# MULTILEVEL METHODS FOR THE APPROXIMATION OF SINGULAR SOLUTIONS OF COMPLETELY CONTINUOUS OPERATOR EQUATIONS 

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#### Abstract

The direct techniques for approximating singular solutions of nonlinear equations depending on parameters transform the original problem into that of solving a suitable augmented system. Dealing with equations involving completely continuous operators, a multilevel approach to the solution of these larger systems is presented.


1. Introduction. We are concerned with the approximation of singular points in branches of solutions of nonlinear parameterdependent equations having the form

$$
\begin{equation*}
F\left(u, \beta_{1}, \ldots, \beta_{p-1}, \gamma\right)=u-K\left(u, \beta_{1}, \ldots, \beta_{p-1}, \gamma\right)=0 \tag{1.1}
\end{equation*}
$$

where $u \in U$, with $U$ a real Banach space, $\gamma \in \mathbf{R}$ and the $\beta_{i}$ 's are ( $p-1$ )-additional real parameters, for some $p \geq 1$. From now on, we assume that $F$ is a $C^{\nu}$-mapping, with $\nu \geq 3$, from an open set $D \subset U \times \mathbf{R}^{p}$ into $U$ and that the operator $K$ is completely continuous.

As is well known, the direct techniques for the approximation of a singular solution are based on the construction of an augmented system having it as a regular solution (cf. [9] for a general discussion). Then, moving from a suitable starting point, obtained for instance by a continuation procedure, the arising system is solved by an iterative method. In several practical cases, this procedure is carried out using operator approximations, on which the accuracy of the computed solution depends (cf. [11, 20]).

In this paper we deal with singular solutions of (1.1), such as turning points and, when unfolded, bifurcation points, which fall within the unifying theory recently developed by Griewank and Reddien in [10, 11], where the corresponding augmented system is built up through the minimum number of additional scalar equations characterizing the singularity. For solving this larger system, we present an approach based

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on the use of iterative multilevel methods, with special attention to the case where collectively compact approximations of $K$ are employed.

The basic approximation scheme is presented and studied in Section 2. In Section 4 we illustrate how it can be suitably combined with classical linear multigrid algorithms, briefly outlined in Section 3.
2. The basic approximation scheme. From now on, for ease of notation, we set $X=U \times \mathbf{R}^{p}$ and $Z=U \times \mathbf{R}^{p-1}(Z=U$, if $p=1)$. Moreover, we denote by $Q_{U}$ and $Q_{Z}$ the coordinate projections from $X$ into $U$ and into $Z$, respectively. That is, we set $Q_{U} \mathbf{x}=u, Q_{Z} \mathbf{x}=$ $\left(u, \beta_{1}, \ldots, \beta_{p-1}\right)^{T}$, for $\mathbf{x}=\left(u, \beta_{1}, \ldots, \beta_{p-1}, \gamma\right)^{T} \in X$. Accordingly, we rewrite equation (1.1) in the form

$$
\begin{equation*}
F(\mathbf{x})=Q_{U} \mathbf{x}-K(\mathbf{x})=0, \quad \mathbf{x} \in X \tag{2.1}
\end{equation*}
$$

Throughout, the first and second derivatives of $F$ with respect to $\mathbf{x}$ will be indicated by $F^{\prime}$ and $F^{\prime \prime}$, respectively, and we will use subscripts for indicating the partial derivatives with respect to each variable. In particular, $F_{z}$ will denote the derivative of $F$ with respect to $Q_{Z} \mathbf{x}$. As usual, for $i=1, \ldots, p, e_{i}$ will be the $i^{\text {th }}$ coordinate vector in $\mathbf{R}^{p}$; any product space $U \times \mathbf{R}^{m}$ will be endowed with a product norm; $|\cdot|_{E}$ will stand for the norm in a Banach space $E$, except for the space $\mathbf{L}\left(E, E^{\prime}\right)$ ( $E$ and $E^{\prime}$ Banach spaces) of all bounded linear operators from $E$ into $E^{\prime}$ whose norm will be denoted by $\|\cdot\|_{E \rightarrow E^{\prime}} ; E^{*}$ will stand for the dual space of $E$ and $B^{*}$ will be the adjoint of a linear operator $B$. The notation $S(\mathbf{x}, \delta)$ will indicate the closed ball in $X$ of center $\mathbf{x}$ and radius $\delta$. Finally, $\mathbf{e} \in X^{*}$ will be such that

$$
\mathbf{e x}=\gamma, \quad \text { for } \mathbf{x}=\left(u, \beta_{1}, \ldots, \beta_{p-1}, \gamma\right)^{T} \in X
$$

Let $\mathbf{x}^{s} \in D$ be a singular point of (2.1) in the sense that:

$$
\begin{equation*}
F\left(\mathbf{x}^{s}\right)=0 \tag{1}
\end{equation*}
$$

$$
\begin{equation*}
\operatorname{Ker}\left(F_{z}\left(\mathbf{x}^{s}\right)\right)=\operatorname{span}\left\{Q_{Z} \mathbf{y}_{1}^{s}, \ldots, Q_{Z} \mathbf{y}_{p}^{s}\right\} \tag{2}
\end{equation*}
$$

with $e_{i}^{T} M Q_{Z} \mathbf{y}_{j}^{s}=\delta_{i j}$ (the Kronecker's delta), for $i, j=1, \ldots, p$, where $M$ is a suitable linear mapping acting from $Z$ into $\mathbf{R}^{p}$;

$$
\begin{gather*}
F_{\gamma}\left(\mathbf{x}^{s}\right) \notin \text { Range } F_{z}\left(\mathbf{x}^{s}\right) \quad \text { and }  \tag{3}\\
\text { Range }\left(F_{z}\left(\mathbf{x}^{s}\right)\right)=\left\{u \in U: \psi^{s} u=0\right\}
\end{gather*}
$$

for some $\psi^{s}(\neq 0) \in U^{*}$. We also assume that
the $p \times p$ matrix $R^{s}=\left[r_{i j}^{s}\right]$ having entries

$$
\begin{equation*}
r_{i j}^{s}=\psi^{s} F_{z z}\left(\mathbf{x}^{s}\right) Q_{Z} \mathbf{y}_{i}^{s} Q_{Z} \mathbf{y}_{j}^{s}, \quad \text { for } i, j=1, \ldots, p \tag{4}
\end{equation*}
$$

is nonsingular.
These conditions characterize various types of singular solutions (see $[\mathbf{7}, \mathbf{1 0}-\mathbf{1 2}, 15,18,20]$ ), which Griewank and Reddien in $[10,11]$ called generalized turning points. As mentioned in the introduction, for characterizing $\mathbf{x}^{s}$ as a regular solution of a suitable augmented system, we adopt here the following approach proposed in $[\mathbf{1 0}, \mathbf{1 1}]$.

For any $\mathbf{x} \in D$, we consider the mapping $A(\mathbf{x}): X \rightarrow X$ defined by

$$
\begin{equation*}
A(\mathbf{x}) \mathbf{y}=\left(F^{\prime}(\mathbf{x}) \mathbf{y}, M Q_{Z} \mathbf{y}\right)^{T}, \quad \text { for } \mathbf{y} \in X \tag{2.2}
\end{equation*}
$$

Then, the following result can be proved (see [11]).

Proposition 2.1. There exists $\delta^{*}$ such that, for every $\mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$, A(x) has a bounded inverse. Moreover, $\mathbf{x}^{s}$ is a regular solution of the following augmented system (from $X$ into itself)

$$
\begin{gather*}
F(\mathbf{x})=0  \tag{2.3}\\
\mathbf{e y}_{i}(\mathbf{x})=0, \quad \text { for } i=1, \ldots, p \tag{2.4}
\end{gather*}
$$

where, for each $i=1, \ldots, p, \mathbf{y}_{i}(\mathbf{x}) \in X$ solves the equation

$$
A(\mathbf{x}) \mathbf{y}_{i}(\mathbf{x})=\left(0, e_{i}\right)^{T}, \quad 0 \in U
$$

From now on, let $\left\{F^{n}\right\}, n=0,1, \ldots$, be a sequence of nonlinear operators from $D$ into $U$ satisfying the following assumptions:
$\left(\mathrm{k}_{1}\right) \quad F^{n}$ is pointwise convergent to $F$ on $D$, i.e., for each $\mathbf{x} \in D$,

$$
F^{n}(\mathbf{x}) \rightarrow F(\mathbf{x}), \quad \text { as } n \rightarrow \infty
$$

$\left(\mathrm{k}_{2}\right)$
$\left\{Q_{U}-F^{n}\right\}$ is a family of collectively compact (c.c.) operators (cf. [1]);
$\mathrm{k}_{3}$ for every $n, F^{n}$ is of class $C^{\nu}, \nu \geq 3$, on $S\left(\mathbf{x}^{s}, \delta^{*}\right)$ and the derivatives $F^{n^{\prime}}(\mathbf{x}), F^{n^{\prime \prime}}(\mathbf{x})$ and $F^{n^{\prime \prime \prime}}(\mathbf{x})$ are uniformly bounded with respect to $\mathbf{x}$ and $n$.

Accordingly, we define the sequence of operators $\left\{A^{n}\right\}$ as

$$
A^{n}(\mathbf{x}) \mathbf{v}=\left(F^{n \prime}(\mathbf{x}) \mathbf{v}, M Q_{Z} \mathbf{v}\right)^{T}, \quad \text { for } \mathbf{v} \in X
$$

We assume that for all $n$ sufficiently large:
$\left(\mathrm{k}_{4}\right)$ for every $\mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$, the operators $A^{n}(\mathbf{x})^{-1}$ exist and they are uniformly bounded with respect to $\mathbf{x}$ and $n$.

Then for solving (2.3)-(2.4), we consider the following approximation scheme:

Scheme 1. Select $n^{*}$ and take a subsequence of $\left\{F^{n}: n \geq n^{*}\right\}$, we relabel here as $\left\{F^{k}\right\}$, for $k=0,1, \ldots$. Starting from $\mathbf{x}^{0} \in D$, for $k=0,1, \ldots$, do the following:

1. Select an approximation $B^{k}$ of $A^{k}\left(\mathbf{x}^{k}\right)^{-1}$.
2. Take

$$
\mathbf{s}^{k}=-B^{k}\left(F^{k}\left(\mathbf{x}^{k}\right), 0\right)^{T}
$$

3. For each $i=1, \ldots, p$, take

$$
\mathbf{y}_{i}^{k}=B^{k}\left(0, e_{i}\right)^{T} .
$$

4. Construct the $p \times p$ matrix $R^{k}$ having entries

$$
\rho_{i j}^{k}=\left[\left(B^{k}\right)^{*} \mathbf{e}\right]\left(-F^{k \prime \prime}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k} \mathbf{y}_{j}^{k}, 0\right)^{T}, \quad \text { for } i, j=1, \ldots, p
$$

and the $p$-dimensional vector $c^{k}$ having entries

$$
\eta_{i}^{k}=-\mathbf{e} \mathbf{y}_{i}^{k}+\left[\left(B^{k}\right)^{*} \mathbf{e}\right]\left[\left(F^{k \prime \prime}\left(\mathbf{x}^{k}\right) \mathbf{s}^{k} \mathbf{y}_{i}^{k}, 0\right)^{T}+A^{k}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}-\left(0, e_{i}\right)^{T}\right]
$$

for $i=1, \ldots, p$.
5. Find $a^{k}=\left(\alpha_{1}^{k}, \ldots, \alpha_{p}^{k}\right)^{T}$ which solves

$$
\begin{equation*}
\mathbf{R}^{k} a^{k}=c^{k} \tag{2.5}
\end{equation*}
$$

6. Set

$$
\mathrm{x}^{k+1}=\mathrm{x}^{k}+\mathrm{s}^{k}+\Sigma_{1 \leq i \leq p} \alpha_{i}^{k} \mathbf{y}_{i}^{k}
$$

In order to simplify the notation, we set

$$
\begin{aligned}
& (\Delta F)^{k}:=\left|F^{k}\left(\mathbf{x}^{s}\right)\right|_{U} \\
& (\Delta A)^{k}=\max _{1 \leq i \leq p}\left|\left(A^{k}\left(\mathbf{x}^{s}\right)-A\left(\mathbf{x}^{s}\right)\right) \mathbf{y}_{i}^{s}\right|_{X} \\
& \vartheta^{k}=\| I-\left.B^{k} A^{k}\left(\mathbf{x}^{k}\right)\right|_{X \rightarrow X} \\
& \mathbf{e}_{\mathbf{x}}^{k}=\left|\mathbf{x}^{s}-\mathbf{x}^{k}\right|_{X} \\
& \mathbf{e}_{\mathbf{y}}^{k}=\max _{1 \leq i \leq p}\left|\mathbf{y}_{i}^{s}-\mathbf{y}_{i}^{k}\right|_{X}
\end{aligned}
$$

Theorem 2.1. Let $\left(\mathrm{h}_{1}\right)-\left(\mathrm{h}_{4}\right)$ and $\left(\mathrm{k}_{1}\right)-\left(\mathrm{k}_{4}\right)$ hold. There exist constants $\delta$ and $\vartheta$ such that, if $\mathbf{x}^{0} \in S\left(\mathbf{x}^{s}, \delta\right), \vartheta^{k} \leq \vartheta$ for every $k$, and $n^{*}$ is sufficiently large, then Scheme 1 is well defined and $\mathbf{x}^{k} \in S\left(\mathbf{x}^{s}, \delta\right)$, for every $k$. Moreover, there are two constants $c_{\mathbf{y}}$ and $c_{\mathbf{x}}$ such that, for every $k$,

$$
\begin{equation*}
e_{\mathbf{y}}^{k} \leq c_{\mathbf{y}}\left[\vartheta^{k}+e_{\mathbf{x}}^{k}+(\Delta A)^{k}\right] \tag{2.6}
\end{equation*}
$$

and

$$
\begin{equation*}
e_{\mathbf{x}}^{k+1} \leq c_{\mathbf{x}}\left[\left(e_{\mathbf{x}}^{k}+\vartheta^{k}\right)\left(e_{\mathbf{x}}^{k}+e_{\mathbf{y}}^{k}\right)+(\Delta F)^{k}+(\Delta A)^{k}\right] \tag{2.7}
\end{equation*}
$$

Proof. If $n^{*}$ is sufficiently large, then, by $\left(\mathrm{k}_{4}\right)$, there is a constant $\mu$ such that, for every $k$, if $\mathbf{x}^{k} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$, then $A^{k}\left(\mathbf{x}^{k}\right)$ is invertible and

$$
\left\|\left(A^{k}\left(\mathbf{x}^{k}\right)\right)^{-1}\right\|_{X \rightarrow X} \leq \mu
$$

Then

$$
B^{k}-\left(A^{k}\left(\mathrm{x}^{k}\right)\right)^{-1}=-\left[I-B^{k} A^{k}\left(\mathrm{x}^{k}\right)\right]\left(A^{k}\left(\mathrm{x}^{k}\right)\right)^{-1}
$$

and, since $\vartheta^{k} \leq \vartheta$, it follows that

$$
\left\|B^{k}\right\|_{X \rightarrow X} \leq(1+\vartheta) \mu
$$

By $\left(\mathrm{h}_{1}\right)-\left(\mathrm{h}_{4}\right)$, using well-known results about approximations of completely continuous operators (see $[\mathbf{1}, \mathbf{1 7}]$ ), since for each $i=1, \ldots, p$, $\mathbf{y}_{i}\left(\mathbf{x}^{s}\right)=\mathbf{y}_{i}^{s}\left(\right.$ cf. (2.4)), we realize that, if $\mathbf{x}^{k} \in S\left(\mathbf{x}^{s}, \delta\right)$, if $\delta$ and $\vartheta$ are sufficiently small and if $n^{*}$ is sufficiently large, then the coefficient matrix $R^{k}$ of (2.5) is nonsingular and there is a constant $\tau$, independent of $k$, such that

$$
\begin{equation*}
\left\|\left(R^{k}\right)^{-1}\right\| \leq \tau \tag{2.8}
\end{equation*}
$$

(Here || || denotes the matrix norm induced by the 1-norm on $\mathbf{R}^{p}$ ). Hence, relation (2.6) easily follows and, moreover, by (2.8), $\mathbf{x}^{k+1}$ is well defined. Now let us prove (2.7).

Let $\alpha_{i}^{\prime k}$, for $i=1, \ldots, p$, be such that

$$
\begin{equation*}
\Sigma_{1 \leq i \leq p} \alpha_{i}^{\prime k} e_{i}=M Q_{Z}\left(\mathbf{x}^{s}-\mathbf{x}^{k}\right) \tag{2.9}
\end{equation*}
$$

Clearly, the $\alpha_{i}^{\prime k}$ 's are uniquely determined.
Then we define $\mathbf{x}^{\prime k}$ by

$$
\mathbf{x}^{\prime k}=\mathbf{x}^{k}+\mathbf{s}^{k}+\Sigma_{1 \leq i \leq p} \alpha_{i}^{\prime k} \mathbf{y}_{i}^{k}
$$

By (2.9), we have

$$
\begin{aligned}
\mathbf{x}^{\prime k}-\mathbf{x}^{k}= & -B^{k}\left[\left(F^{k}\left(\mathbf{x}^{k}\right), 0\right)^{T}-\Sigma_{1 \leq i \leq p} \alpha_{i}^{\prime k}\left(0, e_{i}\right)^{T}\right] \\
= & B^{k}\left(F^{k}\left(\mathbf{x}^{s}\right)-F^{k}\left(\mathbf{x}^{k}\right)-F^{k \prime}\left(\mathbf{x}^{k}\right)\left(\mathbf{x}^{s}-\mathbf{x}^{k}\right), 0\right)^{T} \\
& \left.+B^{k} A^{k}\left(\mathbf{x}^{k}\right)\left(\mathbf{x}^{s}-\mathbf{x}^{k}\right)-B^{k}\left(F^{k}\left(\mathbf{x}^{s}\right)\right), 0\right)^{T}
\end{aligned}
$$

Hence, we easily get

$$
\begin{equation*}
\left|x^{\prime k}-\mathbf{x}^{s}\right|_{X} \leq c\left[\left|\mathbf{x}^{s}-\mathbf{x}^{k}\right|_{X}^{2}+(\Delta F)^{k}+\vartheta^{k}\left|\mathbf{x}^{s}-\mathbf{x}^{k}\right|_{X}\right] \tag{2.10}
\end{equation*}
$$

Observing that the $\mathbf{y}_{i}^{k}$ 's are uniformly bounded, we have, for some constant $c_{0}$, independent of $k$,

$$
\begin{equation*}
\left|\mathbf{x}^{s}-\mathbf{x}^{k+1}\right|_{X} \leq c_{0}\left[\left|\mathbf{x}^{s}-\mathbf{x}^{\prime k}\right|_{X}+\Sigma_{1 \leq i \leq p}\left|\alpha_{i}^{k}-\alpha_{i}^{\prime k}\right|\right] \tag{2.11}
\end{equation*}
$$

Now we provide an upper bound for $\Sigma_{1 \leq i \leq p}\left|\alpha_{i}^{k}-\alpha_{i}^{\prime k}\right|$. From the identities

$$
\Sigma_{1 \leq j \leq p} \rho_{i j}^{k} \alpha_{j}^{k}=\mathbf{e}\left[-\mathbf{y}_{i}^{k}+B^{k}\left(\left(F^{k \prime \prime}\left(\mathbf{x}^{k}\right) \mathbf{s}^{k} \mathbf{y}_{i}^{k}, 0\right)^{T}+A^{k}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}-\left(0, e_{i}\right)^{T}\right)\right]
$$

and

$$
\Sigma_{1 \leq j \leq p} \rho_{i j}^{k} \alpha_{j}^{\prime k}=\Sigma_{1 \leq j \leq p} \alpha_{j}^{\prime k} \mathbf{e} B^{k}\left(-F^{k^{\prime \prime}}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k} \mathbf{y}_{j}^{k}, 0\right)^{T}
$$

for every $i=1, \ldots, p$, we get

$$
\begin{align*}
& \Sigma_{1 \leq j \leq p} \rho_{i j}^{k}\left(\alpha_{j}^{k}-\alpha_{j}^{\prime k}\right)  \tag{2.12}\\
& =e\left[-\mathbf{y}_{i}^{k}+B^{k}\left(\left(F^{k \prime \prime}\left(\mathbf{x}^{k}\right)\left(\mathbf{x}^{k^{\prime}}-\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}, 0\right)^{T}+A^{k}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}-\left(0, e_{i}\right)^{T}\right)\right]
\end{align*}
$$

Now, recalling that, by $\left(\mathrm{h}_{3}\right), \mathbf{e y}_{i}^{s}=0$ and that $\left.A\left(\mathbf{x}^{s}\right)\right) \mathbf{y}_{i}^{s}=\left(0, e_{i}\right)^{T}$, for every $i=1, \ldots, p$, we obtain the following inequality

$$
\begin{aligned}
& \left|\mathbf{e}\left[-\mathbf{y}_{i}^{k}+B^{k}\left(\left(F^{k^{\prime \prime}}\left(\mathbf{x}^{k}\right)\left(\mathbf{x}^{s}-\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}, 0\right)^{T}+A^{k}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}-\left(0, e_{i}\right)^{T}\right)\right]\right| \\
& \quad \leq\left|\mathbf{e} B^{k}\left(-F^{k^{\prime}}\left(\mathbf{x}^{s}\right) \mathbf{y}_{i}^{k}+F^{k^{\prime \prime}}\left(\mathbf{x}^{k}\right)\left(\mathbf{x}^{s}-\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}+F^{k^{\prime}}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}, 0\right)^{T}\right| \\
& +\mid \mathbf{e}\left[\left(\mathbf{y}_{i}^{s}-\mathbf{y}_{i}^{k}\right)-B^{k}\left(A^{k}\left(\mathbf{x}^{s}\right)\left(\mathbf{y}_{i}^{s}-\mathbf{y}_{i}^{k}\right)+B^{k}\left(A^{k}\left(\mathbf{x}^{s}\right)-A\left(\mathbf{x}^{s}\right)\right) \mathbf{y}_{i}^{s}\right] \mid\right.
\end{aligned}
$$

Hence, we easily get, for every $i$,

$$
\begin{align*}
\mid \mathbf{e}\left[-\mathbf{y}_{i}^{k}+B^{k}\left(\left(F^{k^{\prime \prime}}\right.\right.\right. & \left.\left.\left.\left.\left(\mathbf{x}^{k}\right)\left(\mathbf{x}^{k^{\prime}}-\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}\right) 0\right)^{T}+A^{k}\left(\mathbf{x}^{k}\right) \mathbf{y}_{i}^{k}-\left(0, e_{i}\right)^{T}\right)\right] \mid  \tag{2.13}\\
& \leq c_{1}\left[\left(e_{\mathbf{x}}^{k}\right)^{2}+\vartheta^{k} e_{\mathbf{y}}^{k}+e_{\mathbf{x}}^{k} e_{\mathbf{y}}^{k}+(\Delta A)^{k}+\left|\mathbf{x}^{\prime k}-\mathbf{x}^{s}\right|_{X}\right]
\end{align*}
$$

for some constant $c_{1}$ independent of $k$. Thus, by (2.10)-(2.13), we obtain (2.7). Clearly, $\mathbf{x}^{k+1} \in S\left(\mathbf{x}^{s}, \delta\right)$, for a sufficiently large $n^{*}$ and for sufficiently small $\delta$ and $\vartheta$.
3. A review of linear multilevel methods. In this section we will consider some linear multilevel (here we will use the term multigrid) procedures which can be employed for constructing the operators $B^{k}$ and $\left(B^{k}\right)^{*}$ in Scheme 1.
At first, we give a general formulation of the linear multigrid methods which includes various classical procedures, like for instance those proposed by Hackbusch [12] and Hemker and Schippers [13].

Let $\left\{X_{i}\right\}, i=0,1, \ldots$, be a sequence of Banach spaces, let $I_{i}$ denote the identity operator on $X_{i}$, and let us suppose that:

1) $\left\{\pi_{i}\right\}, \pi_{i}: X_{i-1} \rightarrow X_{i}, i=1,2, \ldots$, is a sequence of linear operators (prolongations), such that

$$
\left\|\pi_{i}\right\|_{X_{i-1} \rightarrow X_{i}} \leq c_{\pi}, \quad \text { for every } i
$$

2) $\left\{\rho_{i}\right\}, \rho_{i}: X_{i} \rightarrow X_{i-1}, i=1,2, \ldots$, is a sequence of linear operators (restrictions) such that, for every $i$,

$$
\left\|\rho_{i}\right\|_{X i \rightarrow X i-1} \leq c_{\rho}
$$

and

$$
\rho_{i} \pi_{i}=I_{i-1}
$$

3) $\left\{A_{i}\right\},\left\{\Gamma_{i}\right\},\left\{\Psi_{i}\right\},\left\{\Phi_{i}\right\}$, for $i=0,1, \ldots$, are uniformly bounded sequences of linear operators, with $A_{i}, \Gamma_{i}, \Psi_{i}, \Phi_{i}: X_{i} \rightarrow X_{i}$. We assume that the inverses of the operators $A_{i}$ exist and are uniformly bounded.
Then, for every $i=1,2, \ldots$, we consider the approximation $B_{i}$ of $A_{i}^{-1}$ generated by the following

General linear multigrid scheme. Let an approximation $B_{0}$ of $A_{0}^{-1}$ and an integer $\gamma \geq 1$ be given. For $i=1,2, \ldots$, set

$$
\begin{gather*}
S_{i-1}=I_{i-1}-B_{i-1} A_{i-1},  \tag{3.1}\\
Q_{i-1}=\Sigma_{0 \leq j \leq \gamma-1}\left(S_{i-1}\right)^{j} B_{i-1}  \tag{3.2}\\
B_{i}=\Gamma_{i}+\Psi_{i} \pi_{i} Q_{i-1} \rho_{i} \Phi_{i} . \tag{3.3}
\end{gather*}
$$

For the study of this scheme, the following operators will also be considered:

$$
\begin{gathered}
R_{i}=I_{i}-\Gamma_{i} A_{i}-\Psi_{i} \pi_{i} A_{i-1}^{-1} \rho_{i} \Phi_{i} A_{i} \\
W_{i}=I_{i}-\Gamma_{i} A_{i}+\left(I_{i}-\Psi_{i}\right) \pi_{i} A_{i-1}^{-1} \rho_{i} \Phi_{i} A_{i}
\end{gathered}
$$

Straightforward computations show that, since $Q_{i-1}=\left(I_{i-1}-\left(S_{i-1}\right)^{\gamma}\right)$. $A_{i-1}^{-1}$ and $\rho_{i} \pi_{i}=I_{i-1}$, we have

$$
\begin{align*}
S_{i} & =R_{i}+\Psi_{i} \pi_{i}\left(S_{i-1}\right)^{\gamma} A_{i-1}^{-1} \rho_{i} \Phi_{i} A_{i} \\
& =R_{i}+\Psi_{i} \pi_{i}\left(S_{i-1}\right)^{\gamma} \rho_{i}\left(W_{i}-R_{i}\right) \tag{3.4}
\end{align*}
$$

Remark. The above scheme contains the classical methods described in [12], where the coarse-grid corrections are combined with some presmoothing and post-smoothing steps. Indeed, for every $i$, let $M_{i}$ and $N_{i}$ be linear operators, from $X_{i}$ into itself, such that $M_{i}=\left(I_{i}-N_{i} A_{i}\right)$. Then, for some integers $m \geq q>0$, set

$$
\begin{aligned}
\Gamma_{i} & =\Sigma_{0 \leq j \leq m-1} M_{i}^{j} N_{i} \\
\Psi_{i} & =M_{i}^{m-q} \\
\Phi_{i} & =I_{i}-A_{i} \Sigma_{0 \leq j \leq q-1} M_{i}^{j} N_{i}
\end{aligned}
$$

Accordingly, we have

$$
\begin{gathered}
R_{i}=M_{i}^{m-q}\left[I_{i}-\pi_{i} A_{i-1}^{-1} \rho_{i} A_{i}\right] M_{i}^{q} \\
W_{i}=M_{i}^{m-q}-\left[\left(I_{i}-M_{i}^{m-q}\right) \pi_{i} A_{i-1}^{-1} \rho_{i} A_{i}\right] M_{i}^{q}
\end{gathered}
$$

In particular, choosing, for each $i, M_{i}=I_{i}, m=q=1$, we obtain certain multigrid methods for equations of the second kind discussed in [13].

We come back to the study of the scheme defined by (3.1)-(3.3). Let us assume that $\left\{v_{i}\right\}, i=0,1, \ldots$, is a sequence such that, for each $i=1,2, \ldots$,

$$
\left\|R_{i}\right\|_{X_{i} \rightarrow X_{i}} \leq v_{i}
$$

and

$$
\left\|S_{0}\right\|_{X_{0} \rightarrow X_{0}} \leq v_{0}
$$

Thus, setting

$$
\kappa=c_{\pi} c_{\rho} \operatorname{Sup}_{i}\left\|\Psi_{i}\right\|_{X_{i} \rightarrow X_{i}}
$$

and

$$
s_{i}=\left\|S_{i}\right\|_{X_{i} \rightarrow X_{i}}, \quad w=\operatorname{Sup}_{i}\left\|W_{i}\right\|_{X_{i} \rightarrow X_{i}}
$$

from (3.4) it follows that

$$
\begin{equation*}
s_{i} \leq v_{i}+\kappa\left(s_{i-1}\right)^{\gamma}\left(w+v_{i}\right), \quad \text { for } i=1,2, \ldots, \tag{3.5}
\end{equation*}
$$

with

$$
\begin{equation*}
s_{0} \leq v_{0} \tag{3.6}
\end{equation*}
$$

Arguing as in [13, Lemma 3.3], one proves the following result.

Proposition 3.1. Let (3.5) and (3.6) hold. Assume that $v_{i} \leq v$ for every $i$, set $d_{i}=v_{i} / v_{i-1}$ and then set $d=\inf _{i}\left(d_{i}\right)$. If either,

$$
\gamma \geq 1, \quad 0<\kappa w<d<1, \quad \text { and } \quad v \leq(d-\kappa w) /(d+\kappa d)
$$

or

$$
\begin{equation*}
\gamma=2, \quad \text { and } \quad 4 \kappa v_{i-1}\left(w / d_{i}+v_{i-1}\right) \leq 1, \quad \text { for every } i \tag{3.7}
\end{equation*}
$$

then there is a $c>0$ such that, for every $i$, we have

$$
s_{i} \leq c v_{i}
$$

More precisely, we can take, in the first case, $c=1+\kappa(w+d) /(d-\kappa w)$ and, in the second case, $c=2$.

The above proposition contains the results given in [13], where the sequence $\left\{v_{i}\right\}$ is assumed to be nonincreasing and $\kappa=1$. Moreover, in the case $\gamma=2$, here the ratios $d_{i}$ are allowed to become arbitrarily small, provided that their rate of decay is controlled by (3.7).

In the next section, we will combine Scheme 1 of Section 2 with linear multigrid methods where

$$
\begin{equation*}
X_{i}=X, \quad \pi_{i}=\rho_{i}=I, \quad \Psi_{i}=I \tag{3.8}
\end{equation*}
$$

Precisely, given $B_{0}$, we will consider two pairs of methods of the type (3.1)-(3.3), defined by (3.8) and by particular choices of $\Gamma_{i}$ and $\Phi_{i}$ at the $i$-th step.

The first pair, we call linear multigrid methods 1 , is defined as follows (together with (3.8)):

Method 1a. $\quad \Gamma_{i}=I, \Phi_{i}=I-A_{i}$.
Method 1b. $\quad \Gamma_{i}=0, \Phi=I-A_{i}+A_{i-1}$.
Procedures of this type were proposed and discussed in [13]. In particular, they also yield the classical methods of Atkinson [2] and Brakhage [6] (cf. also Kelley [14]). For both Methods 1a and 1b, we get

$$
\begin{equation*}
R_{i}=\left(A_{i-1}\right)^{-1}\left[A_{i-1}-A_{i}\right]\left(I-A_{i}\right) \tag{3.9}
\end{equation*}
$$

The second pair of multigrid methods, we call here linear multigrid methods 2, are defined (together with (3.8)) by:
Method 2a. $\quad \Gamma_{i}=I, \Phi_{i}=I-A_{i-1}$.
Method 2b. $\Gamma_{i}=0, \Phi_{i}=I$.
Now, for both Methods 2a and 2b, we have

$$
\begin{equation*}
R_{i}=\left(A_{i-1}\right)^{-1}\left[A_{i-1}-A_{i}\right] \tag{3.10}
\end{equation*}
$$

In any case $\kappa=1$. For Methods 1a and 2a we can take in (3.5) $w=\operatorname{Sup}_{i}\left\|I-A_{i}\right\|$ and for Methods 1 b and 2 b , we take $w=1$. Accordingly (cf. Proposition 3.1), we can choose for Methods 1a and 2a even the value $\gamma=1$, provided that $w$ is sufficiently small, while for Methods 1 b and 2 b , we must take $\gamma=2$.
It is easy to verify that, for each case, the corresponding adjoint operators $B_{i}^{*}$, for $i=1,2, \ldots$, are defined through the following multigrid scheme:

Starting from $B_{0}^{*}$, for $i=1,2, \ldots$, set

$$
\begin{aligned}
S_{i-1}^{*} & =I-A_{i-1}^{*} B_{i-1}^{*} \\
Q_{i-1}^{*} & =\Sigma_{0 \leq j \leq \gamma} B_{i-1}^{*}\left(S_{i-1}^{*}\right)^{j} \\
B_{i}^{*} & =\Gamma_{i}^{*}+\Phi_{i}^{*} Q_{i-1}^{*}
\end{aligned}
$$

where now $I$ is the identity on $X^{*}$.
4. Nonlinear multilevel procedures. In this section we will examine some procedures obtained combining Scheme 1 with the two pairs of linear multigrid methods described in the previous section. The use of such procedures will be motivated in the light of the following result which is an immediate consequence of Theorem 2.1 (to which we refer for the notation).

Proposition 4.1. Let $\left\{\Delta_{k}\right\}, k=0,1, \ldots$, be a nonincreasing null sequence of positive numbers such that, for some $0<d<1$,

$$
\Delta_{k+1} \geq d \Delta_{k}, \quad \text { for every } k=0,1, \ldots
$$

Assume that there exist constants $C_{0}$ and $C_{1}$ such that, for every $k=0,1, \ldots$,

$$
(\Delta A)^{k}+(\Delta F)^{k} \leq C_{0} \Delta_{k+1}
$$

and

$$
\vartheta^{k} \leq C_{1} \Delta_{k}
$$

If $\Delta_{0}$ and $e_{\mathbf{x}}^{0}$ are sufficiently small, then Scheme 1 is well defined and there exists a null sequence $\left\{q_{k}\right\}, k=0,1, \ldots$, such that, for every $k$ we have

$$
\begin{equation*}
e_{\mathbf{x}}^{k+1} \leq c_{\mathbf{x}}\left(C_{0}+q_{k}\right) \Delta_{k+1} . \tag{4.1}
\end{equation*}
$$

Proof. Proceed by induction, using estimates (2.6) and (2.7).
Let us return to the equation

$$
Q_{U} \mathbf{x}-K(\mathbf{x})=0
$$

In this case the operator $A$, defined in (2.2), is given by

$$
\left.A(\mathbf{x}) \mathbf{v}=Q_{U} \mathbf{v}-K^{\prime}(\mathbf{x}) \mathbf{v}, M Q_{Z} \mathbf{v}\right), \quad \text { for } \mathbf{v} \in X
$$

Here we propose two different approaches for constructing the approximating operators involved in Scheme 1. The first approach is based only on the use of collectively compact (c.c. for brevity) approximations of $K$; the second one also employs projection methods.
a) Approximations through c.c. operators. Let $\left\{K_{n}\right\}$ be a sequence of completely continuous nonlinear operators from $D$ into $U$, which satisfies the following assumptions:
( $\left.\mathrm{i}_{1}\right)\left\{K_{n}\right\}$ is a c.c. family on $D$;
( $\mathrm{i}_{2}$ ) $K_{n}$ is pointwise convergent to $K$ on $D$, i.e., for each $\mathbf{x} \in D$,

$$
K_{n}(\mathbf{x}) \rightarrow K(\mathbf{x}), \quad \text { as } n \rightarrow \infty
$$

( $\mathrm{i}_{3}$ ) for every $n, K_{n}$ is of class $C^{\nu}, \nu \geq 3$, on $S\left(\mathbf{x}^{s}, \delta^{*}\right)$, and the derivatives $K_{n}^{\prime}(\mathbf{x}), K_{n}^{\prime \prime}(\mathbf{x})$ and $K_{n}^{\prime \prime \prime}(\mathbf{x})$ are uniformly bounded with respect to $\mathbf{x}$ and $n$.

Then, we consider the sequence of operators $\left\{A^{n}\right\}, A^{n}: X \rightarrow X$, defined by
(4.2) $A^{n}(\mathbf{x}) \mathbf{v}=\left(Q_{U} \mathbf{v}-K_{n}^{\prime}(\mathbf{x}) \mathbf{v}, M Q_{Z} \mathbf{v}\right), \quad$ for $\mathbf{v} \in X, \mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$.

As well-known (cf. [1, 2]) assumptions ( $\left.\mathrm{i}_{1}\right)-\left(\mathrm{i}_{3}\right)$ ensure that:
( $\mathrm{j}_{1}$ ) $\left\{I-A^{n}\right\}$ is a c.c. family of linear operators on $S\left(\mathbf{x}^{s}, \delta^{*}\right)$;
$\left(\mathrm{j}_{2}\right) \quad A^{n}(\mathbf{x}) \mathbf{v} \rightarrow A(\mathbf{x}) \mathbf{v}$, for every $\mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$ and $\mathbf{v} \in X$;
$\left(\mathrm{j}_{3}\right)$ there is a constant $C_{A}$ such that, for every $n$ and for every $\mathbf{u}, \mathbf{v} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$,

$$
\left\|A^{n}(\mathbf{u})-A^{n}(\mathbf{v})\right\|_{X \rightarrow X} \leq C_{A}|\mathbf{u}-\mathbf{v}|_{X}
$$

( $\mathrm{j}_{4}$ ) the sequence $\left\{A^{n}(\mathbf{x})\right\}$ is uniformly bounded, with respect to $\mathbf{x}$ and $n$;
( $\mathrm{j}_{5}$ ) for any compact operator $T: X \rightarrow X, \lim _{n \rightarrow \infty} \|(A(\mathbf{x})-$ $\left.A^{n}(\mathbf{x})\right) T \|_{X \rightarrow X}=0$, for each $\mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$;
$\left(\mathrm{j}_{6}\right)$ for all sufficiently large $n$ and for every $\mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right), A^{n}(\mathbf{x})$ is invertible and the inverses are uniformly bounded with respect to $\mathbf{x}$ and $n$. Then, we consider a suitable subsequence of $\left\{K_{n}\right\}$, we relabel as $\left\{K_{k}\right\}$, for $k=0,1, \ldots$. Accordingly, in Scheme 1, we take

$$
F^{k}(\mathbf{x})=Q_{U} \mathbf{x}-K_{k}(\mathbf{x})
$$

Assumptions $\left(\mathrm{k}_{1}\right)-\left(\mathrm{k}_{4}\right)$ of Theorem 2.1 are clearly fulfilled.
Now we discuss the construction of $B^{k}\left(\right.$ and $\left.\left(B^{k}\right)^{*}\right)$ in Scheme 1, by the linear multigrid methods of Section 3. Referring to (4.2), a first procedure could be defined by taking

$$
\begin{equation*}
B_{0}=A^{0}\left(\mathbf{x}^{0}\right)^{-1} \tag{4.3}
\end{equation*}
$$

and, for each $k=1,2, \ldots$,

$$
\begin{equation*}
A_{i}=A^{i}\left(\mathbf{x}^{k}\right), \quad \text { for } i=0,1, \ldots, k \tag{4.4}
\end{equation*}
$$

Then we employ one of the two linear multigrid schemes 1 considered in Section 3. We set, in Scheme $1, B^{k}=B_{k}$, for $k=0,1, \ldots$ Accordingly, for each $k=1,2, \ldots$, from (3.9) we have

$$
\begin{equation*}
R_{i}=\left[A^{i-1}\left(\mathbf{x}^{k}\right)\right]^{-1}\left[A^{i-1}\left(\mathbf{x}^{k}\right)-A^{i}\left(\mathbf{x}^{k}\right)\right]\left(I-A^{i}\left(\mathbf{x}^{k}\right)\right), \quad \text { for } i=1, \ldots, k, \tag{4.5}
\end{equation*}
$$

and

$$
\begin{equation*}
S_{0}=\left[A^{0}\left(\mathbf{x}^{0}\right)\right]^{-1}\left[A^{0}\left(\mathbf{x}^{0}\right)-A^{0}\left(\mathbf{x}^{k}\right)\right] . \tag{4.6}
\end{equation*}
$$

A second procedure could be defined taking $B_{0}$ as in (4.3) and, for each $k=1,2, \ldots$,

$$
\begin{equation*}
A_{i}=A^{i}\left(\mathrm{x}^{i}\right), \quad \text { for } i=0,1, \ldots, k \tag{4.7}
\end{equation*}
$$

Again, we use one of the linear multigrid schemes 1 considered in Section 3, setting $B^{k}=B_{k}$, for each $k$. In this case, we have, for $k=1,2, \ldots$,

$$
\begin{equation*}
R_{k}=\left[A^{k-1}\left(\mathrm{x}^{k-1}\right)\right]^{-1}\left[A^{k-1}\left(\mathrm{x}^{k-1}\right)-A^{k}\left(\mathrm{x}^{k}\right)\right]\left(I-A^{k}\left(\mathrm{x}^{k}\right)\right) \tag{4.8}
\end{equation*}
$$

and

$$
\begin{equation*}
S_{0}=0 \tag{4.9}
\end{equation*}
$$

Clearly, using (4.7) instead of (4.4), one has the advantage of constructing each $A_{i}$ once and for all.

Now we show how Proposition 4.1 can be applied to the cases considered above. In the following $c_{j}$, for $j=1,2,3,4,5$, will be constants independent of the index $k$.

Assume ( $\mathrm{i}_{1}$ )- $\mathrm{i}_{3}$ ) and $\left(\mathrm{j}_{1}\right)-\left(\mathrm{j}_{6}\right)$, and suppose that

$$
\begin{gathered}
\|\left.\left[K^{\prime}\left(\mathbf{x}^{s}\right)-K_{k-1}^{\prime}\left(\mathbf{x}^{s}\right)\right]\left[I-A^{k}\left(\mathbf{x}^{s}\right)\right]\right|_{X \rightarrow U} \leq c_{1} \Delta_{k} \\
\left|K\left(\mathbf{x}^{s}\right)-K_{k}\left(\mathbf{x}^{s}\right)\right|_{U}+\operatorname{Max}_{1 \leq i \leq p}\left|\left[K^{\prime}\left(\mathbf{x}^{s}\right)-K_{k}^{\prime}\left(\mathbf{x}^{s}\right)\right] \mathbf{y}_{i}^{s}\right|_{U} \leq c_{2} \Delta_{k+1}
\end{gathered}
$$

for every $k=1,2, \ldots$, where the sequence $\left\{\Delta_{k}\right\}$ is chosen as in Proposition 4.1. Then, for both choices (4.3)-(4.4) and (4.3)-(4.7), using (4.5)-(4.6), or (4.8)-(4.9), and Proposition 3.1, one can easily prove by induction that (4.1) holds provided that $\Delta_{0}$ and $e_{\mathbf{x}}^{0}$ are sufficiently small.
b) Approximations through c.c. operators and projections. Let $\left\{\mathbf{P}_{n}\right\}$ be a sequence of linear projections from $X$ into itself which is pointwise convergent to $I$. Let the sequence $\left\{K_{n}\right\}$ be given as before. As is wellknown (cf. [4]), the family of operators $\left\{K_{n}\left(\mathbf{P}_{n} \cdot\right)\right\}$ satisfies (i $\left.\mathrm{i}_{1}\right)-\left(\mathrm{i}_{3}\right)$. Assume that $\left\{K_{k}\right\}$ and $\left\{\mathbf{P}_{k}\right\}$ are (relabelled) subsequences of $\left\{K_{n}\right\}$ and $\left\{\mathbf{P}_{n}\right\}$, respectively, such that:
$\left(\mathrm{m}_{1}\right)$ for every $\mathbf{x} \in S\left(\mathbf{x}^{s}, \delta^{*}\right)$, the operators $A^{k}(\mathbf{x})$ defined by

$$
\begin{equation*}
A^{k}(\mathbf{x}) \mathbf{v}=\left(Q_{U} \mathbf{v}-K_{k}^{\prime}\left(\mathbf{P}_{k} \mathbf{x}\right) \mathbf{P}_{k} \mathbf{v}, M Q_{Z} \mathbf{v}\right), \quad \text { for } \mathbf{v} \in X \tag{4.10}
\end{equation*}
$$

are invertible and their inverses are uniformly bounded with respect to $\mathbf{x}$ and $k$.

Then, in Scheme 1 we take

$$
F^{k}(\mathbf{x})=Q_{U} \mathbf{x}-K_{k}\left(\mathbf{P}_{k} \mathbf{x}\right)
$$

Moreover, referring to the operators defined in (4.10), we can use one of the linear multigrid schemes 1 , taking $B_{0}$ as in (4.3), $A_{i}$ as in (4.7) and $B^{k}=B_{k}$. Accordingly, (4.8) and (4.9) follow.

Let $\left\{\Delta_{k}\right\}$ be a sequence as in Proposition 4.1 and suppose that, for every $k=1,2, \ldots$,

$$
\begin{equation*}
\left\|\left(K^{\prime}\left(\mathbf{x}^{s}\right)-K_{k-1}^{\prime}\left(\mathbf{P}_{k-1} \mathbf{x}^{s}\right) \mathbf{P}_{k-1}\right)\left[I-A^{k}\left(\mathbf{x}^{s}\right)\right]\right\|_{X \rightarrow U} \leq c_{3} \Delta_{k} \tag{4.11}
\end{equation*}
$$

$\left|K\left(\mathbf{x}^{s}\right)-K_{k}\left(\mathbf{P}_{k} \mathbf{x}^{s}\right)\right|_{U}+\operatorname{Max}_{1 \leq i \leq p}\left|\left[K^{\prime}\left(\mathbf{x}^{s}\right)-K_{k}^{\prime}\left(\mathbf{P}_{k} \mathbf{x}^{s}\right) \mathbf{P}_{k}\right] \mathbf{y}_{i}^{s}\right|_{U} \leq c_{4} \Delta_{k+1}$.
Then, by assumption $\left(\mathrm{m}_{1}\right)$ and invoking again Propositions 3.1 and 4.1, the usual inductive argument shows that relation (4.1) holds for every $k$, provided that $\Delta_{0}$ and $e_{x}^{0}$ are sufficiently small.

Let us consider the particular case where $K_{k}=K$, for every $k$, and therefore

$$
F^{k}(\mathbf{x})=Q_{U} \mathbf{x}-K\left(\mathbf{P}_{k} \mathbf{x}\right)
$$

and

$$
A^{k}(\mathbf{x}) \mathbf{v}=\left(Q_{U} \mathbf{v}-K^{\prime}\left(\mathbf{P}_{k} \mathbf{x}\right) \mathbf{P}_{k} \mathbf{v}, M Q_{Z} \mathbf{v}\right), \quad \text { for } \mathbf{v} \in X
$$

If

$$
\begin{equation*}
\|\left(K^{\prime}\left(\mathbf{x}^{s}\right)\left(I-\mathbf{P}_{k}\right) \|_{X \rightarrow U} \rightarrow 0\right. \tag{4.12}
\end{equation*}
$$

we could also use the linear multigrid schemes 2 of Section 3, taking, as before, $B^{0}=B_{0}=A^{0}\left(\mathbf{x}^{0}\right)^{-1}$, for each $k, A_{i}=A^{i}\left(\mathbf{x}^{i}\right)$, for $i=0,1, \ldots, k$ and $B^{k}=B_{k}$. In this case, from (3.10), we get

$$
R_{k}=\left[A^{k-1}\left(\mathrm{x}^{k-1}\right)\right]^{-1}\left[A^{k-1}\left(\mathrm{x}^{k-1}\right)-A^{k}\left(\mathrm{x}^{k}\right)\right]
$$

and $S_{0}=0$. Then, assuming that, for every $k$, relation (4.11) (with $K_{k}=K$ ) holds and that

$$
\left\|\left(K^{\prime}\left(\mathbf{x}^{s}\right)-K^{\prime}\left(\mathbf{P}_{k-1} \mathbf{x}^{s}\right) \mathbf{P}_{k-1}\right)\right\|_{X \rightarrow U} \leq c_{5} \Delta_{k}
$$

the conclusions of Proposition 4.1 easily follow.
We notice that, for all the cases considered here, it was understood that, in Scheme 1, the adjoint $B^{k *}$ is obtained by the same multigrid method which yields $B^{k}$, as shown at the end of Section 3.
There are several results, on the approximation of integral operators, which can be used for checking the validity of the various assumptions
above considered as well as for estimating the values of $\left\{\Delta_{k}\right\}$. Among the more recent papers, we refer to $[\mathbf{2}, \mathbf{3}, \mathbf{4}, \mathbf{5}, \mathbf{8}]$.

Remarks. The more natural way for constructing the sequence of projections $\left\{\mathbf{P}_{n}\right\}$ in $X$ is that of taking a sequence $\left\{P_{n}\right\}$ of linear projections in $U$, and then defining $\mathbf{P}_{n}$ as $\mathbf{P}_{n} \mathbf{v}=\left(P_{n} Q_{U} \mathbf{v}, \mathbf{v}_{p}\right)$. Here, for conciseness, we have expressed any $\mathbf{v} \in X$ in the form $\mathbf{v}=\left(Q_{U} \mathbf{v}, \mathbf{v}_{p}\right)$, with $\mathbf{v}_{p} \in \mathbf{R}^{p}$. In order to illustrate how in this case the application of $B_{0}=A^{0}\left(\mathbf{x}^{0}\right)^{-1}$, with $A^{0}$ given by (4.10), can be carried out, we assume, for simplicity, that the operator $M Q_{Z}$ has the form $M Q_{Z} \mathbf{v}=\left(M_{1} Q_{U} \mathbf{v}, M_{2} \mathbf{v}_{p}\right)$, with $M_{1} \in U^{*}, M_{2} \in L\left(\mathbf{R}^{p}, \mathbf{R}^{p-1}\right)$. This is the case usually occurring in practice. Accordingly, an application of $B_{0}$ consists of solving a system of the type

$$
\begin{gathered}
Q_{U} \mathbf{v}-K_{0}^{\prime}\left(\mathbf{P}_{0} \mathbf{x}^{0}\right) \mathbf{P}_{0} \mathbf{v}=f_{1} \\
M_{1} Q_{U} \mathbf{v}=f_{2} \\
M_{2} \mathbf{v}_{p}=f_{3}
\end{gathered}
$$

with $f_{1} \in U, f_{2} \in \mathbf{R}$ and $f_{3} \in \mathbf{R}^{p-1}$. This is equivalent to solve, at first, the projected system

$$
\begin{gather*}
P_{0}\left(Q_{U} \mathbf{v}-K_{0}^{\prime}\left(\mathbf{P}_{0} \mathbf{x}^{0}\right) \mathbf{P}_{0} \mathbf{v}-f_{1}\right)=0, \\
M_{1}\left(K_{0}^{\prime}\left(\mathbf{P}_{0} \mathbf{x}^{0}\right) \mathbf{P}_{0} \mathbf{v}+f_{1}\right)=f_{2},  \tag{4.13}\\
M_{2} \mathbf{v}_{p}=f_{3},
\end{gather*}
$$

with respect to the unknowns $\mathbf{P}_{0} \mathbf{v}$ and $\mathbf{v}_{p}$ and then to get $Q_{U} \mathbf{v}$ as the iterated solution of (4.13) (in the sense of Sloan [19]), namely,

$$
Q_{U} \mathbf{v}=K_{0}^{\prime}\left(\mathbf{P}_{0} \mathbf{x}^{0}\right) \mathbf{P}_{0} \mathbf{v}+f_{1}
$$

Of course, the coefficient matrix of the above linear system will be factorized once and for all.

In the particular case when equation (2.1) has the classical Hammerstein form

$$
\begin{equation*}
Q_{U} \mathbf{w}-K G(\mathbf{w})=0 \tag{4.14}
\end{equation*}
$$

where now $K$ is a linear compact operator into $U$ and $G$ is a nonlinear mapping from $X$ into $U$, for an efficient use of projection methods, it is convenient to consider, instead of (4.14), the equation

$$
Q_{U} \mathbf{x}-G\left(K Q_{U} \mathbf{x}, \mathbf{x}_{p}\right)=0
$$

obtained through the change of variable $\mathbf{x}=\left(G(\mathbf{w}), \mathbf{w}_{p}\right)$ that is $\mathbf{w}=\left(K Q_{U} \mathbf{x}, \mathbf{x}_{p}\right)$. For more details on this point, we refer to $[\mathbf{1 5}$, 16].

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