# Iterative method for solving the Neumann boundary value problem for biharmonic type equation ${ }^{2 / 2}$ 

Quang A Dang*<br>Institute of Information Technology, 8 Hoang Quoc Viet, Cau giay, Hanoi, Vietnam

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#### Abstract

The solution of boundary value problems (BVP) for fourth order differential equations by their reduction to BVP for second order equations, with the aim to use the achievements for the latter ones attracts attention from many researchers. In this paper, using the technique developed by ourselves in recent works, we construct iterative method for the Neumann BVP for biharmonic type equation. The convergence rate of the method is proved and some numerical experiments are performed for testing it in dependence on the choice of an iterative parameter.


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## 1. Introduction

The solution of fourth order differential equations by their reduction to boundary value problems (BVP) for the second order equations, with the aim of using efficient algorithms for these, attracts attention from many researchers. Namely, for the biharmonic equation with the Dirichlet boundary condition, there is intensively developed the iterative method, which leads the problem to two problems for the Poisson equation at each iteration (see e.g. [4, 8,9,11]). Recently, Abramov and Ulijanova [1] proposed an iterative method for the Dirichlet problem for the biharmonic type equation, but the convergence of the method is not proved. In our previous works [6,7], with the help of boundary or mixed boundary-domain operators appropriately introduced, we constructed iterative methods for biharmonic and biharmonic type equations associated with the Dirichlet boundary condition. It is proved that the methods are convergent with the rate of geometric progression. In this paper, we develop our technique in [4-7] for the Neumann BVP for the biharmonic type equation. Namely, we consider the following problem:

$$
\begin{align*}
& L u \equiv \Delta^{2} u-a \Delta u+b u=f \quad \text { in } \Omega,  \tag{1}\\
& \frac{\partial u}{\partial v}=g_{0}, \quad \frac{\partial \Delta u}{\partial v}=g_{1} \text { on } \Gamma, \tag{2}
\end{align*}
$$

[^0]where $\Delta$ is the Laplace operator, $\Omega$ is a bounded domain in $R^{n}(n \geqslant 2), \Gamma$ is the sufficiently smooth boundary of $\Omega, v$ is the outward normal to $\Gamma$ and $a, b$ are positive constants. Eq. (1) with other boundary conditions are met, for example, in $[2,3]$. An iterative method reducing the problem to a sequence of Neumann problems for Poisson equation will be proposed and investigated. Two different cases will be treated in dependence on the sign of $a^{2}-4 b$.

## 2. Case $a^{2}-4 b \geqslant 0$

In this case, we always can lead the original problem (1)-(2) to two Neumann problems for second order equation. To do this, let $\lambda$ and $\mu$ be the roots of the quadratic equation

$$
\xi^{2}-a \xi+b=0
$$

namely,

$$
\begin{equation*}
\lambda=\frac{a-\sqrt{a^{2}-4 b}}{2}, \quad \mu=\frac{a+\sqrt{a^{2}-4 b}}{2} . \tag{3}
\end{equation*}
$$

Clearly, $\lambda$ and $\mu>0$. Put

$$
\begin{equation*}
L_{1}=\Delta-\mu, \quad L_{2}=\Delta-\lambda . \tag{4}
\end{equation*}
$$

Then we have $L=L_{1} L_{2}$ and the problem (1)-(2) is decomposed to the following problems:

$$
\begin{align*}
& L_{1} v \equiv \Delta v-\mu v=f(x) \quad \text { in } \Omega, \\
& \frac{\partial v}{\partial v}=g_{1}-\lambda g_{0} \quad \text { on } \Gamma,  \tag{5}\\
& L_{2} u \equiv \Delta u-\lambda u=v(x) \quad \text { in } \Omega, \\
& \frac{\partial u}{\partial v}=g_{0} \quad \text { on } \Gamma . \tag{6}
\end{align*}
$$

These Neumann problems can be solved by known methods such as finite element method, boundary element method or finite difference method.

## 3. Case $a^{2}-4 b<0$

### 3.1. Construction of iterative method

Suppose $\lambda$ and $\mu$ are positive numbers such that

$$
\begin{equation*}
\mu \geqslant \lambda, \quad \lambda+\mu=a, \quad b_{1}=\lambda \mu<b . \tag{7}
\end{equation*}
$$

Using the notations $L_{1}$ and $L_{2}$ given by (4) we set

$$
v=L_{2} u=\Delta u-\lambda u .
$$

Then we get

$$
L_{1} v=\Delta v-\mu v=\Delta^{2} u-a \Delta u+b_{1} u .
$$

Now, putting

$$
\begin{equation*}
\varphi=-\alpha u \tag{8}
\end{equation*}
$$

where

$$
\begin{equation*}
\alpha=b-b_{1} \tag{9}
\end{equation*}
$$

we can reduce the original problem (1), (2) to the problems

$$
\begin{align*}
& L_{1} v=f+\varphi \quad \text { in } \Omega, \quad \frac{\partial v}{\partial v}=g_{1}-\lambda g_{0} \quad \text { on } \Gamma,  \tag{10}\\
& L_{2} u=v(x) \quad \text { in } \Omega, \quad \frac{\partial u}{\partial v}=g_{0} \quad \text { on } \Gamma, \tag{11}
\end{align*}
$$

where $\varphi$ as $u$ is an unknown function but it is related with $u$ by (8). Now consider the following iterative process for finding $\varphi$ and simultaneously for finding $u$.

1. Given $\varphi^{(0)} \in L_{2}(\Omega)$, for example, $\varphi^{(0)}=0$ in $\Omega$.
2. Knowing $\varphi^{(k)}(x)$ in $\Omega(k=0,1, \ldots)$ solve consecutively two problems

$$
\begin{align*}
& L_{1} v^{(k)}=f+\varphi^{(k)} \quad \text { in } \Omega, \quad \frac{\partial v^{(k)}}{\partial v}=g_{1}-\lambda g_{0} \quad \text { on } \Gamma,  \tag{12}\\
& L_{2} u^{(k)}=v^{(k)} \quad \text { in } \Omega, \quad \frac{\partial u^{(k)}}{\partial v}=g_{0} \quad \text { on } \Gamma . \tag{13}
\end{align*}
$$

3. Compute the new approximation

$$
\begin{equation*}
\varphi^{(k+1)}=\left(1-\tau_{k+1}\right) \varphi^{(k)}-\alpha \tau_{k+1} u^{(k)}, \tag{14}
\end{equation*}
$$

where $\tau_{k+1}$ is an iterative parameter to be chosen later.

### 3.2. Investigation of convergence

In order to investigate the convergence of the iterative process (12)-(14), firstly we rewrite (14) in the canonical form of two-layer iterative scheme [12]

$$
\begin{equation*}
\frac{\varphi^{(k+1)}-\varphi^{(k)}}{\tau_{k+1}}+\left(\varphi^{(k)}+\alpha u^{(k)}\right)=0 \tag{15}
\end{equation*}
$$

Now, we introduce the operator $A$ defined by the formula

$$
\begin{equation*}
A \varphi=u \tag{16}
\end{equation*}
$$

where $u$ is the function determined from the problems

$$
\begin{array}{lll}
L_{1} v=\varphi \quad \text { in } \Omega, & \frac{\partial v}{\partial v}=0 \quad \text { on } \Gamma, \\
L_{2} u=v \quad \text { in } \Omega, & \frac{\partial u}{\partial v}=0 \quad \text { on } \Gamma . \tag{18}
\end{array}
$$

The properties of the operator $A$ will be investigated in the sequel. Now, let us return to the problems (10), (11). We represent their solutions in the form

$$
\begin{equation*}
u=u_{1}+u_{2}, \quad v=v_{1}+v_{2} \tag{19}
\end{equation*}
$$

where $u_{1}, u_{2}, v_{1}, v_{2}$ are the solutions of the problems

$$
\begin{align*}
& L_{1} v_{1}=\varphi \quad \text { in } \Omega, \quad \frac{\partial v_{1}}{\partial v}=0 \quad \text { on } \Gamma  \tag{20}\\
& L_{2} u_{1}=v_{1} \quad \text { in } \Omega, \quad \frac{\partial u_{1}}{\partial v}=0 \quad \text { on } \Gamma .  \tag{21}\\
& L_{1} v_{2}=f \quad \text { in } \Omega, \quad \frac{\partial v_{2}}{\partial v}=g_{1}-\lambda g_{0} \quad \text { on } \Gamma,  \tag{22}\\
& L_{2} u_{2}=v_{2} \quad \text { in } \Omega, \quad \frac{\partial u_{2}}{\partial v}=g_{0} \quad \text { on } \Gamma . \tag{23}
\end{align*}
$$

According to the definition of $A$ we have

$$
\begin{equation*}
A \varphi=u_{1} \tag{24}
\end{equation*}
$$

Since $\varphi$ should satisfy the relation (8), taking into account the representation (19) we obtain the equation

$$
\begin{equation*}
(I+\alpha A) \varphi=-\alpha u_{2} \tag{25}
\end{equation*}
$$

where $I$ is the identity operator.
Thus, we have reduced the original problem (1), (2) to the operator (25), whose right-hand side is completely defined by the data $f, g_{0}$ and $g_{1}$.

Proposition 3.1. The iterative process (12)-(14) is a realization of the two-layer iterative scheme

$$
\begin{equation*}
\frac{\varphi^{(k+1)}-\varphi^{(k)}}{\tau_{k+1}}+(I+\alpha A) \varphi^{(k)}=-\alpha u_{2}, \quad k=0,1,2, \ldots \tag{26}
\end{equation*}
$$

for the operator equation (25).
Proof. Indeed, if in (12), (13) we put

$$
\begin{equation*}
u^{(k)}=u_{1}^{(k)}+u_{2}, \quad v^{(k)}=v_{1}^{(k)}+v_{2} \tag{27}
\end{equation*}
$$

where $u_{2}, v_{2}$ are the solutions of Problems (22), (23), then we get

$$
\begin{align*}
& L_{1} v_{1}^{(k)}=\varphi^{(k)} \quad \text { in } \Omega, \quad \frac{\partial v_{1}^{(k)}}{\partial v}=0 \quad \text { on } \Gamma,  \tag{28}\\
& L_{2} u_{1}^{(k)}=v_{1}^{(k)} \quad \text { in } \Omega, \quad \frac{\partial u_{1}^{(k)}}{\partial v}=0 \quad \text { on } \Gamma . \tag{29}
\end{align*}
$$

From here it is easy to see that

$$
A \varphi^{(k)}=u_{1}^{(k)}
$$

Therefore, taking into account the first relation of (27) and the above equality, from (15) we obtain (26). Thus, the proposition is proved.

Proposition 1 enables us to lead the investigation of convergence of process (12)-(14) to the study of the scheme (26). For this reason we need some properties of the operator $A$.

Proposition 3.2. If the numbers $\lambda$ and $\mu$ satisfy the condition (7) then the operator $A$ defined by (16)-(18) is linear, symmetric, positive and compact operator in the space $L_{2}(\Omega)$.

Proof. The linearity of $A$ is obvious. To establish the other properties of $A$ let us consider the inner product $(A \varphi, \bar{\varphi})$ for two arbitrary functions $\varphi$ and $\bar{\varphi}$ in $L_{2}(\Omega)$. Recall that the operator $A$ is defined by (16)-(18). We denote by $\bar{u}$ and $\bar{v}$ the solutions of (17) and (18), where instead of $\varphi$ there stands $\bar{\varphi}$. We have

$$
\begin{aligned}
(A \varphi, \bar{\varphi}) & =\int_{\Omega} u \cdot \bar{\varphi} \mathrm{~d} x=\int_{\Omega} u(\Delta \bar{v}-\mu \bar{v}) \mathrm{d} x \\
& =\int_{\Omega} u \Delta \bar{v} \mathrm{~d} x-\mu \int_{\Omega} u \bar{v} \mathrm{~d} x=-\int_{\Omega} \nabla u \nabla \bar{v} \mathrm{~d} x-\mu \int_{\Omega} u \bar{v} \mathrm{~d} x .
\end{aligned}
$$

Noting that

$$
\int_{\Omega} \bar{v} \Delta u \mathrm{~d} x=-\int_{\Omega} \nabla \bar{v} \nabla u \mathrm{~d} x
$$

from the latter equality we get

$$
(A \varphi, \bar{\varphi})=\int_{\Omega} \bar{v} \Delta u \mathrm{~d} x-\mu \int_{\Omega} u \bar{v} \mathrm{~d} x .
$$

In the above relation replacing $\Delta u=v+\lambda u$ we obtain

$$
(A \varphi, \bar{\varphi})=\int_{\Omega} v \bar{v} \mathrm{~d} x+(\lambda-\mu) \int_{\Omega} u \bar{v} \mathrm{~d} x .
$$

Further, since $\bar{v}=\Delta \bar{u}-\lambda \bar{u}$ we have

$$
\int_{\Omega} u \bar{v} \mathrm{~d} x=\int_{\Omega} u(\Delta \bar{u}-\lambda \bar{u}) \mathrm{d} x=-\int_{\Omega} u \nabla u \nabla \bar{u} \mathrm{~d} x-\lambda \int_{\Omega} u \bar{u} \mathrm{~d} x .
$$

Hence,

$$
(A \varphi, \bar{\varphi})=\int_{\Omega} v \bar{v} \mathrm{~d} x+(\mu-\lambda)\left(\int_{\Omega} \nabla u \nabla \bar{u} \mathrm{~d} x+\lambda \int_{\Omega} u \bar{u} \mathrm{~d} x\right) .
$$

Clearly,

$$
(A \varphi, \bar{\varphi})=(A \bar{\varphi}, \varphi)
$$

and

$$
(A \varphi, \varphi)=\int_{\Omega} v^{2} \mathrm{~d} x+(\mu-\lambda)\left(\int_{\Omega}|\nabla u|^{2} \mathrm{~d} x+\lambda \int_{\Omega} u^{2} \mathrm{~d} x\right) \geqslant 0
$$

due to $\mu \geqslant \lambda$. Moreover, it is easy seen that $(A \varphi, \varphi)=0$ if and only if $\varphi=0$. So, we have shown that the operator $A$ is symmetric and positive in $L_{2}(\Omega)$.

It remains to show the compactness of $A$. As is well known that if $\varphi \in L_{2}(\Omega)$ then the problem (17) has a unique solution $v \in H^{2}(\Omega)$, therefore, the problem (18) has a unique solution $v \in H^{4}(\Omega)$. Consequently, the operator $A$ maps $L_{2}(\Omega)$ into $H^{4}(\Omega)$. In view of the compact imbedding $H^{4}(\Omega)$ into $L_{2}(\Omega)$ it follows that $A$ is compact operator in $L_{2}(\Omega)$.

Thus, the proof of Proposition 3.2 is complete.
Due to the above proposition the operator

$$
\begin{equation*}
B=I+\alpha A \tag{30}
\end{equation*}
$$

is linear, symmetric, positive definite and bounded operator in the space $L_{2}(\Omega)$. More precisely, we have

$$
\begin{equation*}
\gamma_{1} I<B \leqslant \gamma_{2} I, \tag{31}
\end{equation*}
$$

where

$$
\begin{equation*}
\gamma_{1}=1, \quad \gamma_{2}=1+\alpha\|A\| . \tag{32}
\end{equation*}
$$

In the future, we need the following estimate for the function $u$ given by (17) and (18)

$$
\begin{equation*}
\|u\|_{H^{4}(\Omega)} \leqslant C\|\varphi\|_{L_{2}(\Omega)} \tag{33}
\end{equation*}
$$

which follows from the general theory of BVP [10].
Before stating the result of convergence of the iterative process (12)-(14) we assume the following regularity of the data of the original problem (1), (2):

$$
f \in L_{2}(\Omega), \quad g_{0} \in H^{5 / 2}(\Gamma) \text { and } g_{1} \in H^{1 / 2}(\Gamma) .
$$

Then the problem (1), (2) has a unique solution $u \in H^{4}(\Omega)$. For the function $u_{2}$ determined by (22), (23) we have also $u_{2} \in H^{4}(\Omega)$.

Theorem 3.3. Let $u$ be the solution of Problems (1), (2) and $\varphi$ be the solution of Eq. (25). Then, if $\left\{\tau_{k+1}\right\}$ is the Chebyshev collection of parameters, constructed by the bounds $\gamma_{1}$ and $\gamma_{2}$ of the operator $B$, namely

$$
\begin{equation*}
\tau_{0}=\frac{2}{\gamma_{1}+\gamma_{2}}, \quad \tau_{k}=\frac{\tau_{0}}{\rho_{0} t_{k}+1}, \quad t_{k}=\cos \frac{2 k-1}{2 M} \pi, \quad k=1, \ldots, M \quad \rho_{0}=\frac{1-\xi}{1+\xi}, \quad \xi=\frac{\gamma_{1}}{\gamma_{2}} \tag{34}
\end{equation*}
$$

we have

$$
\begin{equation*}
\left\|u^{(M)}-u\right\|_{H^{4}(\Omega)} \leqslant C_{1} q_{M}, \tag{35}
\end{equation*}
$$

where

$$
\begin{equation*}
C_{1}=C\left\|\varphi^{(0)}-\varphi\right\|_{L_{2}(\Omega)} \tag{36}
\end{equation*}
$$

$C$ being the constant in (33),

$$
\begin{equation*}
q_{M}=\frac{2 \rho_{1}^{M}}{1+\rho_{1}^{2 M}}, \quad \rho_{1}=\frac{1-\sqrt{\xi}}{1+\sqrt{\xi}} . \tag{37}
\end{equation*}
$$

In the case of the stationary iterative process, $\tau_{k}=\tau_{0}(k=1,2, \ldots)$ we have

$$
\begin{equation*}
\left\|u^{(k)}-u\right\|_{H^{4}(\Omega)} \leqslant C_{1} \rho_{0}^{k}, \quad k=1,2, \ldots \tag{38}
\end{equation*}
$$

Proof. According to the general theory of the two-layer iterative schemes (see [12]) for the operator equation (25) we have

$$
\left\|\varphi^{(M)}-\varphi\right\|_{L_{2}(\Omega)} \leqslant q_{M}\left\|\varphi^{(0)}-\varphi\right\|_{L_{2}(\Omega)},
$$

if the parameter $\left\{\tau_{k+1}\right\}$ is chosen by the formulae (34) and

$$
\left\|\varphi^{(k)}-\varphi^{*}\right\|_{H^{4}(\Omega)} \leqslant \rho_{0}^{k}\left\|\varphi^{(0)}-\varphi^{*}\right\|_{L_{2}(\Omega)}, \quad k=1,2, \ldots
$$

if $\tau_{k}=\tau_{0}(k=1,2, \ldots)$. In view of these estimates the corresponding estimates (35) and (38) follow from (33) applied to the problems

$$
\begin{array}{llll}
L_{1}\left(v^{(k)}-v\right)=\varphi^{(k)}-\varphi & \text { in } \Omega, & \frac{\partial}{\partial v}\left(v^{(k)}-v\right)=0 & \text { on } \Gamma, \\
L_{2}\left(u^{(k)}-u\right)=v^{(k)}-v & \text { in } \Omega, & \frac{\partial}{\partial v}\left(u^{(k)}-u\right)=0 & \text { on } \Gamma,
\end{array}
$$

which are obtained in the result of the subtraction (10) and (11) from (12) and (13), respectively. The theorem is proved.

Remark. Theorem 3.3 theoretically ensures the error estimate (35) and the convergence of the iterative method (26) when the computing process is ideal, that is, all the computations are carried out with an infinity number of significant digits. But any computer makes calculations with a finite speed and a finite number of digits. In this case rounding errors may accumulate and the computing process may be unstable. In order to overcome this computational instability following Samarskii and Nikolaev [13,12] it is needed to use a stable collection of parameters $\left\{\tau_{k}^{*}\right\}$ calculated by the formula

$$
\tau_{k}^{*}=\frac{\tau_{0}}{\rho_{0} t_{k}^{*}+1}, \quad t_{k}^{*}=-\cos \left(\frac{\pi}{2 M} \theta_{M}^{*}(k)\right), \quad k=1, \ldots, M
$$

where $\theta_{M}^{*}$ is the sequence of the odd integers from 1 to $2 M-1$, determined by a rule given there.
An alternative way to treat the numerical instability of the Richardon's iterative method (26) is the use of the semiiterative method based on stationary Richardon method as was recommended in [14] for solving a system of linear algebraic equations. But it requires extra computational work in computing sums of iterations.

In the case of the stationary process the value $\tau_{0}$, as well known in [12,13], is optimal. But for calculating it we need to know $\|A\|$, which is difficult to be determined. Therefore, in the next section, first we shall determine $\|A\|$ experimentally from one particular case, and then use it to find nearly optimal value of the iterative parameter in some other cases.

## 4. Numerical results

We performed some limited experiments in MATLAB for testing the convergence of the iterative process (12)-(14) in dependence on iterative parameters, which are taken fixed at all iterations, i.e. $\tau_{k} \equiv \tau$. In the examples considered below the computational domain is a rectangle with the uniform grid including $65 \times 65$ nodes. The Neumann problems for the second order equations (12), (13) are discretized by difference schemes of second order approximation obtained by a variational method. The stopping criterion for the iterative process is $\left\|\varphi^{(k)}\right\|_{\infty}<10^{-4}$.

Example 1. We take an exact solution $u=\sin x \sin y$ in the rectangle $[0, \pi] \times[0, \pi]$ and vary the coefficients $a$ and $b$ in the equation (1). Hence, the right-hand side of the equation is $f=(4+2 a+b) \sin x \sin y$.

In all two examples we take $a=1, \lambda=\mu=0.5$ and change only $b$. First, we take $b=1$ and make experiments for testing the convergence of the iterative method in dependence on the choice of the iterative parameter $\tau$. The results of computation are presented in Table 1, where $k$ is the number of iterations, err is the error of approximate solution $u_{a p}$, err $=\left\|u_{a p}-u\right\|_{\infty}$.

From Table 1 we see that it is possible to adopt experimentally $\tau_{\mathrm{opt}}=0.36$. Meanwhile, the formulae (32), (34) give

$$
\begin{equation*}
\tau_{\mathrm{opt}}=\frac{2}{2+\alpha\|A\|}, \tag{39}
\end{equation*}
$$

where $\alpha=b-\lambda \mu$. From here we can calculate $\|A\| \approx 4.7407$. Note that for a fixed domain $\|A\|$ depends only on $\lambda$ and $\mu$. Therefore, this value of $\|A\|$ can be used for calculating good iterative parameter for the cases of other values

Table 1
Case $b=1$

| $\tau$ | $k$ | err |
| :--- | :--- | :--- |
| 0.20 | 30 | 0.0006 |
| 0.30 | 20 | 0.0006 |
| 0.35 | 16 | 0.0006 |
| 0.36 | 15 | 0.0006 |
| 0.37 | 16 | 0.0006 |
| 0.40 | 18 | 0.0005 |
| 0.45 | 39 | 0.0008 |
| 0.50 | $\infty$ |  |

Table 2
Case $b=0.5$

| $\tau$ | $k$ | err |
| :--- | :---: | :---: |
| 0.40 | 13 | 0.0011 |
| 0.50 | 13 | 0.0011 |
| 0.55 | 9 | 0.0011 |
| 0.60 | 8 | 0.0011 |
| $\mathbf{0 . 6 2 7 9}$ | 7 | 0.0011 |
| 0.65 | 7 | 0.0011 |
| 0.70 | 7 | 0.0009 |
| 0.80 | 8 | 0.0013 |

Table 3
Case $b=1.5$

| $\tau$ | $k$ | err |
| :--- | :--- | :--- |
| 0.20 | 31 | 0.0005 |
| 0.22 | 29 | 0.0005 |
| 0.25 | 25 | 0.0005 |
| 0.27 | 24 | 0.0005 |
| $\mathbf{0 . 2 5 2 3}$ | 25 | 0.0005 |
| 0.30 | 41 | 0.0006 |
| 0.32 | 111 | 0.0001 |
| 0.35 | $\infty$ |  |

Table 4
Case $b=1$

| $\tau$ | $k$ | err |
| :--- | :--- | :--- |
| 0.20 | 49 | $1.2 e-5$ |
| 0.30 | 32 | $6.6 e-6$ |
| 0.35 | 25 | $4.9 e-6$ |
| $\mathbf{0 . 3 6}$ | 26 | $4.9 e-4$ |
| 0.37 | 25 | $5.2 e-6$ |
| 0.40 | 24 | $6.8 e-5$ |
| 0.45 | 45 | $2.1 e-4$ |
| 0.50 | $\infty$ | 4 |

Table 5
Case $b=0.5$

| $\tau$ | $k$ | $e r r$ |
| :--- | :--- | :--- |
| 0.4 | 21 | $1.4 e-5$ |
| 0.5 | 16 | $8.1 e-6$ |
| 0.5 | 16 | $8.1 e-6$ |
| 0.55 | 14 | $6.9 e-6$ |
| $\mathbf{0 . 6 2 7 9}$ | 12 | $2.8 e-6$ |
| 0.65 | 11 | $6.9 e-6$ |
| 0.7 | 10 | $2.1 e-4$ |
| 0.8 | 18 | $2.1 e-4$ |

Table 6
Case $b=1.5$

| $\tau$ | $k$ | err |
| :--- | :---: | :---: |
| 0.20 | 50 | $8.1 e-6$ |
| 0.22 | 46 | $5.8 e-6$ |
| 0.25 | 40 | $5.2 e-6$ |
| $\mathbf{0 . 2 5 2 3}$ | 40 | $4.5 e-6$ |
| 0.27 | 38 | $4.5 e-6$ |
| 0.30 | 48 | $2.1 e-4$ |
| 0.32 | 127 | $2.0 e-4$ |
| 0.35 | $\infty$ |  |

of $b$ provided that $a, \lambda, \mu$ are unchanged. The results presented in Tables 2 and 3 support this assertion. In these tables as in Tables 4-6 below the values of $\tau$ in bold face are calculated by the formula (39) with $\|A\|=4.7407$.

Example 2. We take an exact solution $u=\left(x^{2}-4\right)\left(y^{2}-1\right)$ in the rectangle $[-2,-2] \times[-1,1]$. The right-hand side of the equation (1) is $f=8-2 a\left(x^{2}+y^{2}-5\right)+b\left(x^{2}-4\right)\left(y^{2}-1\right)$.

The results of the convergence of the method are presented in Tables 4-6.
From Tables $1-6$ it is clear that the values of iterative parameter $\tau$ computed by (39), where $\|A\|$ is found experimentally, are nearly optimal.

## 5. Concluding remark

In the paper, an iterative method was proposed for reducing the Neumann problem for biharmonic type equation to a sequence of Neumann problems for second order equations. The convergence rate of the method depends on the choice of the iterative parameter $\tau$, whose optimal value is determined by the norm of the operator $A$. In its turn, this norm is fully defined for each domain and each value of the coefficient $a$ in Eq. (1) if a choice of $\lambda$ and $\mu$ is fixed, for example, $\lambda=\mu=a / 2$. By experimental way it is possible to find this quantity. Theoretical estimate for $\|A\|$ is an interesting problem to be studied.

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## References

[1] A.A. Abramov, V.I. Ulijanova, On a method for solving biharmonic type equation with singularly small parameter, J Comput. Math. Math. Phys. 32 (4) (1992) 567-575 (Russian).
[2] J.P. Aubin, Approximation of Elliptic Boundary-value Problems, Wiley-Interscience, New York, 1971.
[3] D. Begis, A. Perronet, The club modulef, in: G. Marchuc, J.L. Lions (Eds.), Numerical Methods in Applied Mathematics, Nauka, Novosibirsk, 1982, pp. 212-236, (Russian).
[4] Q.A. Dang, On an iterative method for solving a boundary value problem for fourth order differential equation, Math. Phys. Nonlinear Mech. 10 (44) (1988) 54-59 (Russian).
[5] Q.A. Dang, Approximate method for solving an elliptic problem with discontinuous coefficients, J. Comput. Appl. Math. 51 (2) (1994) 193-203.
[6] Q.A. Dang, Boundary operator method for approximate solution of biharmonic type equation, J. Math. 22 (1\& 2) (1994) 114-120.
[7] Q.A. Dang, Mixed boundary-domain operator in approximate solution of biharmonic type equation, Vietnam J. Math. 26 (3) (1998) $243-252$.
[8] A. Dorodnisyn, N. Meller, On some approaches to the solution of the stationary Navier-Stoke equation, J. Comput. Math. Math. Phys. 8 (2) (1968) 393-402 (Russian).
[9] R. Glowinski, J-L. Lions, R. Tremoliere, Analyse numerique des inequations variationelles, Dunod, Paris, 1976.
[10] J.-L. Lions, E. Magenes, Problemes aux limites non homogenes et applications, vol. 1, Dunod, Paris, 1968.
[11] B.V. Palsev, On the expansion of the Dirichlet problem and a mixed problem for biharmonic equation into a seriaes of decomposed problems, J. Comput. Math. Math. Phys. 6 (1) (1966) 43-51 (Russian).
[12] A.A. Samarskii, The Theory of Difference Schemes, Marcel Dekker, New York, 2001.
[13] A. Samarskii, E. Nikolaev, Numerical Methods for Grid Equations, vol. 2, Birkhäuser, Basel, 1989.
[14] D.M. Young, Iterative Solution of Large Linear Systems, Academic Press, New York, 1971.


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    * Tel.: 844836 1770; fax: 8447564217.

    E-mail addresses: dangqa@ioit.ncst.ac.vn, dangqa2004@yahoo.com (Q.A. Dang).

