h_{k+1}}$, $\Pi_{h_{k}} \widetilde{Q}_{k}=Q_{h_{k}} \Pi_{h_{0}}, Q_{k}=P_{k}^{0} \widetilde{Q}_{k}$, the fact that $Q_{h_{k}}: H^{1}(\Omega) \rightarrow V_{h_{k}}$ is an $L_{2}(\Omega)-$ orthogonal projection and $V_{h_{k+1}} \subset V_{h_{k}}$ :

$$
\begin{aligned}
\left\|\left(Q_{k}-Q_{k+1}\right) \mathbf{v}\right\|_{k}^{2} & =\left\|\left(P_{k}^{0} \widetilde{Q}_{k}-P_{k+1}^{0} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{k}^{2} \\
& =\left\|P_{k}^{0}\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{k}^{2} \\
& =\left\|\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|^{2} \\
& \approx h_{k}^{-d}\left\|\Pi_{h_{k}}\left(\widetilde{Q}_{k}-P_{k+1}^{k} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{L_{2}(\Omega)}^{2} \\
& =h_{k}^{-d}\left\|\left(\Pi_{h_{k}} \widetilde{Q}_{k}-\Pi_{h_{k+1}} \widetilde{Q}_{k+1}\right) \mathbf{v}\right\|_{L_{2}(\Omega)}^{2} \\
& =h_{k}^{-d}\left\|\left(Q_{h_{k}}-Q_{h_{k+1}}\right) \Pi_{h_{0}} \mathbf{v}\right\|_{L_{2}(\Omega)}^{2} \\
& \leqslant h_{k}^{-d}\left(\left\|\left(I-Q_{h_{k}}\right) \Pi_{h_{0}} \mathbf{v}\right\|_{L_{2}(\Omega)}^{2}+\left\|\left(Q_{h_{k}}-Q_{h_{k+1}}\right) \Pi_{h_{0}} \mathbf{v}\right\|_{L_{2}(\Omega)}^{2}\right) \\
& =h_{k}^{-d}\left\|\left(I-Q_{h_{k+1}}\right) \Pi_{h_{0}} \mathbf{v}\right\|_{L_{2}(\Omega)}^{2} \\
& \leqslant C \frac{h_{k+1}^{2}}{h_{k}^{d}}\left|\Pi_{h_{0}} \mathbf{v}\right|_{H^{1}(\Omega)}^{2} .
\end{aligned}
$$

Since

$$
\left|\Pi_{h_{0}} \mathbf{v}\right|_{H^{1}(\Omega)}^{2} \approx a\left(\Pi_{h_{0}} \mathbf{v}, \Pi_{h_{0}} \mathbf{v}\right)=\|\mathbf{v}\|_{A}^{2},
$$

and from (5.10), i.e.

$$
\lambda_{k+1, k} \leqslant \bar{\lambda}_{k+1, k} \equiv C \frac{h_{k}^{d}}{h_{k+1}^{2}}
$$

we obtain

$$
\left\|\left(Q_{k}-Q_{k+1}\right) \mathbf{v}\right\|_{k}^{2} \leqslant \frac{C}{\overline{\lambda_{k+1, k}}}\|\mathbf{v}\|_{A}^{2}
$$

proving (2.14).
To verify (2.15), we use (5.6), (5.2) and $\Pi_{h_{k}} \widetilde{Q}_{k}=Q_{h_{k}} \Pi_{h_{0}}$ and observe that

$$
\begin{aligned}
\left\|Q_{k} \mathbf{v}\right\|_{A}^{2} & =\left\|\widetilde{Q}_{k} \mathbf{v}\right\|_{A_{k}}^{2}=a\left(\Pi_{h_{k}} \widetilde{Q}_{k} \mathbf{v}, \Pi_{h_{k}} \widetilde{Q}_{k} \mathbf{v}\right)=a\left(Q_{h_{k}} \Pi_{h_{0}} \mathbf{v}, Q_{h_{k}} \Pi_{h_{0}} \mathbf{v}\right) \\
& \leqslant C\left|Q_{h_{k}} \Pi_{h_{0}} \mathbf{v}\right|_{H^{1}(\Omega)}^{2} \leqslant C\left|\Pi_{h_{0}} \mathbf{v}\right|_{H^{1}(\Omega)}^{2} \leqslant C a\left(\Pi_{h_{0}} \mathbf{v}, \Pi_{h_{0}} \mathbf{v}\right)=C\|\mathbf{v}\|_{A}^{2}
\end{aligned}
$$

To satisfy (4.1), it is sufficient to use, in the definition (3.11) of $S_{k}$, a polynomial $p_{\bar{\lambda}_{k}, N_{k}}$ of sufficiently large degree. From (5.10), we have

$$
\bar{\lambda}_{k+1, k} \equiv C \frac{h_{k}^{d}}{h_{k+1}^{2}} \geqslant \lambda_{k+1, k}
$$

Further, due to (5.8), we can take

$$
\bar{\lambda}_{k}=C h_{k}^{d-2} \geqslant \varrho\left(A_{k}\right)
$$

Thus, we have

$$
\frac{\bar{\lambda}_{k}}{\bar{\lambda}_{k+1, k}}=C\left(\frac{h_{k+1}}{h_{k}}\right)^{2}
$$

and to satisfy (4.1), we need

$$
\operatorname{deg}\left(S_{k}\right) \geqslant C \frac{h_{k+1}}{h_{k}}
$$

which is guaranteed by (5.4). Thus the assumptions of Theorem 4.1 are verified whenever $\gamma=1$ or $\gamma$ is even.

We summarize the above results in the following theorem:
Theorem 5.1. Consider the model elliptic problem and coarse spaces derived from nested quasiuniform triangulations as described in this section, with the intergrid transfer operators defined by the natural embedding of the spaces (5.3). Assume the error propagation operators of both the pre- and post-smoother are given on each level $k=0, \ldots, l-1$ by (3.1), with $S_{k}$ defined by (3.11), (3.8), and either $\gamma=1$ or $\gamma>0$ even. We assume $\gamma$ is bounded. In addition, we assume that the degree of the smoothing polynomial satisfies (5.4). Then the resulting multigrid operator, $B$, is nearly spectrally equivalent to $A$, that is,

$$
\mathbf{v}^{T} A \mathbf{v} \leqslant \mathbf{v}^{T} B \mathbf{v} \leqslant C l \mathbf{v}^{T} A \mathbf{v} \quad \forall \mathbf{v} \in \mathbb{R}^{n_{0}}
$$

where the constant $C$ is independent on the meshsizes $h_{k}$ (and the coarsening ratio $\left.h_{k+1} / h_{k}\right)$.

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