# SUPERCONVERGENCE OF MODIFIED PROJECTION METHOD FOR INTEGRAL EQUATIONS OF THE SECOND KIND 

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#### Abstract

Recently, Nair proposed a modification for the projection (Galerkin) method for the approximate solution of second kind integral equations and established conditions under which it was superconvergent (relative to the corresponding projection approximation). Its advantage over other superconvergent methods for such equations is its simplicity of implementation in requiring virtually no more computation than that implicit in the determination of the corresponding projection approximation. In this paper, we propose two variants of Nair's modification and establish conditions under which they are superconvergent.


1. Introduction. Using operator notation, second kind Fredholm integral equations take the form

$$
\begin{equation*}
x-K x=y \tag{1.1}
\end{equation*}
$$

where the operator $K$ is defined, in terms of its (known) kernel $k(s, t)$, by

$$
\begin{equation*}
(K v)(s)=\int_{a}^{b} k(s, t) v(t) d t, \quad a \leq s \leq b \tag{1.2}
\end{equation*}
$$

It is assumed that $K$ is defined on a suitable Banach space $X$, that $K$ is compact as a mapping from $X$ into $X$, and that 1 is not an eigenvalue of $K$. Then, by Fredholm theory, (1.1) has a unique solution $x$ in $X$ for every $y$ in $X$.

Approximate solutions are sought using a sequence of finite dimensional spaces $X_{n}$ which, in practice, will be generated by polynomials or piecewise polynomials. Let $\pi_{n}$ be a projection of $X$ onto $X_{n}$. Then, assuming $\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}$ exists, the projection (Galerkin) approximation $x_{n}^{\mathrm{P}}$ of $x$ is defined by the equation

$$
x_{n}^{P}-\boldsymbol{\pi}_{n} K x_{n}^{P}=\boldsymbol{\pi}_{n} y .
$$

The iterated projection (or Sloan) approximation $\tilde{x}_{n}^{\mathrm{P}}=x_{n}^{S}$ and the Kantorovich approximation $x_{n}^{K}$ are defined by

$$
x_{n}^{\mathrm{S}}=y+K x_{n}^{\mathrm{P}} \quad \text { and } \quad x_{n}^{K}=y+z_{n}^{\mathrm{P}}
$$

respectively, where $z_{n}^{\mathrm{P}}$ is the projection solution of the equation $z-$ $K z=K y$; i.e., $z_{n}^{\mathrm{P}}$ satisfies

$$
z_{n}^{\mathrm{P}}-\boldsymbol{\pi}_{n} K z_{n}^{\mathrm{P}}=\boldsymbol{\pi}_{n} K y
$$

Recently, Nair [7] introduced the modified projection approximation $x_{n}^{M}$, defined by

$$
x_{n}^{M}=x_{n}^{\mathrm{P}}+\left(I-\boldsymbol{\pi}_{n}\right) y,
$$

and showed that it satisfies the relation

$$
\begin{equation*}
x-x_{n}^{M}=\boldsymbol{\pi}_{n}\left(x-x_{n}^{\mathrm{S}}\right)+\left(I-\boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{K}\right) \tag{1.3}
\end{equation*}
$$

Thus, if $\boldsymbol{\pi}_{n}$ is uniformly bounded, then (Theorem 2.1 in $[\mathbf{7}]$ )

$$
\begin{equation*}
\left\|x-x_{n}^{M}\right\| \leq c \max \left\{\left\|x-x_{n}^{\mathrm{S}}\right\|,\left\|x-x_{n}^{K}\right\|\right\} \tag{1.4}
\end{equation*}
$$

where $c$ denotes (here and below) a generic constant.
In the standard superconvergence arguments (e.g., Sloan and Thomeé $[\mathbf{1 3}])$, it is assumed that the projections are evaluated exactly. When numerical approximations are used, the proof of such results usually becomes technically more difficult. The advantage of Nair's [7] modification is that it is performed exactly using the information derived in constructing the projection approximation $x_{n}^{P}$. However, this practical advantage is achieved through a partial loss of full superconvergence, which is the topic of Nair $[\mathbf{7}]$ as well as of this paper.

It is known that, under certain smoothness assumptions on $K$ and $y$, the Sloan approximation $x_{n}^{\mathrm{S}}$ converges to the exact solution $x$ faster than $x_{n}^{\mathrm{P}}$; i.e., $x_{n}^{\mathrm{S}}$ exhibits superconvergence (see Chandler [1], Chatelin and Lebbar [2], Graham et al. [6], Joe [5], and Sloan [12]). In addition, the speed of convergence of $x_{n}^{S}$ also depends on the form of the projection method used (e.g., $x_{n}^{S}$ does not always have the better convergence when collocation methods are applied).

If $y$ has less regularity, then it is preferable to use the Kantorovich approximation $x_{n}^{K}$ rather than $x_{n}^{P}$, for in the Kantorovich method we
approximate $z=K x$ rather than $x$ itself (see Schock [8] and Sloan [11]). Moreover, the Kantorovich approximation, in general, has better global (uniform) convergence behaviour than the projection and Sloan approximations (see Schock [9]). In the light of these comments, $x_{n}^{M}$ is likely to have better convergence rates than $x_{n}^{P}$ under conditions favorable to $x_{n}^{\mathrm{S}}$ and $x_{n}^{K}$ and, in addition, $x_{n}^{M}$ requires (virtually) no more computation than $x_{n}^{\mathrm{P}}$ (see computational results in $[7]$ ). It should be noted, however, that $x_{n}^{M}$ can never be better than $x_{n}^{K}$. This is an immediate consequence of the result that

$$
\left(I-\boldsymbol{\pi}_{n} K\right)\left(x-x_{n}^{K}\right)=\left(I-\boldsymbol{\pi}_{n}\right) K x=\left(I-\boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{K}\right),
$$

since, using (1.3), it yields

$$
\begin{equation*}
\left(I-\boldsymbol{\pi}_{n} K\right)\left(x-x_{n}^{K}\right)=\left(I-\boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{M}\right) \tag{1.5}
\end{equation*}
$$

Thus, whenever $\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}$ and $\boldsymbol{\pi}_{n}$ are uniformly bounded, we obtain, using (1.3), (1.4) and (1.5), that

$$
c\left\|x-x_{n}^{K}\right\| \leq\left\|x-x_{n}^{M}\right\| \leq c^{\prime} \max \left\{\left\|x-x_{n}^{\mathrm{S}}\right\|,\left\|x-x_{n}^{K}\right\|\right\}
$$

where $c^{\prime}$ is a constant satisfying $\left\|\left(I-\pi_{n} K\right)^{-1}\left(I-\pi_{n}\right)\right\| \leq 1 / c^{\prime}$. Thus, whenever $x_{n}^{S}$ has a better order of convergence than $x_{n}^{K}$, the order of convergence of $x_{n}^{M}$ is exactly that of $x_{n}^{K}$. An example given by Sloan [11] can be used to illustrate this point. However, one iteration of $x_{n}^{K}$ (i.e., the iterated Kantorovich approximation) $\tilde{x}_{n}^{K}=y+K x_{n}^{K}$ is best among all the former approximations (see Sloan [11, 12]). In [13], Sloan and Thomeé have shown that, in situations where the smoothing of powers of $K$ is not optimal, partial superconvergence of $x_{n}^{\mathrm{S}}$ can be guaranteed under appropriate circumstances. They have also shown that the higher iterates of $x_{n}^{P}$, namely

$$
x_{n}^{(k)}=y+K x_{n}^{(k-1)}, \quad x_{n}^{(0)}=x_{n}^{\mathrm{P}}, \quad k=1,2, \ldots
$$

have improving rates of convergence provided the powers of $K$ have appropriate smoothing properties.

In this paper, we introduce two new approximations which are variants of the modified approximation of Nair [7] and establish conditions for their superconvergence in terms of that for

$$
\tilde{x}_{n}^{K}=y+K x_{n}^{K} \quad \text { and } \quad \tilde{x}_{n}^{\mathrm{S}}=y+K x_{n}^{\mathrm{S}}
$$

The two new approximations are
(1) $\hat{x}_{n}^{M}=y+z_{n}^{M}$,
(2) $\tilde{x}_{n}^{M}=y+K x_{n}^{M}$,
where

$$
z_{n}^{M}=z_{n}^{\mathrm{P}}+\left(I-\boldsymbol{\pi}_{n}\right) K y
$$

is the modified projection solution of $z-K z=K y$. The approximations $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$ can also be written as

$$
\begin{align*}
& \hat{x}_{n}^{M}=\left(y+z_{n}^{\mathrm{P}}\right)+\left(I-\boldsymbol{\pi}_{n}\right) K y=x_{n}^{K}+\left(I-\boldsymbol{\pi}_{n}\right) K y  \tag{1.6}\\
& \tilde{x}_{n}^{M}=y+K\left(x_{n}^{\mathrm{P}}+\left(I-\boldsymbol{\pi}_{n}\right) y\right)=x_{n}^{\mathrm{S}}+K\left(I-\boldsymbol{\pi}_{n}\right) y \tag{1.7}
\end{align*}
$$

In fact, both of these approximations can be viewed as particular cases of

$$
w_{n}=x_{n}+\left(K-K_{n}\right) y
$$

where $x_{n}$ is the solution of

$$
\begin{equation*}
x_{n}-K_{n} x_{n}=y \tag{1.8}
\end{equation*}
$$

with $K_{n}=\boldsymbol{\pi}_{n} K$, for $\hat{x}_{n}^{M}$, and $K_{n}=K \boldsymbol{\pi}_{n}$ for $\tilde{x}_{n}^{M}$. Because of the form of expressions (1.6) and (1.7) for $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$, we shall refer to these approximations as the modified Kantorovich approximation and the modified Sloan approximation, respectively.
In the next section, we obtain error estimates for $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$ in an abstract setting, while, in the final section, we examine their application to the integral equation (1.1).
2. Error bounds for $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$. We recall the basic assumption which we will invoke in the remainder of this paper. Let $K$ denote a compact operator on a Banach space $X$ and $\boldsymbol{\pi}_{n}$ a projection of $X$ onto $X_{n}$. In addition, it is assumed that
(i) 1 is not an eigenvalue of $K$, and
(ii) for sufficiently large $n$, the $\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}$ exist and are uniformly bounded.
Under these assumptions, each of the equations

$$
u-K u=v
$$

and

$$
\begin{equation*}
u-\pi_{n} K u=v \tag{2.1}
\end{equation*}
$$

has a unique solution in $X$ for any given $v$ in $X$. Thus, $x_{n}^{P}$ and $z_{n}^{\mathrm{P}}$ exist uniquely as solutions of (2.1), with $v=\boldsymbol{\pi}_{n} y$ and $v=\boldsymbol{\pi}_{n} K y$, respectively. Thereby, the approximations

$$
\hat{x}_{n}^{M}=y+z_{n}^{M}, \quad z_{n}^{M}=z_{n}^{\mathrm{P}}+\left(I-\boldsymbol{\pi}_{n}\right) K y
$$

and

$$
\tilde{x}_{n}^{M}=y+K x_{n}^{M}, \quad x_{n}^{M}=x_{n}^{\mathrm{P}}+\left(I-\boldsymbol{\pi}_{n}\right) y
$$

are well defined, for sufficiently large $n$.
The error bounds to be given in Theorem 2.1 involve the iterated Kantorovich approximation

$$
\tilde{x}_{n}^{K}=y+K x_{n}^{K}, \quad x_{n}^{K}=y+z_{n}^{\mathrm{P}}
$$

and the iterated Sloan approximation

$$
\tilde{x}_{n}^{\mathrm{S}}=y+K x_{n}^{\mathrm{S}}, \quad x_{n}^{\mathrm{S}}=y+K x_{n}^{\mathrm{P}}
$$

The following relations are easily derived from their definitions:

$$
\begin{align*}
x-\tilde{x}_{n}^{K} & =K\left(x-x_{n}^{K}\right), \\
x-\tilde{x}_{n}^{\mathrm{S}} & =K\left(x-x_{n}^{s}\right)=K^{2}\left(x-x_{n}^{\mathrm{P}}\right), \\
K\left(I-\boldsymbol{\pi}_{n}\right) x & =\left(I-K \boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{S}\right), \\
\left(I-\boldsymbol{\pi}_{n}\right) K x & =\left(I-\boldsymbol{\pi}_{n} K\right)\left(x-x_{n}^{K}\right),  \tag{2.2}\\
K\left(I-\boldsymbol{\pi}_{n}\right) K x & =\left(I-K \boldsymbol{\pi}_{n}\right)\left(x-\tilde{x}_{n}^{K}\right) . \tag{2.3}
\end{align*}
$$

Theorem 2.1.
(a) If the sequence $\boldsymbol{\pi}_{n}$ is uniformly bounded, then
$c_{1}\left\|\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x\right\| \leq\left\|x-\hat{x}_{n}^{M}\right\| \leq c_{2} \max \left\{\left\|x-\tilde{x}_{n}^{K}\right\|,\left\|\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x\right\|\right\}$
for some constants $c_{1}$ and $c_{2}$.
(b) If the sequence ( $K \boldsymbol{\pi}_{n}$ ) is uniformly bounded, then

$$
\left\|x-\tilde{x}_{n}^{M}\right\| \leq c \max \left\{\left\|K\left(I-\pi_{n}\right)\left(x-x_{n}^{\mathrm{S}}\right)\right\|,\left\|x-\tilde{x}_{n}^{K}\right\|,\left\|x-\tilde{x}_{n}^{\mathrm{S}}\right\|\right\}
$$

for some constant $c$.

Proof. Let $K_{n}$ represent $\boldsymbol{\pi}_{n} K$ or $K \boldsymbol{\pi}_{n}$, and let $x_{n}$ be the unique solution of (1.8). Then

$$
\left(I-K_{n}\right)\left(x-x_{n}\right)=\left(K-K_{n}\right) x
$$

Denoting $w_{n}=x_{n}+\left(K-K_{n}\right) y$, we obtain

$$
\begin{align*}
x-w_{n} & =\left(x-x_{n}\right)-\left(K-K_{n}\right) y \\
& =\left(x-x_{n}\right)-\left(K-K_{n}\right)(x-K x)  \tag{2.4}\\
& =\left(x-x_{n}\right)-\left(I-K_{n}\right)\left(x-x_{n}\right)+\left(K-K_{n}\right) K x \\
& =K_{n}\left(x-x_{n}\right)+\left(K-K_{n}\right) K x
\end{align*}
$$

Now, taking $K_{n}=\pi_{n} K$, we obtain

$$
x-\hat{x}_{n}^{M}=\boldsymbol{\pi}_{n} K\left(x-x_{n}^{K}\right)+\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x
$$

In addition,

$$
\left(I-\pi_{n}\right)\left(x-\hat{x}_{n}^{M}\right)=\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x
$$

Together these relations yield

$$
\left\|x-\hat{x}_{n}^{M}\right\| \leq\left\|\boldsymbol{\pi}_{n}\right\|\left\|x-\tilde{x}_{n}^{K}\right\|+\left\|\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x\right\|
$$

and

$$
\left\|\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x\right\| \leq\left(1+\left\|\boldsymbol{\pi}_{n}\right\|\right)\left\|x-\hat{x}_{n}^{M}\right\|
$$

Thus, the result in (a) follows by taking $c_{1}$ and $c_{2}$ such that

$$
1+\left\|\boldsymbol{\pi}_{n}\right\| \leq 1 / c_{1} \quad \text { and } \quad c_{2}=1 / c_{1}
$$

Now, taking $K_{n}=K \boldsymbol{\pi}_{n}$ in (1.8), we obtain, using (2.3),

$$
\begin{aligned}
x-\hat{x}_{n}^{M} & =K \boldsymbol{\pi}_{n}\left(x-x_{n}^{\mathrm{S}}\right)+K\left(I-\boldsymbol{\pi}_{n}\right) K_{x} \\
& =K\left(x-x_{n}^{s}\right)-\left(K-K \boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{\mathrm{S}}\right)+K\left(I-\boldsymbol{\pi}_{n}\right) K x \\
& =\left(x-\tilde{x}_{n}^{\mathrm{S}}\right)-K\left(I-\boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{\mathrm{S}}\right)+\left(I-K \boldsymbol{\pi}_{n}\right)\left(x-\tilde{x}_{n}^{K}\right)
\end{aligned}
$$

Thus, the result in (b) follows on taking a $c$ satisfying

$$
3+\left\|K \boldsymbol{\pi}_{n}\right\| \leq c . \square
$$

Remark 2.2. It can be seen that the order of convergence of $\|(I-$ $\left.\pi_{n}\right) K^{2} x \|$ to zero is the same as that of $\left\|x-\hat{x}_{n}^{K}\right\|$, where $\hat{x}_{n}^{K}$ denotes the second order Kantorovich approximation considered by Schock [10]:

$$
\hat{x}_{n}^{K}=y+K y+\hat{z}_{n}, \quad \hat{z}_{n}-\pi_{n} K \hat{z}_{n}=\pi_{n} K^{2} y
$$

Schock [10] has also characterized a class of operators $K$ for which $\hat{x}_{n}^{K}$ is a better approximation of $x$ than $x_{n}^{K}$.

We remark that the assumptions in Theorem 2.1 are satisfied if $\left\|\left(I-\boldsymbol{\pi}_{n}\right) u\right\| \rightarrow 0$ as $n \rightarrow \infty$, for every $u$ in $X$. However, this is a much stronger condition than is required. For example, in [14], Sloan and Burn considered the (Lagrange) interpolatory projection $\boldsymbol{\pi}_{n}$ with $n$ Tchebychev's points as the nodes in $[-1,1]$ and showed that $\left\|K\left(I-\pi_{n}\right) u\right\| \rightarrow 0$ as $n \rightarrow \infty$ for every $u$ in $X=C[-1,1]$, which implies that $\left(K \boldsymbol{\pi}_{n}\right)$ is uniformly bounded. In this case, $\boldsymbol{\pi}_{n}$ is not uniformly bounded. On the other hand, if we take $X_{n}$ to be the space of piecewise polynomials of degree $\leq r-1$ with $r$ Gaussian quadrature points placed on each subinterval as the nodes, then the interpolatory projections $\boldsymbol{\pi}_{n}$ onto $X_{n}$ are known to be uniformly bounded (see Graham et al. [6]). But, $\left\|\left(I-\boldsymbol{\pi}_{n}\right) u\right\| \rightarrow 0$ does not hold for all $u$ in $X=C[-1,1]$.

The key point is that the assumptions used to obtain the results in Theorem 2.1 are minimal. In fact, Theorem 2.1 has been structured to give the interrelationships between the errors associated with the different approximations without considering whether convergence is guaranteed.
3. Superconvergence of $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$. In this section, we apply the results of $\S 2$ to the integral equation (1.1). We now assume that the Banach space $X=L_{p}=L_{p}[a, b], 1 \leq p \leq \infty$, and take $X_{n}$ as the space $S_{r, n}$ of all piecewise polynomials of degree $\leq r-1$ defined on the mesh

$$
\Delta_{n}:=\left\{a=t_{0}<t_{1}<\cdots<t_{n}=b\right\}
$$

with no continuity requirement at the nodes. Let $\boldsymbol{\pi}_{n}$ be the projection of $L_{p}$ onto $S_{r, n}$ defined by

$$
\left\langle\boldsymbol{\pi}_{n} v, u_{n}\right\rangle=\left\langle v, u_{n}\right\rangle, \quad \text { for all } v \in L_{p}, \quad u_{n} \in S_{r, n}
$$

where $\langle.,$.$\rangle denotes the L_{2}$-inner product. We further assume that the sequence of meshes $\Delta_{n}$ is quasiuniform (i.e., $h_{i} \geq c h$, where $\left.h=\max _{1 \leq i \leq n} h_{i}\right), h_{i}=t_{i}-t_{i-1}$. As a consequence, $\boldsymbol{\pi}_{n}$ is uniformly bounded (see Werschulz [15, Lemma 3.1]). It also follows that

$$
\left\|\left(I-\boldsymbol{\pi}_{n}\right) v\right\|=\left\|\left(I-\boldsymbol{\pi}_{n}\right)\left(v-u_{n}\right)\right\|, \quad \text { for all } u_{n} \in S_{r, n}
$$

and hence

$$
\begin{equation*}
\left\|\left(I-\boldsymbol{\pi}_{n}\right) v\right\| \leq c \inf _{u_{n} \in S_{r, n}}\left\|v-u_{n}\right\| \tag{3.1}
\end{equation*}
$$

where $c \geq 1+\left\|\boldsymbol{\pi}_{n}\right\|$, for all $n$. Here $\|v\|$ denotes the $L_{p}$-norm of the function $v$.

For non-negative integers $l$, we let $W_{p}^{l}$ denote the Sobolev space of functions defined on $[a, b]$ and equipped with the norm

$$
\|v\|_{l, p}=\max _{0 \leq j \leq l}\left\|v^{(j)}\right\|_{p}
$$

where $v^{(j)}$ is the $j$-th distributional derivative of $v$. In order to describe the smoothing properties of $K$, we consider, as in [13], a class $\mathcal{K}_{p, m, l}$ of operators $\mathbf{K}$ which maps $L_{p}$ into $W_{p}^{m}$ and satisfies

$$
\|K v\|_{m, p} \leq c\|v\|_{-l, p}, \quad \vartheta \in L_{p}
$$

for some constant $c$. Here, the negative Sobolev norm $\|\cdot\|_{-l, p}$ is defined by

$$
\begin{equation*}
\|v\|_{-l, p}:=\sup _{u \in W_{q}^{l}} \frac{|\langle v, u\rangle|}{\|u\|_{l, q}}, \quad v \in L_{p} \tag{3.2}
\end{equation*}
$$

with $1 / p+1 / q=1$. We recall the following well- known result on the approximation properties of $S_{r, n}$ (see de Vore [3] or Demko [4]).

Proposition 3.1. If $\vartheta \in W_{p}^{l}$, then, for each $n \geq 0$,

$$
d_{p}\left(v, s_{r, n}\right)=\inf _{u_{n} \in S_{r, n}}\left\|v-u_{n}\right\| \leq c h^{\nu}\|v\|_{l, p}
$$

where $\nu=\min \{l, r\}, 1 \leq p \leq \infty$, and $c$ is a constant.

From this result, using the compactness of $K$ and uniform boundedness of the $\boldsymbol{\pi}_{n}$, one gets $\left\|\left(I-\boldsymbol{\pi}_{n}\right) K\right\| \rightarrow 0$ as $n \rightarrow \infty$, so that, for $n$ sufficiently large, $\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}$ exists that is uniformly bounded. Here, and in what follows, $\|A\|$ denotes the norm of the operator $A$ induced by the norm of the $L_{p}$ space. Unless otherwise specified, $\|v\|$ denotes the $L_{p}$-norm of the function $v$.

Throughout the remainder of this paper, we shall impose the following smoothness conditions on $K$ :
(1) $K^{*}$, the adjoint of $K$, maps $W_{p}^{l}$ into itself for every non-negative integer $l \leq r$, where $1 / p+1 / q=1$;
(2) $K$ is of class of $\mathcal{K}_{p, 0, l_{1}}$ and $\mathcal{K}_{p, m_{1}, 0}$ for some non-negative integers $l_{1}, m_{1} \leq r$; and
(3) $K^{2}$ is of class $\mathcal{K}_{p, 0, l_{2}}$ and $\mathcal{K}_{p, m_{2}, 0}$ for some non-negative integers $l_{2}, m_{2} \leq r$.

For later reference, we state the following result. The proof follows from Theorem 1 of Sloan and Thomeé [13], on using the fact that $K$ is of class $\mathcal{K}_{p, O, l}$ if $K^{*}$ maps $L_{p}$ into $W_{q}^{l}[\mathbf{1 3}, \mathrm{p} .9]$.

Proposition 3.2. Let $v \in L_{p}$ and $u_{n}^{P}$ be the projection solution of the equation $u-K u=v$. Then

$$
\left\|u-u_{n}^{P}\right\|_{-l, p} \leq c h^{l}\left\|u-u_{n}^{P}\right\|_{p}
$$

where $l \in\left\{l_{1}, l_{2}\right\}$, and $v \in L_{p}$.

We are now in a position to state and prove

Theorem 3.3. Under the above smoothness assumptions on $K$, it follows that
(i) $\left\|x-x_{n}^{K}\right\| \leq c h^{m_{1}}\|x\|$,
(ii) $\left\|x-x_{n}^{\mathrm{S}}\right\| \leq \operatorname{ch}^{l_{1}}\left\|x-x_{n}^{\mathrm{P}}\right\|$,
(iii) $\left\|x-\tilde{x}_{n}^{K}\right\| \leq c h^{l_{1}+m_{1}}\|x\|$,
(iv) $\left\|x-\tilde{x}_{n}^{\mathrm{S}}\right\| \leq c h^{l_{2}}\left\|x-x_{n}^{\mathrm{P}}\right\|$,
where $c$ denotes a generic constant.

Proof. (i). From the relations (2.2) and (3.1), we have

$$
\begin{aligned}
\left\|x-x_{n}^{K}\right\| & \leq\left\|\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}\right\|\left\|\left(I-\boldsymbol{\pi}_{n}\right) K x\right\| \\
& \leq\left\|\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}\right\|\left(1+\left\|\boldsymbol{\pi}_{n}\right\|\right) d_{p}\left(K x, S_{r, n}\right)
\end{aligned}
$$

Since $K$ is of class $\mathcal{K}_{p, m_{1}, 0}$, we have

$$
K x \in W_{p}^{m_{1}}, \quad\|K x\|_{m_{1}, p} \leq c\|x\|
$$

Hence, by Proposition 3.1, we obtain

$$
d_{p}\left(K x, S_{r, n}\right) \leq c h^{m_{1}}\|K x\|_{m_{1}, p} \leq c h^{m_{1}}\|x\|
$$

which yields the required result.
(ii). Since $K$ is of class $\mathcal{K}_{p, 0, l_{1}}$,

$$
\left\|x-x_{n}^{\mathrm{S}}\right\|=\left\|K\left(x-x_{n}^{\mathrm{P}}\right)\right\| \leq c\left\|x-x_{n}^{\mathrm{P}}\right\|_{-l_{1}, p}
$$

Applying Proposition 3.2, we obtain the required result.
(iii) Since $K$ is of class $\mathcal{K}_{p, 0, l_{1}}$,

$$
\left\|x-\tilde{x}_{n}^{K}\right\|=\left\|K\left(x-x_{n}^{K}\right)\right\|=\left\|K\left(z-z_{n}^{\mathrm{P}}\right)\right\| \leq c\left\|z-z_{n}^{\mathrm{P}}\right\|_{-l_{1}, p}
$$

where $z=K x$ and $z_{n}^{\mathrm{P}}$ is the projection approximation of the solution of $z-K z=K y$. Now, by Proposition 3.3 and (i), we obtain the required result,

$$
\left\|x-\tilde{x}_{n}^{K}\right\| \leq c h^{l_{1}}\left\|z-z_{n}^{\mathrm{P}}\right\|=c h^{l_{1}}\left\|x-x_{n}^{K}\right\| \leq c h^{l_{1}+m_{1}}\|x\| .
$$

(iv). Since $K^{2}$ is of class $\mathcal{K}_{p, 0, l_{2}}$, we have, on using Proposition 3.2, that

$$
\left\|x-\tilde{x}_{n}^{\mathrm{S}}\right\|=\left\|K^{2}\left(x-x_{n}^{\mathrm{P}}\right)\right\| \leq c\left\|x-x_{n}^{\mathrm{P}}\right\|_{-l_{2}, p} \leq c h^{l_{2}}\left\|x-x_{n}^{\mathrm{P}}\right\|
$$

which is the required result.

Next, we state and prove the main theorem of this section.

Theorem 3.4. Under the above smoothness assumptions on $K$, it follows that
(i) $\left\|x-\hat{x}_{n}^{M}\right\| \leq c \max \left\{h^{l_{1}+m_{1}}, h^{m_{2}}\right\}$,
(ii) $\left\|x-\tilde{x}_{n}^{M}\right\| \leq c \max \left\{h^{l_{1}+m_{1}}, h^{l_{2}}\left\|x-x_{n}^{\mathrm{P}}\right\|, h^{l_{1}+m_{1}}\left\|x-x_{n}^{\mathrm{P}}\right\|\right\}$.

Proof. Initially, we show that
(a) $\left\|\left(I-\pi_{n}\right) K^{2} x\right\| \leq c h^{m_{2}}\|x\|$ and
(b) $\left\|K\left(I-\pi_{n}\right)\left(x-x_{n}^{\mathrm{S}}\right) \leq c h^{l_{1}+m_{1}}\right\| x-x_{n}^{\mathrm{P}} \|$.

The required results (i) and (ii) then follow from Theorems 2.1 and 3.3.
Using Proposition 3.1 and the fact that $K^{2}$ is of class $\mathcal{K}_{p, m_{2}, 0}$, we obtain

$$
\left\|\left(I-\boldsymbol{\pi}_{n}\right) K^{2} x\right\| \leq c h^{m_{2}}\left\|K^{2} x\right\|_{l_{2}, p} \leq c h^{m_{2}}\|x\|
$$

which proves (a). Since $K$ is of class $\mathcal{K}_{p, 0, l_{1}}$,

$$
\begin{aligned}
\left\|K\left(I-\boldsymbol{\pi}_{n}\right)\left(x-x_{n}^{\mathrm{S}}\right)\right\| & =\left\|K\left(I-\boldsymbol{\pi}_{n}\right) \mathbf{K}\left(x-x_{n}^{\mathrm{P}}\right)\right\| \\
& \leq c\left\|\left(I-\boldsymbol{\pi}_{n}\right) K\left(x-x_{n}^{\mathrm{P}}\right)\right\|_{-l_{1}, p} .
\end{aligned}
$$

It now follows from (3.2) (the definition of the negative norm $\|\cdot\|_{-l_{1}, p}$ ), that it is sufficient to prove that

$$
\left|\left\langle\left(I-\boldsymbol{\pi}_{n}\right) K\left(x-x_{n}^{\mathrm{P}}\right), v\right\rangle\right| \leq c h^{l_{1}+m_{1}}\left\|x-x_{n}^{\mathrm{P}}\right\|\|v\|_{l_{1}, q}
$$

for every $v \in W_{q}^{l_{1}}$. Note that

$$
\begin{aligned}
\left\langle\left(I-\boldsymbol{\pi}_{n}\right) \mathbf{K}\left(x-x_{n}^{\mathrm{P}}\right), v\right\rangle & =\left\langle\left(I-\boldsymbol{\pi}_{n}\right) K\left(x-x_{n}^{\mathrm{P}}\right), v-u_{n}\right\rangle \\
& \leq\left\|\left(I-\boldsymbol{\pi}_{n}\right) K\left(x-x_{n}^{\mathrm{P}}\right)\right\|_{0, p}\left\|v-u_{n}\right\|_{o, q}
\end{aligned}
$$

for every $u_{n} \in S_{r, n}$. Hence,

$$
\begin{aligned}
\left|\left\langle\left(I-\boldsymbol{\pi}_{n}\right) K\left(x-x_{n}^{\mathrm{P}}\right), v\right\rangle\right| & \leq\left\|\left(I-\boldsymbol{\pi}_{n}\right) K\left(x-x_{n}^{\mathrm{P}}\right)\right\| d_{q}\left(v, S_{r, n}\right) \\
& \leq c d_{p}\left(K\left(x-x_{n}^{\mathrm{P}}\right), S_{r, n}\right) d_{q}\left(v, S_{r, n}\right) \\
& \leq c h^{l_{1}+m_{1}}\left\|K\left(x-x_{n}^{\mathrm{P}}\right)\right\|_{m_{1}, p}\|v\|_{l_{1}, p} \\
& \leq c h^{l_{1}+m_{1}}\left\|x-x_{n}^{\mathrm{P}}\right\|\|v\|_{l_{1}, q}
\end{aligned}
$$

Here we have used the relation (3.1), Proposition 3.1, and the fact that $K$ is of class $\mathcal{K}_{p, m_{1}, 0} . \square$

If the non-homogeneous part $y$ in (1.1) belongs to $W_{p}^{\rho}$, for some nonnegative integer $\rho$, then we find that

$$
x=y+K x \in W_{p}^{\beta}, \quad \beta=\min \left\{\rho, m_{1}\right\}
$$

Thus,

$$
\begin{aligned}
\left\|\left(I-\boldsymbol{\pi}_{n}\right) x\right\| & \leq\left(1+\left\|\boldsymbol{\pi}_{n}\right\|\right) d_{p}\left(x, S_{r, n}\right) \\
& \leq c h^{\beta}\|x\|_{\beta, p}
\end{aligned}
$$

so that

$$
\left\|x-x_{n}^{\mathrm{P}}\right\|=\left\|\left(I-\boldsymbol{\pi}_{n} K\right)^{-1}\left(I-\boldsymbol{\pi}_{n}\right) x\right\| \leq c h^{\beta}\|x\|_{\beta, p}
$$

Now, as a corollary to the above two theorems we can state

Theorem 3.5. In addition to the above smoothness conditions on $K$, let $y \in W_{p}^{\rho}$ for some non-negative integer $\rho$. Then, for a generic constant $c$ depending on $x$, we obtain the following inequalities:
(i) $\left\|x-x_{n}^{\mathrm{P}}\right\| \leq c h^{\beta}, \quad \beta=\min \left\{\rho, m_{1}\right\}$,
(ii) $\left\|x-x_{n}^{K}\right\| \leq c h^{m_{1}}$,
(iii) $\left\|x-x_{n}^{\mathrm{S}}\right\| \leq c h^{l_{1}+\beta}$,
(iv) $\left\|x-\tilde{x}_{n}^{K}\right\| \leq c h^{l_{1}+m_{1}}$,
(v) $\left\|x-\tilde{x}_{n}^{\mathrm{S}}\right\| \leq c h^{l_{2}+\beta}$,
(vi) $\left\|x-\hat{x}_{n}^{M}\right\| \leq c \max \left\{h^{l_{1}+m_{1}}, h^{m_{2}}\right\}$,
(vii) $\left\|x-\tilde{x}_{n}^{M}\right\| \leq c \max \left\{h^{l_{1}+m_{1}}, h^{l_{2}+\beta}\right\}$.

Remark 3.6. From Theorem 3.5, we find that if $m_{2} \geq l_{1}+m_{1}$ (respectively $l_{2}+\beta \geq l_{1}+m_{1}$ ), then $\hat{x}_{n}^{M}$ (respectively $\tilde{x}_{n}^{M}$ ) has a better order of accuracy than $x_{n}^{\mathrm{P}}, x_{n}^{K}$, and $x_{n}^{\mathrm{S}}$ and has the same order of accuracy as $\tilde{x}_{n}^{K}$.

Thus, if $y$ and $K$ are not so smooth, but $K^{2}$ is, then $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$ have better order of convergence than $x_{n}^{K}$ and $x_{n}^{\mathrm{S}}$, and they do not require as much computation as for $\tilde{x}_{n}^{K}$ and $\tilde{x}_{n}^{\mathrm{S}}$. We also remark that if $X=L_{\infty}, y \in C^{\rho}, k \in C^{l_{1}}$ with $\rho<l_{1}=m_{1}$, then the results in Theorem 3.5, (i)-(iv), are the same as those in Sloan [11, (4.4)-(4.7)].

Example 3.7. Consider the following example of Sloan and Thomeé [13], where $X=L_{2}[0,2 \pi]$ and

$$
(K v)(s)=\lambda \int_{0}^{2 \pi} \log \left|\left(\frac{s-t}{2}\right)\right| v(s) d s
$$

assuming $\lambda$ is not an eigenvalue of $K$. They show that
(i) $K=K^{*}$ is of class $\mathcal{K}_{p, 0,1}$ and $\mathcal{K}_{p, 1,0}$, and
(ii) $K^{i}$ is of class $\mathcal{K}_{p, m, l}$ for $i \geq m+l$.

Thus we have $l_{1}=1=m_{1}$ and $l_{2}=2=m_{2}=l_{1}+m_{1}$.

Computational Remarks 3.8. The following expressions show explicitly the quantities to be computed:

$$
\begin{aligned}
\tilde{x}_{n}^{K} & =y+K x_{n}^{K}=y+K y+K z_{n}^{\mathrm{P}} \\
\tilde{x}_{n}^{\mathrm{S}} & =y+K x_{n}^{\mathrm{S}}=y+K y+K^{2} x_{n}^{\mathrm{P}} \\
\hat{x}_{n}^{M} & =y+z_{n}^{M}=y+K y+\left(z_{n}^{\mathrm{P}}-\boldsymbol{\pi}_{n} K y\right) \\
\tilde{x}_{n}^{M} & =y+K x_{n}^{M}=y+K y+K\left(x_{n}^{\mathrm{P}}-\boldsymbol{\pi}_{n} y\right)
\end{aligned}
$$

We note that $\boldsymbol{\pi}_{n} K y$ is required to compute $z_{n}^{\mathrm{P}}$, and $\boldsymbol{\pi}_{n} y$ is needed to compute $x_{n}^{\mathrm{P}}$. Thus $\hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$ require less computations than $\tilde{x}_{n}^{K}$ and $\tilde{x}_{n}^{S}$.

The following table shows the operations to be performed in obtaining $\tilde{x}_{n}^{K}, \tilde{x}_{n}^{\mathrm{S}}, \hat{x}_{n}^{M}$ and $\tilde{x}_{n}^{M}$ except for the common expression $y+K y$ :

| Method | Solve | Compute |
| :---: | :---: | :---: |
| $\tilde{x}_{n}^{K}$ | $u_{n}-\boldsymbol{\pi}_{n} K u_{n}=\boldsymbol{\pi}_{n} K y$ | $K u_{n}$ |
| $\tilde{x}_{n}^{\mathrm{S}}$ | $v_{n}-\boldsymbol{\pi}_{n} K v_{n}=\boldsymbol{\pi}_{n} y$ | $K^{2} v_{n}$ |
| $\hat{x}_{n}^{M}$ | $u_{n}-\boldsymbol{\pi}_{n} K u_{n}=\boldsymbol{\pi}_{n} K y$ | $u_{n}-\boldsymbol{\pi}_{n} K y$ |
| $\tilde{x}_{n}^{M}$ | $v_{n}-\boldsymbol{\pi}_{n} K v_{n}=\boldsymbol{\pi}_{n} y$ | $K\left(v_{n}-\boldsymbol{\pi}_{n} y\right)$ |

Acknowledgments. One of the authors (M. Thamban Nair) did this work while visiting the Centre for Mathematical Analysis, Australian National University, in March-June, 1989.

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