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## MULTIGRID APPLIED TO SINGULAR PERTURBATION PROBLEMS

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## 1 Introduction

This report discusses the application of a multigrid algorithm to the solution of the following one dimensional singular perturbation problem:

$$
\begin{equation*}
-\varepsilon u^{\prime \prime}+b(x) u^{\prime}=f, \quad 0<x<1 \tag{1}
\end{equation*}
$$

with

$$
1 \gg \varepsilon>0, \quad u(0)=u_{0}, \quad u(1)=u_{1} .
$$

Many other authors have discussed the application of various methods of solution to the algebraic problem; in particular see Dorr [Dorr70a], Babuska [Babu69a], Ervin and Layton [Ervi85a], and Kellogg and Tsan [Kell78a].

Much of the literature regarding multigrid methods restricts itself to the solution of nice problems. Indeed most authors require that the linear system be well conditioned in addition to symmetric and positive definite. However, in the case of singular perturbation problems as the coefficient of the second order term tends to zero the usual symmetric discretization fails to be of positive type. Thus the first measure taken in the numerical discretization of these problems is to replace the usual symmetric difference of the first order term with some form of skewed differencing. In particular for problems with turning points this may lead to a system of equations which is ill-conditioned for small $\varepsilon$.

From the standpoint of calculating a numerical approximation to the solution of problem (1) the first question is: does the discretization converge to the continuous solution? Then, assuming it does, how does the multigrid algorithm perform as a solver for this system of linear equations? What modifications, if any, are necessary in the multigrid algorithm?

The main result shows that if the original system of equations resulting from the discretization of problem (1) is of positive type then the theoretical results for the multigrid algorithm developed in Kamowitz and Parter [Kamo85a] and in McCormick and Ruge [McCo82a] apply.

From a computational perspective it is convenient to use for the coarse grid operators of the multigrid algorithm the operators that are the finite difference analogue of the original operator. In Section 6 the rate of convergence of the algorithm using the finite difference version of the coarse grid operators is considered. It is shown that the new rate of convergence is an $O\left(h^{2}\right)$ perturbation of the rate obtained using the Galerkin choice.

## 2 The Discrete Problem

Three model problems are chosen for study, one where the sign of $b(x)$ does not change and two where $b(x)$ changes sign. The three problems are designated

Problem BL

$$
\begin{equation*}
L u \equiv-\varepsilon u^{\prime \prime}+u^{\prime}=0, \quad 0<x<1 \tag{2}
\end{equation*}
$$

Problem TP-1

$$
\begin{equation*}
L u \equiv-\varepsilon u^{\prime \prime}+\left(x-\frac{1}{2}\right) u^{\prime}=0, \quad 0<x<1 \tag{3}
\end{equation*}
$$

Problem TP-2

$$
\begin{equation*}
L u \equiv-\varepsilon u^{\prime \prime}-\left(x-\frac{1}{2}\right) u^{\prime}=0 \quad 0<x<1 . \tag{4}
\end{equation*}
$$

For all three problems the boundary conditions are

$$
u(0)=1, \quad u(1)=3
$$

For the discrete problem as usual let $N>0$ be chosen and set $h=1 /(N+1)$. The interval

$$
\Omega=(0,1)
$$

is discretized to form

$$
\Omega_{h} \equiv\{i h: 1 \leq i \leq N\}
$$

The notation $x_{i}$ refers to the point $i h$ and $u_{i}$ refers to $u\left(x_{i}\right)$. In addition the usual notation for finite differences is used:

$$
\begin{gathered}
D_{+} u_{i}=\frac{u_{i+1}-u_{i}}{h}, \quad D_{-} u_{i}=\frac{u_{i}-u_{i-1}}{h} \\
D_{0} u_{i}=\frac{u_{i+1}-u_{i-1}}{2 h}, \quad D_{+} D_{-} u_{i}=\frac{u_{i+1}-2 u_{i}+u_{i-1}}{h^{2}}
\end{gathered}
$$

For problem ( $B L$ ) the following two discretizations are considered:

$$
\begin{array}{cc}
L_{h}^{1} u_{i} \equiv-\varepsilon D_{+} D_{-} u_{i}+D_{-} u_{i} & \text { (upwind differencing) } \\
L_{h}^{2} u_{i} \equiv \frac{-\varepsilon}{1+h / 2 \varepsilon} D_{+} D_{-} u_{i}+D_{-} u_{i} & \text { (see Kellogg and Tsan). }
\end{array}
$$

Since problems ( $T P-1$ ) and ( $T P-2$ ) have turning points at $x=\frac{1}{2}$ it is necessary to change the direction of the discretization of the first order term. For problem ( $T P-1$ ) the discretization tested is

$$
\begin{align*}
L_{h}^{3} u_{i} \equiv-\varepsilon D_{+} D_{-} u_{i}+\left(x_{i}-\frac{1}{2}\right) D_{+} u_{i}, & 1 \leq i \leq \frac{N+1}{2}  \tag{5}\\
L_{h}^{3} u_{i} \equiv-\varepsilon D_{+} D_{-} u_{i}+\left(x_{i}-\frac{1}{2}\right) D_{-} u_{i}, & \frac{N+1}{2}+1 \leq i \leq N . \tag{6}
\end{align*}
$$

Similarly, for problem (TP-2) the discretization is

$$
\begin{align*}
L_{h}^{4} u_{i} \equiv-\varepsilon D_{+} D_{-} u_{i}-\left(x_{i}-\frac{1}{2}\right) D_{-} u_{i}, & 1 \leq i \leq \frac{N+1}{2}  \tag{7}\\
L_{h}^{4} u_{i} \equiv-\varepsilon D_{+} D_{-} u_{i}-\left(x_{i}-\frac{1}{2}\right) D_{+} u_{i}, & \frac{N+1}{2}+1 \leq i \leq N . \tag{8}
\end{align*}
$$

Note that each of the discretizations $L_{h}^{k}, k=1, \ldots, 4$ results in a tridiagonal linear system of equations whose coefficient matrix may be denoted by

$$
L_{h}^{k} \equiv\left[\begin{array}{lll}
-\alpha_{i} & \beta_{i} & -\gamma_{i}
\end{array}\right]
$$

with

$$
\alpha_{i} \geq 0, \quad \gamma_{i} \geq 0, \quad \text { and } \quad \beta_{i} \geq \alpha_{i}+\gamma_{i}
$$

Thus each $L_{h}^{k}, k=1, \ldots, 4$ is an M-matrix and the linear system of equations

$$
L_{h}^{k} u_{i}=f_{i}, \quad 1 \leq i \leq N,
$$

$u_{0}$ and $u_{N+1}$ fixed has a solution; see for example Berman and Plemmons [Berm79a].

## 3 Review of the Multigrid Algorithm

The particular multigrid algorithm used to solve (1) has been discussed in detail in Kamowitz and Parter [Kamo85a]; thus only a cursory review is given here. In particular the details of the theory behind the convergence results can be found there. What is important to realize is that although the algorithm and convergence results discussed in the previous work apply to well conditioned two point boundary value problems the same bounds on the rate of convergence apply to the singular perturbation problems discussed here.

In order to completely describe the algorithm a number of spaces and operators need to be defined. For the various spaces choose $g$ nested grids

$$
\Omega_{h} \supset \Omega_{2 h} \supset \ldots \supset \Omega_{p h}
$$

where

$$
p=2^{g-1}
$$

Denote by $S_{k h}$ the space of grid functions defined on $\Omega_{k h}$. From now on the notation $k h$ will be replaced by $k$. Denote by $G_{k}$ the smoother, $I_{k}^{k+1}$ the restriction map, and $I_{k+1}^{k}$ the interpolation map. Associated with each space $S_{k}$ there is also a nonsingular tridiagonal operator

$$
L_{k}: S_{k} \mapsto S_{k},
$$

where $L_{k}$ is of positive type. The operators $L_{k}$ are again denoted

$$
L_{k}=\left[\begin{array}{lll}
-\alpha_{i}^{k} & \beta_{i}^{k} & -\gamma_{i}^{k}
\end{array}\right] .
$$

These operators will be formally defined after the statement of the algorithm.
The following is an outline of the algorithm used. Assuming $U_{1}^{n}$ is the $n$th approximation to the solution of the system of equations

$$
L_{1} U_{1}=F_{1}
$$

algorithm $M G\left(U_{k}^{n}, L_{k}, F_{k}, k\right)$ returns $U_{k}^{n+1}$, the next iteration of the multigrid algorithm. The grid layer is denoted by $k$; set $k=1$ to start the algorithm on the finest grid layer.

Algorithm $M G\left(U_{k}^{n}, L_{k}, F_{k}, k\right)$

1. Coarse Solve: If $k=g$ (coarsest grid) then return

$$
M G \leftarrow L_{g}^{-1} F_{g}
$$

otherwise
2. Smooth: Apply the smoother, $G_{k}$, call the result of this step $\tilde{U}_{k}^{n}$.

If $k=1$ use $U_{1}^{n}$ as the initial guess, otherwise use 0 .
3. Recursively Apply the Algorithm: Set
$U^{n+1} \leftarrow \tilde{U}_{k}^{n}+I_{k+1}^{k} M G\left(0, L_{k+1}, I_{k}^{k+1}\left(F_{k}-L_{k} \tilde{U}_{k}^{n}\right), k+1\right)$.
4. (optional) Smooth Again: Set $U^{n+1} \leftarrow G_{k} U^{n+1}$.
5. Return: Set $M G \leftarrow U^{n+1}$.

As defined algorithm $M G$ is called a symmetric V-cycle if step 4 is used; otherwise algorithm $M G$ is referred to as a slash cycle (following the notation of McCormick and Ruge).

For the smoother, $G$, choose $m$-applications of damped Jacobi iteration with parameter $a$. Formally, repeat for $1 \leq r \leq m$

$$
U_{j}^{r+1} \leftarrow U_{j}^{r}+\frac{1}{1+a}\left(L_{j, j}\right)^{-1}\left(F-L U^{r}\right)_{j} .
$$

The interpolation operator, $I_{k+1}^{k}$, is defined as follows. For points common to $\Omega_{k+1}$ and to $\Omega_{k}$ set

$$
\left[I_{k+1}^{k} U\right]_{2 j}=U_{j}
$$

and at odd (new) points of $\Omega_{k}$ require

$$
\left\{L_{k}\left[I_{k+1}^{k} U\right]\right\}_{2 j-1}=0 .
$$

This results in the explicit system of equations at the point $x_{2 j-1}$

$$
U_{2 j-1}^{k}=\frac{1}{\beta_{2 j-1}^{k}}\left[\alpha_{2 j-1}^{k} U_{j-1}^{k+1}+\gamma_{2 j-1}^{k} U_{j+1}^{k+1}\right] .
$$

The restriction operator, $I_{k}^{k+1}$, is

$$
\left[I_{k}^{k+1} U^{k}\right]_{j}=\frac{1}{2}\left[\frac{\alpha_{2 j}^{k}}{\beta_{2 j-1}^{k}} U_{2 j-1}^{k}+U_{2 j}^{k}+\frac{\gamma_{2 j}^{k}}{\beta_{2 j+1}^{k}} U_{2 j+1}^{k}\right]
$$

Note that if $L_{k}$ is symmetric then

$$
I_{k}^{k+1}=\frac{1}{2}\left(I_{k+1}^{k}\right)^{T}
$$

A fundamental observation due to McCormick and Ruge [McCo82a] in the symmetric case is that $S_{k}$ can be written

$$
\begin{equation*}
S_{k}=\text { Range } I_{k+1}^{k} \oplus \text { Nullspace } I_{k}^{k+1} L_{k} \tag{9}
\end{equation*}
$$

For non-symmetric problems the above decomposition follows directly from the characterization of Range $I_{k+1}^{k}$ and Nullspace $I_{k}^{k+1} L_{k}$; see Kamowitz and Parter.

Finally the coarse grid operators, $L_{k+1}$, are chosen by setting

$$
L_{k+1}=\hat{L}_{k+1} \equiv I_{k}^{k+1} L_{k} I_{k+1}^{k}
$$

A direct computation shows that

$$
L_{k+1}=\left[\begin{array}{lll}
-\alpha_{j}^{k+1} & \beta_{j}^{k+1} & -\gamma_{j}^{k+1}
\end{array}\right]
$$

with

$$
\begin{align*}
\alpha_{j}^{k+1} & =\frac{1}{2}\left[\frac{\alpha_{2 j}^{k} \alpha_{2 j-1}^{k}}{\beta_{2 j-1}^{k}}\right]  \tag{10}\\
\beta_{j}^{k+1} & =\frac{1}{2}\left[\beta_{2 j}^{k}-\frac{\alpha_{2 j}^{k} \gamma_{2 j-1}^{k}}{\beta_{2 j-1}^{k}}-\frac{\gamma_{2 j}^{k} \alpha_{2 j+1}^{k}}{\beta_{2 j+1}^{k}}\right]  \tag{11}\\
\gamma_{j}^{k+1} & =\frac{1}{2}\left[\frac{\gamma_{2 j}^{k} \gamma_{2 j+1}^{k}}{\beta_{2 j+1}^{k}}\right] . \tag{12}
\end{align*}
$$

Note that the choice of $\hat{L}_{k+1}$ is the 'natural' choice; however it is not the only feasible choice for $L_{k+1}$.

### 3.1 Convergence Theory

For completeness some results from [Kamo85a] are repeated.
Let

$$
\epsilon^{n}=U_{t r u e}-U^{n}
$$

be the error in the $n$th iterate. Here $U_{\text {true }}$ corresponds to the true solution of the algebraic problem

$$
L_{h} U_{t r u e}=f
$$

The rate of improvement of the $n$th iterate of algorithm $M G$ is then

$$
\rho^{n}=\frac{\left\|\epsilon^{n}\right\|}{\left\|\epsilon^{n-1}\right\|}
$$

To estimate $\rho$, the asymptotic rate as $n \rightarrow \infty$, one needs to bound the $\rho^{n}$.
For later use define

$$
\|x\|_{L_{h}}^{2} \equiv\left\langle L_{h} x, x\right\rangle=\sum_{j}\left(L_{h} x\right)_{j} x_{j} .
$$

First the two grid process is considered. Given an initial error

$$
\epsilon^{0}=U_{t r u e}-U^{0},
$$

the two grid process yields:

1. Smoothing

$$
\epsilon^{0} \rightarrow \tilde{\epsilon}^{0}=G^{\prime} \epsilon^{0}
$$

where $G^{\prime}$ corresponds to the linear part of the smoother $G$. Note that from (9)

$$
\tilde{\epsilon}^{0}=\eta+I_{2 h}^{h} w_{2 h}
$$

where

$$
\eta \in N u l l \text { space } I_{h}^{2 h} L_{h} \text { and } w_{2 h} \in S_{2 h} .
$$

2. Restricting the residual $r_{h}=L_{h} \tilde{\epsilon}^{0}$ to $\Omega_{2 h}$ yields

$$
R_{2 h}=I_{h}^{2 h} L_{h} \tilde{\epsilon}^{0}=I_{h}^{2 h} L_{h} \eta+I_{h}^{2 h} L_{h} I_{2 h}^{h} w_{2 h}=I_{h}^{2 h} L_{h} I_{2 h}^{h} w_{2 h}
$$

since

$$
\eta \in \text { Nullspace } I_{h}^{2 h} L_{h} .
$$

3. Computing the coarse grid correction directly (as is the case for the two grid algorithm) is equivalent to solving

$$
L_{2 h} \psi_{2 h}=R_{2 h}=\hat{L}_{2 h} w_{2 h}
$$

Here

$$
L_{2 h}=\hat{L}_{2 h}=I_{h}^{2 h} L_{h} I_{2 h}^{h}
$$

so

$$
\psi_{2 h}=w_{2 h} .
$$

4. Finally, correcting $\tilde{\epsilon}^{0}$ using the coarse grid correction yields $\epsilon^{\mathbf{1}}$ :

$$
\begin{aligned}
\epsilon^{1} & =U_{t r u e}-U^{1} \\
& =U_{t r u c}-\left(\tilde{U}^{0}+I_{2 h}^{h} \psi_{2 h}\right) \\
& =\tilde{\epsilon}^{0}-I_{2 h}^{h} \psi_{2 h} \\
& =\eta+I_{2 h}^{h} w_{2 h}-I_{2 h}^{h} \psi_{2 h} \\
& =\eta .
\end{aligned}
$$

Note that if

$$
\tilde{\epsilon}^{0} \in S_{2 h}
$$

then one step of the two grid procedure solves the problem! In general, since $\left\|G^{\prime}\right\| \leq 1$ the rate $\rho^{1}$ satisnies

$$
\rho^{1}=\frac{\left\|\epsilon^{1}\right\|}{\left\|\epsilon^{0}\right\|} \leq \frac{\left\|\epsilon^{1}\right\|}{\left\|\tilde{\epsilon}^{0}\right\|} .
$$

In this case

$$
\epsilon^{1}=\eta
$$

so

$$
\rho^{1}=\frac{\|\eta\|}{\left\|\epsilon^{0}\right\|} .
$$

Notice that the $\eta$ term is related to the action of the smoother while the $I_{2 h}^{h} w_{2 h}$ term is eliminated by the multigrid process itself.

In Kamowitz and Parter [Kamo85a] an explicit decomposition of $S_{h}$ was found in terms of the eigenvalues and eigenvectors of the damped Jacobi scheme. This decomposition was exploited to compute bounds on the rate of convergence.

For the two grid scheme, the bounds on $\rho$, for a given $m$ and $a$, are given in Table 1. Note that

| $m$ | $a$ |  |  |  |  |  |
| :--- | :---: | :---: | :---: | :---: | :---: | :---: |
|  | .333 | .500 | .667 | .750 | 1.00 | 1.333 |
| 1 | .500 | .333 | .400 | .429 | .500 | .572 |
| 2 | .250 | .111 | .160 | .184 | .250 | .326 |
| 3 | .125 | .078 | .088 | .093 | .125 | .187 |
| 4 | .068 | .062 | .068 | .072 | .083 | .109 |

Table 1: Predicted Rates for 2 Grids

| $m$ | $a$ |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | .333 | .500 | .667 | .750 | 1.00 | 1.333 |
| 1 | .633 | .577 | .561 | .561 | .577 | .614 |
| 2 | .435 | .408 | .417 | .424 | .447 | .475 |
| 3 | .336 | .335 | .349 | .357 | .378 | .403 |
| 4 | .283 | .293 | .307 | .314 | .333 | .357 |

Table 2: Predicted Asymptotic Rate
the optimal rate is obtained for $a=.5$. In the succeeding sections the bound in the $a=.5$ case will be obtained experimentally for problems $(B L),(T P-1)$, and (TP-2).

In the case where the number of grids is arbitrary, estimates based on the ideas in [ $\mathrm{McCo82a}$ ] are given in Table 2. As indicated in [Kamo85a] these bounds are not sharp.

## 4 The Singular Perturbation Case

For the study of singular perturbation problems it is necessary that if $L_{h}$ is an M -matrix then the coarse grid operators $L_{k+1}, k=1,2, \ldots$ are also M -matrices. This is necessary to insure that each subproblem

$$
L_{k+1} \psi_{k+1}=f_{k+1}
$$

is solvable. In particular since the smoother $G_{k+1}$ depends on $L_{k+1}$, then if $L_{k+1}$ is an M-matrix, $\left\|G_{k+1}\right\|<1$.

Lemma 4.1 If

$$
L_{k}=\left[\begin{array}{lll}
-\alpha_{j}^{k} & \beta_{j}^{k} & -\gamma_{j}^{k}
\end{array}\right]
$$

is an M-Matrix then

$$
L_{k+1}=\left[\begin{array}{lll}
-\alpha_{j}^{k+1} & \beta_{j}^{k+1} & -\gamma_{j}^{k+1}
\end{array}\right]
$$

is also an $M$-matrix, where $\alpha_{j}^{k+1}, \beta_{j}^{k+1}$ and $\gamma_{j}^{k+1}$ are given by equations (10-12).

Proof: The hypotheses that $L_{k}$ is an M-matrix insures that

$$
\alpha_{j}^{k}>0, \quad \beta_{j}^{k}>0, \quad \text { and } \gamma_{j}^{k}>0
$$

so $\alpha_{j}^{k+1}$ and $\gamma_{j}^{k+1}$ are also positive. In addition, $\beta_{j}^{k+1}$ must satisfy

$$
\beta_{j}^{k+1} \geq \alpha_{j}^{k+1}+\gamma_{j}^{k+1}
$$

Recalling the expressions for $\alpha_{j}^{k}, \beta_{j}^{k}$ and $\gamma_{j}^{k}$ note that this is equivalent to requiring

$$
\beta_{2 j}^{k}-\frac{\alpha_{2 j}^{k} \gamma_{2 j-1}^{k}}{\beta_{2 j-1}^{k}}-\frac{\gamma_{2 j}^{k} \alpha_{2 j+1}^{k}}{\beta_{2 j+1}^{k}} \geq \frac{\alpha_{2 j}^{k} \alpha_{2 j-1}^{k}}{\beta_{2 j-1}^{k}}+\frac{\gamma_{2 j}^{k} \gamma_{2 j+1}^{k}}{\beta_{2 j+1}^{k}}
$$

or,

$$
\beta_{2 j}^{k} \geq \frac{\alpha_{2 j}^{k}}{\beta_{2 j-1}^{k}}\left[\alpha_{2 j-1}^{k}+\gamma_{2 j-1}^{k}\right]+\frac{\gamma_{2 j}^{k}}{\beta_{2 j+1}^{k}}\left[\alpha_{2 j+1}^{k}+\gamma_{2 j+1}^{k}\right]
$$

which follows directly from the fact that $L_{k}$ is diagonally dominant.
Since each of the $L_{k+1}$ is an M-matrix, the theory developed in Section 3.1 can be applied in the singular perturbation case.

## 5 Experimental Results

The algorithm of the previous sections was applied to problems (BL), (TP-1) and (TP-2). Problem ( $D L$ ) exhibits a boundary layer at $x=1$. From a computational point of view the system of linear equations that is solved is well conditioned even for small $\varepsilon$; however the fact that the linear system is not symmetric leads to computational difficulties in computing the experimental rate of convergence.

In the case of problem ( $T P-1$ ) there are two boundary layers; at $x=0$ and at $x=1$. In addition the system of equations that is solved is ill-conditioned.

The discrete problem corresponding to problem (TP-2) is well conditioned. The fact that there is an interior turning point at $x=1 / 2$ does not appear to lead to computational difficulties.

By using the one sided difference schemes discussed in the previous section the linear systems arising from the discretization of the problems are of positive type.

From a practical standpoint, however, what effect does the ill-conditioning have on the observed solution? In particular, how should the rate of convergence (reduction in the error) be measured?

### 5.1 General Remarks Regarding the Experiments

In all the experiments $N$, the number of points on the fine grid, was 63 . The initial guess, $U^{0}$, was constructed so that the initial error could be chosen in advance. In other words,

$$
U_{j}^{0}=u\left(x_{j}\right)-\epsilon_{j} .
$$

For the experiments discussed here

$$
\epsilon_{j}=(-1)^{j} .
$$

The number of smoothing steps for the experiments discussed here was set to one. Experiments were run where the number of smoothing steps was greater than one. When more than one smoothing step was used there was no observed qualitative difference in the behavior of the algorithm compared to the behavior for non-singular perturbation problems. Unless otherwise noted the results are for the two grid case. The computer used was a VAX-11/780 and the tests were run in double precision (roughly 16 decimal digits of accuracy).

The damped Jacobi parameter, $a$, was chosen to be .5. The following heuristic argument due to Brandt [Bran77a] suggests why $a=.5$ is optimal for one smoothing step. His suggestion is to choose $a$ so that the range of eigenvalues which are reduced by the damped Jacobi process is equal to the range which is left alone by the process. As noted in [Kamo85a] the eigenvalues of the damped Jacobi scheme come in pairs $\lambda(\mu)$ and $\hat{\lambda}(\mu)$ where

$$
\lambda(\mu)=\frac{\mu+a}{1+a} \quad \text { and } \quad \hat{\lambda}(\mu)=\frac{a-\mu}{1+a}
$$

and $\mu$ are the eigenvalues of the scheme for $a=0$. The eigenvalues $\mu$ are real and distinct and as $h \rightarrow 0$ they fill out the interval $[0,1]$. Brandt's requirement corresponds to choosing $a$ so that

$$
-\lambda(-1)=\lambda(0)
$$

or in other words to taking $a=1 / 2$. Note that this is equivalent to requiring

$$
-\hat{\lambda}(1)=\hat{\lambda}(0)
$$

Indeed in [Kamo85a] for 1 smoothing step the theoretical and experimental results indicate that this choice of $a$ is optimal.

### 5.2 Results for the Boundary Layer Problem

For problem ( $B L$ ) the solution to the analytic problem can be calculated explicitly. For all $\varepsilon>0$ the solution $u_{\varepsilon}(x)$ is

$$
u_{\varepsilon}(x)=C_{1} e^{x / \varepsilon}+C_{2}
$$

| Grid | 1 | .1 | .01 | .005 | .003 | .001 | 0 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1 | $.11 \times 10^{-2}$ | $.54 \times 10^{-2}$ | $.73 \times 10^{-2}$ | $.86 \times 10^{-2}$ | .011 | .012 |
| 1 | $.62 \times 10^{-3}$ | $.11 \times 10^{-2}$ | $.45 \times 10^{-2}$ | $.18 \times 10^{-1}$ | $.21 \times 10^{-1}$ | $.28 \times 10^{-1}$ | .027 |
| .024 |  |  |  |  |  |  |  |
| 2 | $.25 \times 10^{-2}$ | $.12 \times 10^{-1}$ | $.37 \times 10^{-1}$ | $.44 \times 10^{-1}$ | $.46 \times 10^{-1}$ | .048 | .048 |
| 3 | $.98 \times 10^{-2}$ | $.10 \times 10^{-1}$ |  |  |  |  |  |
| 4 | $.39 \times 10^{-1}$ | $.41 \times 10^{-1}$ | $.80 \times 10^{-1}$ | $.92 \times 10^{-1}$ | $.96 \times 10^{-1}$ | .099 | .1 |
| 5 | .15 | .15 | .19 | .20 | .21 | .21 | .21 |
| 6 | 1 | 1 | 1 | 1 | 1 | 1 | 1 |

Table 3: Condition Number - Problem BL
The constants $C_{1}$ and $C_{2}$ satisfy

$$
C_{1}=\frac{2}{e^{1 / \varepsilon}-1}, \quad C_{2}=1-C_{1}=\frac{e^{1 / \varepsilon}-3}{e^{1 / \varepsilon}-1}
$$

Note that

$$
\lim _{\epsilon \rightarrow 0} C_{1}=0, \quad \lim _{\epsilon \rightarrow 0} C_{2}=1
$$

so

$$
\lim _{\varepsilon \rightarrow 0} u_{\varepsilon}(x)=1 .
$$

Denote by $U_{h, \varepsilon}(x)$ the numerical solution of problem ( $B L$ ) for a fixed $h$ and $\varepsilon$. The consistency condition on the solution then requires that for fixed $\varepsilon$

$$
\lim _{h \rightarrow 0} U_{h, \epsilon}(x)=u_{\varepsilon}(x) .
$$

Dorr [Dorr70a] proved that for fixed $h$

$$
\lim _{\varepsilon \rightarrow 0} U_{h, e}(x)=1 .
$$

The discretization of problem ( $B L$ ) using $L_{h}^{1}$ (standard upwind differencing) results in an approximation to the true solution which is an $O(h)$ approximation. To improve on this approximation Kellogg and Tsan point out that the use of $L_{h}^{2}$ results in an $O\left(h^{2}\right)$ approximation for $\varepsilon>0$. For the reduced problem $(\varepsilon=0) L_{h}^{2}$ gives an $O(h)$ approximation.

It is important to note that the linear system of equations arising from the discretization of Problem ( $B L$ ) is not ill-conditioned. Table 3 displays the LINPACK [Dong79a] estimate RCOND (an estimate of the inverse of the condition number) for $L_{h}^{1}, h=1 / 64$ and $\varepsilon=1, .01, .005, .003$, and .001 .

First, some general conclusions about the experimental results. For both $L_{h}^{1}$ and $L_{h}^{2}$ the algorithm converged to the solution of the algebraic problem with an asymptotic convergence rate
identical to the rate predicted in Section 3.1. Moreover there was no observed qualitative difference in the behavior of the algorithm with respect to the choice of $L_{h}^{1}$ or $L_{h}^{2}$. However, as $\varepsilon \rightarrow 0$ the behavior of the iterates changes dramatically.

Define

$$
R_{n}=\frac{\left\|r_{n}\right\|}{\left\|r_{n-1}\right\|}
$$

where

$$
r_{n}=F-L_{h} U^{n}
$$

is the residual after the $n$th iteration and the norm used is the $l_{2}$ norm. From a computational point of view this is a convenient measure of the rate of improvement since one does not know the true solution. From past computational experience this ratio is bounded above by the theoretical error reduction rate. Indeed for large $\varepsilon$, say $\varepsilon=1$, this is the case. For small $\varepsilon$, say $\varepsilon=.001$, after a small number of iterations $R_{n}$ exceeded the rate predicted in Table 1 , and then as $n$ increased $R_{n}$ declined towards the predicted rate; see the solid line in Figure 1. The dashed line will be referred to later. In both cases two grids are used.

An explanation for this behavior is that as the algorithm proceeds the fact that the error is being measured in an asymmetric norm causes problems. The user of the algorithm should be careful to note that while in principle all norms on finite dimensional spaces are equivalent the use of a symmetric norm results in much better observed behavior of this particular algorithm.

To demonstrate this hypothesis the problem

$$
L_{h} U=F
$$

is transformed into the equivalent symmetric problem

$$
\begin{equation*}
D^{-1} L_{h} D\left(D^{-1} U\right)=D^{-1} F=\hat{F} . \tag{13}
\end{equation*}
$$

The matrix $D$ is a positive diagonal matrix whose entries are given by

$$
d_{1}=1, \quad d_{i}=\sqrt{\frac{\alpha_{i}}{\gamma_{i-1}}} d_{i-1}
$$

Applying this transformation to $L_{h}$ results in $L_{h}^{s}$ which is now symmetric. The matrix $L_{h}^{s}$ is denoted

$$
L_{h}^{s}=\left[\begin{array}{lll}
-\alpha_{i} d_{i-1} / d_{i} & \beta_{i} & -\gamma_{i} d_{i+1} / d_{i}
\end{array}\right]
$$

Denote by algorithm $M G^{s}$ algorithm MG applied to problem (13). The error resulting from applying algorithm $M G^{s}$ is related to the error from applying algorithm $M G$ as follows:


Figure 1: $R_{n}$, for $\varepsilon=.001$

Lemma 5.1 If $\epsilon^{n}$ is the error after the n th iteration of applying algorithm $M G$ then the new error, $\epsilon^{n+1}$, is related to the error in the symmetric problem by

$$
\left(M G^{s}\right)\left(D^{-1} \epsilon\right)=D^{-1} M G \epsilon=D^{-1} \epsilon^{n+1}
$$

Proof: By the definition of $\epsilon^{n}$,

$$
\epsilon^{n}=U-U^{n}
$$

and

$$
\epsilon^{n+1}=U-U^{n+1}=M G \epsilon^{n}
$$

Also,

$$
D^{-1} \epsilon^{n}=D^{-1} U-D^{-1} U^{n}
$$

Applying each step of algorithm $M G^{s}$ to $D^{-1} U^{n}$ results in $D^{-1} U^{n+1}$ as a straightforward calculation shows. Hence

$$
M G^{s} D^{-1} U^{n}=D^{-1} M G U^{n}
$$

Applying this transformation to the error $\epsilon^{n}$ and computing the rates

$$
D^{-1} R_{n}=\frac{\left\|D^{-1} \epsilon^{n}\right\|_{L_{h}}}{\left\|D^{-1} \epsilon^{n-1}\right\|_{L_{h}}}
$$

results in the more usual behavior displayed by the dashed line in Figure 1. Figure 2 displays the location where the maximum norm of the error is taken on versus the iteration number. Notice how as the iteration proceeds the location where the maximum occurs 'drifts.' Applying $D^{-1}$ and then measuring $R_{n}$ has the effect of 'fixing' the location where the maximum error is taken on.

Unfortunately the entries of the transformation matrix $D$ satisfy

$$
D_{i} \rightarrow \infty
$$

as $\varepsilon \rightarrow 0$ so this is not a stable method for solving this problem in general.
In summary, for problem (2) where the coefficient matrix was not poorly conditioned but was non-symmetric the only computational difficulty was in choosing the norm in which to measure the rate of improvement of the algorithm.

### 5.3 Turning point Problem - 1

In this section consider the solution of

$$
\begin{equation*}
-\varepsilon u_{\varepsilon}^{\prime \prime}+\left(x-\frac{1}{2}\right) u_{\varepsilon}^{\prime}=0 \tag{14}
\end{equation*}
$$



Figure 2: Index of Maximum Error, Problem BL

| Grid | $\varepsilon$ |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1 | .1 | .01 | .005 | .003 | .001 |  |
| 1 | .00023 | .00032 | $.20 \times 10^{-6}$ | $.19 \times 10^{-9}$ | $.13 \times 10^{-12}$ | $.71 \times 10^{-21}$ |  |
| 2 | .0023 | .0013 | $.76 \times 10^{-6}$ | $.66 \times 10^{-9}$ | $.41 \times 10^{-12}$ | $.23 \times 10^{-18}$ |  |
| 3 | .0096 | .0085 | $.83 \times 10^{-2}$ | $.85 \times 10^{-2}$ | $.80 \times 10^{-2}$ | $.60 \times 10^{-2}$ |  |
| 4 | .038 | .037 | $.46 \times 10^{-1}$ | $.55 \times 10^{-1}$ | $.60 \times 10^{-1}$ | $.65 \times 10^{-1}$ |  |
| 5 | .15 | .146 | .16 | .17 | .18 | .19 |  |
| 6 | 1 | 1 | 1 | 1 | 1 | 1 |  |

Table 4: Condition Number - Problem TP-1
with

$$
u_{\varepsilon}(0)=1, \quad u_{\varepsilon}(1)=3 .
$$

This problem has two boundary layers, at $x=0$ and at $x=1$. The asymptotic solution satisfies (see [Krei74a] and [Dorr70a])

$$
\lim _{\varepsilon \rightarrow 0} u_{\varepsilon}(x)=2 .
$$

For the discrete problem again

$$
\begin{equation*}
L_{h}^{3} U=F \tag{15}
\end{equation*}
$$

is solved where $L_{h}^{3}$ is given by equations (5) and (6). The condition number of $L_{h}^{3}$ tends to infinity as $\varepsilon \rightarrow 0$ for a fixed $h$. Table 4 displays the LINPACK estimate RCOND for $L_{h}^{3}, h=1 / 64$ and $\varepsilon=1$, $.01, .005, .003$, and .001 . It is important to note that although the operator corresponding to the original discretization on the $h$-grid is ill-conditioned, $L_{k}, k=3,4, \ldots$ are not ill-conditioned. In fact in the limiting case where the coarse grid has one point the system of equations to be solved is a $1 \times 1$ system which has condition number 1 . This feature of the multigrid algorithm is reassuring to the user since it means that the actual system being solved directly in step 1 of algorithm $M G$ is well conditioned and no special measures need to be taken.

Although the coarse grids used in the algorithm do not resolve the boundary layers the algorithm still converges. This is because the role of the coarse grids is to solve for the error in the solution, rather than to represent the solution itself. The boundary layers are not seen in the expression for the error. The observed rate of convergence was independent of the value of $\varepsilon$ used.

For all tested values of $\varepsilon$ the multigrid algorithm converged to the solution computed by Gaussian Elimination to the algebraic problem. However, the ill-conditioning of the algebraic system resulting from the discretization of the continuous problem opens up the question of to what solution does the multigrid algorithm converge? Since the original system is ill-conditioned there is a family of solutions $U$ for which the residual is small. Indeed for $\varepsilon=.001$ the condition number

| Grid | $\varepsilon$ |  |  |  |  |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | 1 | .1 | .01 | .005 | .003 | .001 | 0 |  |
| 1 | $.64 \times 10^{-3}$ | $.98 \times 10^{-3}$ | $.45 \times 10^{-2}$ | $.69 \times 10^{-2}$ | $.89 \times 10^{-2}$ | .013 | 0 |  |
| 2 | $.26 \times 10^{-2}$ | $.39 \times 10^{-2}$ | $.17 \times 10^{-1}$ | $.24 \times 10^{-1}$ | $.28 \times 10^{-1}$ | .031 | 0 |  |
| 3 | $.99 \times 10^{-2}$ | $.11 \times 10^{-1}$ | $.25 \times 10^{-1}$ | $.33 \times 10^{-1}$ | $.37 \times 10^{-1}$ | .038 | .035 |  |
| 4 | .039 | .040 | .063 | .080 | .089 | .095 | .096 |  |
| 5 | .15 | .15 | .18 | .20 | .21 | .22 | .22 |  |
| 6 | 1 | 1 | 1 | 1 | 1 | 1 | 1 |  |

Table 5: Condition Number - Problem TP-2
estimate is $.71 \times 10^{-21}$ which is less than the unit roundoff of the VAX-11/780. Thus there is no reason to expect that the solution calculated by solving equation (15) is an adequate approximation to the solution of the continuous problem. In the experimental runs that was indeed the case and both algorithm $M G$ and Gaussian Elimination returned 0 for the solution.

Since for this problem the error is not skewed as in problem ( $B L$ ), it is not surprising to see that the ratio of improvements $R_{n}$ are bounded above by the predicted rate from Table 1. This can be seen in Figure 3, for $\varepsilon=.001$.

### 5.4 Turning point Problem - 2

In this section the solution to

$$
\begin{equation*}
-\varepsilon u_{\varepsilon}^{\prime \prime}-\left(x-\frac{1}{2}\right) u_{\varepsilon}^{\prime}=0 \tag{16}
\end{equation*}
$$

is discussed. The asymptotic solution, see [Kreī $\mathbf{4 a}$ ] is

$$
\begin{array}{ll}
u_{\varepsilon}(x)=1, & 0 \leq x \leq \frac{1}{2} \\
u_{\varepsilon}(x)=3, & \frac{1}{2} \leq x \leq 1
\end{array}
$$

Although one sided differencing, in this case $L_{h}^{4}$, is used to compute the solution, the system of equations that is solved is not ill-conditioned. As in Section 5.3 the LINPACK estimate RCOND of the condition number was computed; see Table 5.

It appears that the condition number is related to the boundary data. For problem (TP-1) the discretization uses information from the boundary layer to estimate $u_{\varepsilon}(x)$ at interior points. The boundary data for problem (TP-2) is well represented (no boundary layers) in the interior.

From the standpoint of the multigrid algorithm again as the grids get coarser the condition number improves. Also, as for problem (TP-1) the rate of convergence of the algorithm was


Figure 3: $R_{n}$ - Problem TP-1, $\varepsilon=.001$
independent of the value of $\varepsilon$, and was identical to the rate predicted in Section 3.1. In addition $R_{n}$ was bounded above by the predicted rate.

## 6 Comparison of $I_{h}^{2 h} L_{h} I_{2 h}^{h}$ and $L_{2 h}$

From a computational point of view one would like to choose for the coarse grid operator $L_{2 h}$ the tridiagonal operator analogous to $L_{h}$ obtained by finite differences. This choice for $L_{2 h}$ will be denoted $L_{2 h}^{f d}$. For now only the symmetric problem

$$
\begin{equation*}
-\left(p u^{\prime}\right)^{\prime}=f, \quad p(x) \geq p_{0}>0 \tag{17}
\end{equation*}
$$

with boundary conditions

$$
u(0)=u(1)=0
$$

is considered. Here $p(x)$ is a smooth function.
Lemma 6.1 The operator $L_{2 h}^{f d}$ is related to $\hat{L}_{2 h}$ by

$$
\hat{L}_{2 h}=L_{2 h}^{f d}+h^{2}\left[\begin{array}{lll}
-A_{j} & B_{j} & -C_{j}
\end{array}\right]
$$

where $A_{j}, C_{j}$ depend on $p(x), p^{\prime}(x)$ and $p^{\prime \prime}(x)$, and $B_{j}=A_{j}+C_{j}$.
Proof: Recall that the coefficients of $\hat{L}_{2 h}$ are given in equations (10-12). In particular

$$
\hat{\alpha}_{k}=\frac{1}{2}\left[\frac{\alpha_{2 k} \alpha_{2 k-1}}{\beta_{2 k-1}}\right] .
$$

Expanding

$$
\begin{aligned}
\alpha_{2 k} & =p\left(\frac{4 k-1}{2} h\right), \\
\alpha_{2 k-1} & =p\left(\frac{4 k-3}{2} h\right) \quad \text { and } \\
\beta_{2 k-1} & =p\left(\frac{4 k-1}{2} h\right)+p\left(\frac{4 k-3}{2} h\right)
\end{aligned}
$$

in Taylor series about the point $\bar{x}=(4 k-1) h$ and collecting terms yields

$$
\begin{aligned}
\hat{\alpha}_{k} & =\frac{1}{2} \frac{\alpha_{2 k} \alpha_{2 k-1}}{\beta_{2 k-1}} \\
& =\frac{1}{2} \frac{p\left(\left(\frac{4 k-1}{2}\right) h\right) p\left(\left(\frac{4 k-3}{2}\right) h\right)}{p\left(\left(\frac{4 k-1}{2}\right) h\right)+p\left(\left(\frac{4 k-3}{2}\right) h\right)}
\end{aligned}
$$

$$
\begin{aligned}
& =\frac{1}{2} \frac{\left[p(\bar{x})+\frac{h}{2} p^{\prime}(\bar{x})+\frac{h^{2}}{8} p^{\prime \prime}(\bar{x})+C_{1} h^{3}\right]\left[p(\bar{x})-\frac{h}{2} p^{\prime}(\bar{x})+\frac{h^{2}}{8} p^{\prime \prime}(\bar{x})-C_{2} h^{3}\right]}{\left[p(\bar{x})+\frac{h}{2} p^{\prime}(\bar{x})+\frac{h^{2}}{8} p^{\prime \prime}(\bar{x})+C_{1} h^{3}\right]+\left[p(\bar{x})-\frac{h}{2} p^{\prime}(\bar{x})+\frac{h^{2}}{8} p^{\prime \prime}(\bar{x})-C_{2} h^{3}\right]} \\
& =\frac{1}{2} \frac{p(\bar{x})^{2}+h^{2}\left[\frac{1}{4} p^{\prime \prime}(\bar{x}) p(\bar{x})-\frac{1}{4}\left(p^{\prime}(\bar{x})\right)^{2}\right]+D h^{3}}{2 p(\bar{x})+\frac{h^{2}}{4} p^{\prime \prime}(\bar{x})+E h^{3}} \\
& =\frac{1}{4} p(\bar{x})+h^{2} A_{k} \\
& =\alpha_{k}^{f d}+h^{2} A_{k}
\end{aligned}
$$

where $A_{k}$ depends on $p(\bar{x}), p^{\prime}(\bar{x})$ and $p^{\prime \prime}(\bar{x})$.
A similar calculation yields

$$
\hat{\gamma}_{k}=\gamma_{k}^{f d}+h^{2} C_{k} .
$$

In addition

$$
\begin{aligned}
\hat{\beta}_{k} & =\hat{\alpha}_{k}+\hat{\gamma}_{k} \\
& =\alpha_{k}^{f d}+h^{2} A_{k}+\gamma_{k}^{f d}+h^{2} C_{k} \\
& =\beta_{k}^{f d}+h^{2} B_{k} .
\end{aligned}
$$

The error analysis in Section 3.1 shows that if $\hat{L}_{2 h}$ is not used then the coarse grid correction $\psi_{2 h}$ will not equal $w_{2 h}$. Thus a bound on the eigenvalues of the eigenvalue problem

$$
\lambda L_{2 h}^{f d} \psi=\hat{L}_{2 h} \psi
$$

is needed.
Lemma 6.2 The eigenvalues $\lambda$ of

$$
\lambda L_{2 h}^{f d} \psi=\hat{L}_{2 h} \psi
$$

satisfy

$$
|\lambda-1| \leq \frac{\max \left|C_{i}\right|}{\min \left|\hat{\gamma}_{i}\right|} h^{2}<\infty .
$$

Proof: By Lemma 6.1

$$
\hat{L}_{2 h}=L_{2 h}^{f d}+h^{2}\left[\begin{array}{lll}
-A_{j} & B_{j} & -C_{j}
\end{array}\right],
$$

thus

$$
(\lambda-1)\left\langle L_{2 h}^{f d} \psi, \psi\right\rangle=h^{2}\left\langle\left[-A_{i} \quad B_{i} \quad-C_{i}\right] \psi, \psi\right\rangle .
$$

Summation by parts yields

$$
\begin{aligned}
\lambda-1 & =h^{2} \frac{\left\langle\left[-A_{i} B_{i}-C_{i}\right] \psi, \psi\right\rangle}{\left\langle L_{2 h}^{f d} \psi, \psi\right\rangle} \\
& =h^{2} \frac{\sum C_{i}\left(\psi_{i}-\psi_{i+1}\right)^{2}}{\sum \gamma_{i}^{f d}\left(\psi_{i}-\psi_{i+1}\right)^{2}} \\
& \leq h^{2} \frac{\max \left|C_{i}\right|}{\min \left|\gamma_{i}^{f d}\right|} \\
& =K h^{2} \\
& <\infty
\end{aligned}
$$

because of the smoothness of $p$ and since $\left|\gamma_{i}^{f d}\right| \geq p_{0}>0$. In addition, since $L_{2 h}^{f d}$ is positive definite,

$$
(\lambda-1)\left\langle L_{2 h}^{f d} \psi, \psi\right\rangle>0
$$

which implies

$$
\lambda-1 \geq-K h^{2}
$$

Therefore

$$
|\lambda-1| \leq K h^{2}
$$

Recall from equation (9)

$$
\tilde{\epsilon}^{0}=\eta+I_{2 h}^{h} w_{2 h}
$$

and if $\hat{L}_{2 h}$ is used then

$$
\epsilon^{1}=\eta
$$

In the case where $L_{2 h}^{f d}$ is used to solve

$$
L_{2 h}^{f d} \psi_{2 h}=\hat{L}_{2 h} w_{2 h}
$$

the error is denoted $\bar{\epsilon}^{1}$. Now

$$
\bar{\epsilon}^{1}=\eta+I_{2 h}^{h}\left(w_{2 h}-\psi_{2 h}\right)
$$

Since Range $I_{2 h}^{h}$ and $N u l l s p a c e ~ I_{h}^{2 h} L_{h}$ are $L_{h}$-orthogonal (see [Kamo85a]),

$$
\left\|\bar{\epsilon}^{1}\right\|_{L_{h}}^{2}=\|\eta\|_{L_{h}}^{2}+\left\|I_{2 h}^{h}\left(w_{2 h}-\psi_{2 h}\right)\right\|_{L_{h}}^{2}
$$

Applying Lemmas 6.1 and 6.2,

$$
\begin{aligned}
\left\|I_{2 h}^{h}\left(w_{2 h}-\psi_{2 h}\right)\right\|_{L_{h}}^{2} & =\left\langle L_{h} I_{2 h}^{h}\left(w_{2 h}-\psi_{2 h}\right), I_{2 h}^{h}\left(w_{2 h}-\psi_{2 h}\right)\right\rangle \\
& =\frac{1}{2}\left\langle\hat{L}_{2 h}\left(w_{2 h}-\psi_{2 h}\right),\left(w_{2 h}-\psi_{2 h}\right)\right\rangle \\
& =\frac{1}{2}\left\langle\hat{L}_{2 h}\left(w_{2 h}-\left(L_{2 h}^{f d}\right)^{-1} \hat{L}_{2 h} w_{2 h}\right),\left(w_{2 h}-\left(L_{2 h}^{f d}\right)^{-1} \hat{L}_{2 h} w_{2 h}\right)\right\rangle \\
& =\frac{1}{2}\left\langle F \hat{L}_{2 h}^{1 / 2} w_{2 h}, F \hat{L}_{2 h}^{1 / 2} w_{2 h}\right\rangle
\end{aligned}
$$

where

$$
F=\left(I-\hat{L}_{2 h}^{1 / 2}\left(L_{2 h}^{f d}\right)^{-1} \hat{L}_{2 h}^{1 / 2}\right)
$$

The eigenvalues of $\hat{L}_{2 h}^{1 / 2}\left(L_{2 h}^{f d}\right)^{-1} \hat{L}_{2 h}^{1 / 2}$ are the same as the eigenvalues of $\left(L_{2 h}^{f d}\right)^{-1} \hat{L}_{2 h}$, so by Lemma 6.2

$$
\begin{aligned}
\frac{1}{2}\left\langle F \hat{L}_{2 h}^{1 / 2} w_{2 h}, F \hat{L}_{2 h}^{1 / 2} w_{2 h}\right\rangle & \leq \frac{1}{2}\left(K h^{2}\right)^{2}\left\langle\hat{L}_{2 h}^{1 / 2} w_{2 h}, \hat{L}_{2 h}^{1 / 2} w_{2 h}\right\rangle \\
& =\frac{1}{2}\left(K h^{2}\right)^{2}\left\langle\hat{L}_{2 h} w_{2 h}, w_{2 h}\right\rangle \\
& =\left(K h^{2}\right)^{2}\left\langle L_{h} I_{2 h}^{h} w_{2 h}, I_{2 h}^{h} w_{2 h}\right\rangle .
\end{aligned}
$$

Therefore

$$
\| I_{2 h}^{h}\left(w_{2 h}-\psi_{2 h}\left\|_{L_{h}}^{2} \leq\left(K h^{2}\right)^{2}\right\| I_{2 h}^{h} w_{2 h} \|_{L_{h}}^{2}\right.
$$

and

$$
\left\|\bar{\epsilon}^{1}\right\|_{L_{h}}^{2} \leq\|\eta\|_{L_{h}}^{2}+\left(K h^{2}\right)^{2}\left\|I_{2 h}^{h} w_{2 h}\right\|_{L_{h}}^{2}
$$

Since $\|G\|_{L_{h}} \leq 1$,

$$
\left\|\bar{\epsilon}^{1}\right\|_{L_{h}} \leq\|\eta\|_{L_{h}}+K h^{2}\left\|\epsilon^{0}\right\|_{L_{h}}
$$

Finally, the rate using $L_{2 h}^{f d}$ is

$$
\begin{aligned}
\rho^{f d} & =\frac{\left\|\bar{\epsilon}^{1}\right\|_{L_{h}}}{\left\|\epsilon^{0}\right\|_{L_{h}}} \\
& \leq \frac{\|\eta\|_{L_{h}}}{\left\|\epsilon^{0}\right\|_{L_{h}}}+K h^{2}
\end{aligned}
$$

Recall the rate using $\hat{L}_{2 h}$ is

$$
\rho \leq \frac{\|\eta\|_{L_{h}}}{\left\|\epsilon^{0}\right\|_{L_{h}}}
$$

$$
\rho^{f d} \leq \rho+K h^{2} .
$$

Remark: Although the discussion in the previous section is restricted to symmetric problems, the same result extends to the general two point boundary value problem

$$
\begin{equation*}
-\left(p u^{\prime}\right)^{\prime}+b u^{\prime}+c u=f, \quad p(x) \geq p_{0}>0, c(x) \geq 0 \tag{18}
\end{equation*}
$$

since problem (18) can be rewritten

$$
\begin{equation*}
-\left(\bar{p} u^{\prime}\right)^{\prime}+\bar{c} u=\bar{f} \tag{19}
\end{equation*}
$$

where

$$
\begin{align*}
\bar{p} & =g(x)  \tag{20}\\
\bar{c} & =g(x) \frac{c(x)}{p(x)}  \tag{21}\\
\bar{f} & =g(x) \frac{f(x)}{p(x)} . \tag{22}
\end{align*}
$$

The function $g(x)$ is

$$
g(x)=\exp \left\{-\int \frac{b(x)-p^{\prime}(x)}{p(x)} d x\right\} .
$$

Note that $\bar{p}, \bar{c}$ and $\bar{f}$ are defined since

$$
p(x) \geq p_{0}>0 .
$$

### 6.1 An Illustration

As an illustration of the effect of replacing $\hat{L}_{2 h}$ with $L_{2 h}^{f d}$ consider using for the smoother $G$ in algorithm $M G$ one sweep of odd Gauss-Seidel smoothing. In other words

$$
(G U)_{i}=U_{i} \quad \text { for } i=0 \bmod 2
$$

and for $i=1 \bmod 2$

$$
(G U)_{i}=\frac{1}{\beta_{i}}\left[\alpha_{i} U_{i-1}+\gamma_{i} U_{i+1}+f_{i}\right] .
$$

This guarantees that the error $\tilde{\epsilon}$ lies completely in Range $I_{2 h}^{h}$. In other words, the $\eta$ term is 0 since for $i=0 \bmod 2$

$$
\left(L_{h} I_{2 h}^{h} w_{2 h}\right)_{i} \neq 0
$$

yet

$$
\begin{aligned}
0 & =\left\langle L_{h} \eta, I_{2 h}^{h} w_{2 h}\right\rangle \\
& =\left\langle\eta, L_{h} I_{2 h}^{h} w_{2 h}\right\rangle .
\end{aligned}
$$

In addition for points $x_{i}$, with $i=1 \bmod 2$,

$$
\begin{aligned}
0=r_{i}=L_{h} \bar{\epsilon} & =\left(L_{h} \eta\right)_{i}+\left(L_{h} I_{2 h}^{h} w_{2 h}\right)_{i} \\
& =\left(L_{h} \eta\right)_{i}
\end{aligned}
$$

which implies $\eta_{i}=0$ since $\eta_{i-1}=\eta_{i+1}=0$.
Thus

$$
\tilde{\epsilon}=I_{2 h}^{h} w_{2 h}
$$

for some $\boldsymbol{w}_{\mathbf{2 h}}$. If $\hat{L}_{\mathbf{2 h}}$ is used then one iteration of the algorithm results in completely solving the problem.

For a numerical example problem (17) is solved with

$$
p(x)=e^{\cos (\pi x)}
$$

The right hand side $f$ is chosen so that the the true solution is

$$
u_{\text {true }}(x)=\sin (\pi x) .
$$

As expected when $\hat{L}_{2 h}$ is used the algorithm converges in one step. When $L_{2 h}^{f d}$ is used then the observed rate of convergence corresponds to the $K h^{2}$ term of Section 6. Table 6 displays the observed rate of convergence versus $h$ for four different choices for the number of grids. The column $\rho(h)$ corresponds to the observed rate of convergence for each value of $h$. The value of $\delta(h)$ is

$$
\delta(h)=\frac{\rho(2 h)}{\rho(h)} .
$$

Since the error in the two grid case is $O\left(h^{2}\right)$ one expects that $\delta(h) \rightarrow 4$ as $h \rightarrow 0$ for two grids. Indeed this is the case. The reason why $\rho(h)$ varies with the number of grids used is that there are 'pollution' effects from not solving each coarse grid equation exactly. In other words

$$
\epsilon^{1}=\eta_{h}+I_{2 h}^{h} \eta_{2 h}+\ldots+O\left(h^{2}\right)
$$

where the $O\left(h^{2}\right)$ term corresponds to the error made by not using $\hat{L}_{2 h}$.

|  | 2 grids |  | 3 grids |  | 4 grids |  | 5 grids |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $h$ | $\rho(h)$ | $\delta(h)$ | $\rho(h)$ | $\delta(h)$ | $\rho(h)$ | $\delta(h)$ | $\rho(h)$ | $\delta(h)$ |
| $1 / 16$ | $.50 \times 10^{-2}$ |  | $.27 \times 10^{-1}$ |  | $.85 \times 10^{-1}$ |  |  |  |
| $1 / 32$ | $.12 \times 10^{-2}$ | 4.17 | $.62 \times 10^{-2}$ | 4.35 | $.28 \times 10^{-1}$ | 3.04 | $.86 \times 10^{-1}$ |  |
| $1 / 64$ | $.30 \times 10^{-3}$ | 4.00 | $.15 \times 10^{-2}$ | 4.13 | $.65 \times 10^{-2}$ | 4.31 | $.29 \times 10^{-1}$ | 3.01 |
| $1 / 127$ | $.75 \times 10^{-4}$ | 4.00 | $.38 \times 10^{-3}$ | 3.95 | $.16 \times 10^{-2}$ | 4.06 | $.66 \times 10^{-2}$ | 4.33 |

Table 6: Rate using $L_{2 h}^{f d}$

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