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## INTRODUCTION

This work describes the development of a computer program for the analysis of the time - dependnet incompressible viscoas flow problems.

Mathematically, the problem is that of solving numerically a partial differential equation in three variables containg non-linear terms. The most natural boundary conditions to impose, those in which velocities are prescribed at the boundary, are also among those which are the most difficult to handle computationally. These various difficulties are illustrated in the recent work. It was found necessary to approximate the boundary conditions in a way which affected accuracy, and also to take such small time steps for reasons of stability and accuracy that the computer time become excessive.

The method here depends on the use of the primitive variables - i.e the velocities and the pressure and is applicable to problems in two and three space dimensions. An analytical disscussion of the properties of the method requires a background of numerical analysis for that reason, we have collected the relevant informations, e.g definitions and theormes in chapter I. In chapter II the general method of solving numerically the Navier - Stokes equations for pressure and velocities in Hydrodynamic is presented. Finally chapter III contains a flow chart and Algol 60 Program for solving
the test example given at the end of chapter II. Finally at the end of this work a list of the references used will be given.

## CHAPTER I

## PRELIMINARIES DEFINITIONS AND NOTATIONS

This chapter includes the required theoremes, notations and definitions which we shall need throught this work.

## $\int$ 1. NORMS AND MATRICES

The norm of a matrix is a number assigned to the matrix which is in some sense a measure of the magnitude of the matrix. The norm of $A$, denoted by $\|A\|$ have the following properties.
(1) $\left\|_{A}\right\| \geqq 0,\left\|_{A}\right\|_{=0}$ if and only if $A=0$
(2) $\|C A\|=|C|\left\|_{A}\right\|$ where $C$ is any real number
(3) $\|A+B\| \leqq\left\|_{A}\right\|+\|B\|$
(4) $\|A B\| \leqq\left\|_{A}\right\|\left\|_{B}\right\|$
among the many possible ways of defining $\|$ A $\|$ which satisfy (I-I) we consider:

$$
\begin{aligned}
& \|A\|_{E}=\left(\sum_{i=1}^{m} \sum_{j=1}^{n} a_{i j}\right)^{I / 2}, \text { the Euclidean norm } \\
& \left\|\|_{S}=\max _{i}\left[\lambda_{i}\left(A A^{I}\right)\right]^{1 / 2},\right. \text { the spectral norm }
\end{aligned}
$$

Both of which is defined for any nxnmatrix. In the definition of the spectral norm the notation $\lambda_{i}\left(A A^{T}\right)$ denotes an eigenvalue of $A A^{T}$. For vectors we define the norm in the Euclidian sense as,

$$
\left\|_{x}\right\|_{=\left(x^{T} x\right)^{I / 2}=|x|| | x \mid} \mid
$$

The spectral radius is defined, by,

Thus

$$
\begin{aligned}
& P(A)=\max _{i}\left|\lambda_{i}(A)\right| \begin{array}{l}
\text { the maximum modulus } \\
\text { eigenvalue of the mat- } \\
\text { riv } A
\end{array} \\
& \|A\|=\left[P\left(A A^{T}\right)\right]^{T / 2}
\end{aligned}
$$

Theorein 1.1

If $A$ is the tridiagonal matrix,

$$
\left[\begin{array}{cccc}
a & b & & 0 \\
c & a & b & \\
& \ddots & \ddots & \vdots \\
0 & \ddots & 0 & b \\
0 & c & a
\end{array}\right] \text { where } a, b \text { and } c \text { are }
$$

real and $b c>0$, then the eigenvalues of $A$ are givenby,

$$
\lambda_{m}=a+2 \sqrt{b c} \quad \cos \frac{m \pi}{n+1},(m=1,2, \ldots, n)
$$

## Theorem 1.2

For any matrix $A,\left\|_{A}\right\| \geqq P(A)$, if $A$ is symmetric, $\|A\|=P(A)$
Def. The matrix $A$ is convergent to zero if the sequance of matrices $A, A^{2}, A^{3}, \ldots$ converges to the null matrix 0 .

Theorem 1.3

$$
\operatorname{LimaA}_{r \rightarrow \infty}=0 \text { if }\|A\|<I
$$

proof:

$$
\left\|A^{r}\right\|=\left\|A A^{r-1}\right\| \leqq\|A\|\left\|A^{r-1}\right\|\|A\|^{2}\left\|A^{r-2}\right\| \cdots \triangleq\|A\|^{x}
$$

Theorem 1.4

proof: Consider the Jordan canonical form of $A$. A Jordan submatrix of $A$ is of the form, $\left[\begin{array}{cccc}\lambda_{i} & & & \\ 1 & \lambda_{i} & & \\ & 1 & & \\ & & & \\ & & & \\ 0 & & & 0 \\ & & & \\ \lambda_{i}\end{array}\right]$
where $\lambda_{i}$ is an eigenvalue of $A$. If the Jordan submatrix is raised to the power $r$, then the result tendes to the null matrix as $r \rightarrow \infty$; if and only if $\mid \lambda_{1} k I$. Theorem 1. 5

If $\lambda_{1} ; \lambda_{2}, \ldots, \lambda_{n}$ are the eigenvalues of $A$, then the eigenvalues of $A^{k}$ are $\lambda_{1}^{k}, \lambda_{2}^{k}, \ldots, \lambda_{n}^{k}$, more generally if $p(x)$ is a polynomial, the eigenvalues of $p(A)$ are $p\left(\lambda_{1}\right), \ldots, p\left(\lambda_{n}\right)$ 。

## Theorem I. 6

If $A$ is real and symmetric, all eigenvalues and eigenvectors are real. Moreover, eigenvectors corresponding to distinct eigenvalues are orthogonal and the left eigenvector corresponding to the eigenvector $X_{i}$ is $x_{i}^{T}$. Theorem I. 7
Any similarity transformation $P A P^{-1}$ applied to A leaves the eigenvalues of the matrix unchanged. proof: Let $\lambda$ be an eigenvalue of $A$ and $x$ the associated eigenvector then,

$$
\begin{equation*}
A x=\lambda x, \text { so that }, P A x=\lambda P \mathrm{x} \tag{I-2}
\end{equation*}
$$

Let $\quad y=P x \quad$ so that $\quad X=P^{-1} y$, subistituting in (1.2) gives

$$
\operatorname{PAP}^{-1} y=\lambda y
$$

Thus $\lambda$ is an eigenvalue of $P A P^{-1}$ and $y$ is the assocrated eigenvector.

Theorem I. 8
Let $f(\lambda)=|A-\lambda I|=0$
be the characteristic equation of $A$ then $f(A)=0$

## Theorem 7.9

Given an arbitrary matrix $A$, there exists a non-singular matrix $P$, whose elements may be complex, such that

$$
P^{-1} A P=\left[\begin{array}{cccccc}
J_{1} & 0 & \cdots & \cdots & \cdots & 0 \\
0 & J_{2} & \ddots & & & \vdots \\
\vdots & \ddots & \ddots & & & \\
\vdots & \ddots & \ddots & \ddots & & \vdots \\
\vdots & & \ddots & \ddots & \ddots & \vdots \\
\vdots & & & \ddots & \ddots & 0 \\
0 & & \cdots & \ddots & \ddots & J_{k}
\end{array}\right]
$$

Where $J_{\mathbb{k}}, k=I, \ldots, \mathbb{M} \leqq n$ is a matrix with an eigenvalue $\lambda_{i}$ of $A$ on its main diagonal and I'S on the diagonal above the main diagonal.

Thus

$$
J_{k}=\left[\begin{array}{cccccc}
\lambda_{i} & 1 & 0 & & & 0 \\
0 & \lambda_{i} & 1 & & & 1 \\
\vdots & \ddots & \ddots & & \\
& \ddots & \ddots & \ddots & 0 \\
& & \ddots & \ddots & & 1 \\
& & & \ddots & \ddots & 1 \\
\theta & \cdots & & & 0 & \lambda_{i}
\end{array}\right]
$$

Note that a given eigenvalue may appear as the dagoal element of more than one $J_{k}$. The matrix in (1-3) is called the Jordan canonical form of $A$. The determinants

$$
\left|\left(J_{k}-\lambda I\right)\right|=\left(\lambda_{i}-\lambda\right)^{\nu_{k}}
$$

where $\nu_{k}$ is the order of $J_{k}$ are called the elementary divisors of $A$.
$\oint$ 2. THE SUCCESSIVE OVERRELAXATION METHOD FOR

## SOLVING A SET OF LINEAR ALGEBRIC EQUATIONS

An implicit finite differente formula which approximates a partial differential equation in any number of space variables involves several grid points at the advanced time level. So it is required to find the solustin of the equations which arises there. A set of simultanous linear equations, which can be written in the form.

$$
\begin{equation*}
P \underline{X}=\underline{c}, \quad(|P| \neq 0) \tag{2.1}
\end{equation*}
$$

where $P$ is a square matrix, with no zero on the main diagonal, $x$, $c^{\text {c are vectors requiers to be solved at each }}$ time step.

Equation (2.1) can be written in the form

$$
\begin{equation*}
A \underline{x}=\underline{b}, \quad(|A| \neq 0) \tag{2.2}
\end{equation*}
$$

where $A=D P, \quad \underline{b}=D \underline{C}$
and $D$ is a diagonal matrix chosen so that the elements of the principal diagonal of $A$ are unity. $A$ can be written in the form

$$
A=I-I-U
$$

where $I$ is the unit matrix, and $I$, $U$ are the lower and upper triangular matrices respectively.

An iterative process is used to solve (2-2). We begin by initial value $\underline{x}_{0}$ and is successively improved by the iterative process until it is arbitrary close to X. (2-2) can b written in the form

$$
\begin{equation*}
(I-I-U)_{\underline{x}}=\underline{b} \tag{2-3}
\end{equation*}
$$

and $X_{i}, X_{i+1}$ are successive approximate solutions of equation (2-2), then using of equation (2-3) give,

$$
I \underline{X}_{i+1}=(L+U) \underline{X}_{i}+\underline{b} \quad,(i=1,2, \ldots)
$$

which is the point Jacobi iterative method, or

$$
(I-I) \underline{x}_{i+1}=U \underline{x}_{i}+\underline{b}
$$

which is the Gauss-Seidel iterative method. These two methods are special cases of the general iterative process

$$
\begin{equation*}
\underline{x}_{i+1}=B \underline{x}_{i}+\underline{c},(i=1,2, \ldots) \tag{2-4}
\end{equation*}
$$

where $B=I+U$ and $(I-I)^{-1}$ in the point Jacobi and Gauss-Seidel processes respectively. An error in the $i^{\text {th }}$ iterate is

$$
\underline{e}_{i}=\underline{x}_{i}-\underline{x} \quad,(i=0,1,2, \ldots)
$$

Then

$$
\begin{aligned}
& e_{i+1}+\underline{x}=B \underline{e}_{1}+B \underline{x}+\underline{c} \\
& \underline{e}_{i+1}=B \underline{e}_{i} \quad(i=0,1,2, \ldots)
\end{aligned}
$$

then

$$
\begin{aligned}
& e_{i}=B^{i} \underline{e}_{0}, \text { and } \text { so, } \\
& \underline{e}_{i} \rightarrow 0, \text { as } i \rightarrow \infty, \text { if } B^{i} \rightarrow 0
\end{aligned}
$$

where 0 is the null matrix $B$ is convergent (ie $B^{i} \rightarrow 0$ as $i \rightarrow \infty$ ) if and only if $\rho(B)<1$ (see theorm l-2).

Thus the iteration (2.4) is convergent if and only if $(P(B)<1$ 。

Thus to solve equations (2-2) by the method of successive overrelaxation, we introduce,

$$
\begin{equation*}
\tilde{\underline{X}}_{i+1}=I \underline{X}_{i+1}+U \underline{\underline{x}}_{i}+\underline{b} \tag{2-5}
\end{equation*}
$$

where

$$
\begin{equation*}
\underline{x}_{i+1}=\omega \tilde{\underline{x}}_{i+1}+(1-\omega) \underline{x}_{i} \tag{2-6}
\end{equation*}
$$

where $\omega>0$ an arbitrary parameter, independent of $i$, called the relaxation factor. Elimination of $\tilde{\underline{X}}_{i+1}$ between $(2-5),(2-6)$, gives,

$$
(I-\omega I)_{\underline{x}_{i+1}}=[U+(I-\omega) I] \underline{X}_{i}+\omega \underline{b}
$$

ie

$$
\begin{equation*}
\underline{x}_{i+1}=(I-\omega I)^{-1}[\omega U+(I-\omega) I] \underline{x}_{i}+\omega(I-\omega I)^{-} \underline{b} \tag{2-7}
\end{equation*}
$$

This is an iterative method of successive overrelaxation and simillar to (2.4) with

$$
B \equiv(I-\omega L)^{-I}[\omega U+(I-\omega) I]
$$

so the method of successive overrelazation will be convergent if and only if,

$$
P\left[(I-\omega I)^{-1}\{\omega U+(I-\omega) I\}\right]<I
$$

we can write

$$
\underset{\omega}{\mathrm{H}_{\omega}}=(I-\omega I)^{-1}\{\omega U+(I-\omega) I\}
$$

Thus, if $\lambda$ is an eigenvalue of $\frac{H_{\omega}}{}$, then,

$$
\begin{equation*}
\left|H_{\omega}-\lambda I\right|=0 \tag{2-8}
\end{equation*}
$$

So, we shall calculate the maximum eigenvalue of $\mathrm{H}_{\omega}$ from equation (2-8) and minimize this with respect tow.

Definition 2.1

A matrix is two -Cyclic if by a siutable permutation of its rows and corresponding columns, it can be written in the form,

$$
\left[\begin{array}{ll}
I & F \\
G & I
\end{array}\right]
$$

where $I$ is a square unit matrix, and $F$, G are rectanguar matrices.

## Definition 2.2

A matrix is weakly two-cyclic if by a siutable permutation of its rows and corresponding columns, it can be written in the form

$$
\left[\begin{array}{ll}
0 & F \\
G & 0
\end{array}\right]
$$

where 0 is a square null matrix.
Definition 2.3

If the matrix ( $I-I-U$ ) is two-cyclic, then it is consistently ordered if all the eigenvalues of the matrix

$$
\alpha I+\frac{1}{\alpha} U \quad(\alpha \neq 0)
$$

are independent of $\alpha$.
Thus returning to equation (2-8), it can be written in the form,

$$
\begin{aligned}
& \left|(I-\omega I)^{-I}\{I+\omega(U-I)\}-\lambda I\right|=0 \\
& |\{I+\omega(U-I)\}-\lambda(I-\omega I)|=0 \\
& \left|(U+\omega I)-\frac{\lambda+\omega-I}{\omega} I\right|=0 \\
& \left|\lambda^{I / 2}\left(\lambda^{I / 2} I+\lambda^{-I / 2} U\right)-\frac{\lambda+\omega-I}{\omega} I\right|=0 \\
& \left.\mid \lambda^{I / 2} I+\lambda^{-I / 2} U\right) \left.-\frac{\lambda+\omega-I}{\lambda^{I / 2}} I \right\rvert\,=0
\end{aligned}
$$

If I - I - U is two-cyclic and consistently ordered, then

$$
\left|(I+U)-\frac{\lambda+\omega-1}{\lambda^{1 / 2}} I\right|=0
$$

Thus for any eigenvalue $\lambda$ of the successive overrelaxation $H_{\omega}$, there corresponds an eigenvalue $\mu$ of the point Jacobi matrix $(I+U)$, where

$$
\begin{equation*}
\mu=\frac{\lambda+\omega-1}{\omega \lambda^{1 / 2}} \tag{2-9}
\end{equation*}
$$

Equation (2-9) conects the eigenvalue of the successive overrelaxation matrix with the eigenvalues of the point Jacobi matrix, provided that I- I-U is two -cyclic and consistently ordered.

If the matrix $I-I-U$ is symmetric as well as being two cyclic and consistently ordered, and so ( $I+U$ ) is symmetric and hence the eigenvalues of ( $I+U$ ) are real.


Since ( $I+U$ ) is weakly two-cyclic, its non-zero eigenvalues occure in pairs different in sign. i.e

$$
-P(I+U) \leqq \mu \leqq P(I+U)
$$

interchanging rows and corresponding columns of ( $I+U$ ), it can be written as,

$$
\left[\begin{array}{ll}
O_{1} & E \\
G & O_{2}
\end{array}\right]
$$

where $\mathrm{O}_{1}, \mathrm{O}_{2}$, are square matrices of order $r, s$ respectively and ( $I+U$ ) is square matrix of order ( $r+s$ ). Since the interchanging of rows and columns does not affect the eigenvalues of a matrix, the eigenvalues of $(I+U)$ are given by,

$$
\left[\begin{array}{cc}
-\mu I_{1} & F \\
G & -\mu I_{2}
\end{array}\right]=0
$$

where $I_{1}, F_{2}$ are unit matrices of order $r$ and $s$ respectively. Thus

$$
\left[\begin{array}{cc}
\mu_{I_{I}} & F \\
G & \mu_{I_{2}}
\end{array}\right]=0
$$

by multiplyaing the first $r$ rows and the last $s$ columns of the determinant by -1. This shows also that $-\mu$ is also an eigenvalues of $(I+U)$. We assume that the point Jacobi method is convergent and hence,

$$
0<P(I+U)<I
$$

Also since $I$ - I - U is a consistently ordered twocyclic matrix, the Gauss-Seidel method is also convergent. From equation (2-9) we consider for a given value of $\mu_{\text {in }}$ the range

$$
0<\mu \leqq P(I+U)<I
$$

the two functions of $\lambda$

$$
f(\lambda)=\frac{\lambda+\omega-1}{\omega}, \quad g(\lambda)=\mu \lambda^{1 / 2}
$$

These function can be shown in the figure (1), where $f_{\omega}(\lambda)$ is a straight line passing through ( $1, I$ ) and $g(\lambda)$ is a parabola.


Figure 1
Thus equation (2-9) geometrically represents the intersection of the curves $f_{\omega}(\lambda)$ and $g(\lambda)$ with the two values of $\lambda$ at the points of intersection $A$ and $B$ given by

$$
\lambda^{2}+2\left[(\omega-1)-1 / 2 \mu^{2} \omega^{2}\right] \lambda+(\omega-1)^{2}=0
$$

i.e $\quad \lambda=1 / 22^{2} \omega^{2}-(\omega-1) \pm \mu \omega\left[1 / \pi \mu^{2} \omega^{2}-(\omega-1)\right]^{1 / 2}$

It is clear that the large abscissa of the two points of intersection decreases with increasing $\boldsymbol{\omega}$, until eventually $f_{w}(\lambda)$ becomes a tangent to $g(\lambda)$ at the point C. Thus

$$
1 / 4 \mu^{2} \omega^{2}-\omega+1=0
$$

ie

$$
\begin{aligned}
& \omega=\frac{I_{-}^{+}\left(1-\mu^{2}\right)^{1 / 2}}{1 / 2 \mu^{2}} \\
& \omega=\frac{2}{I_{F}\left(1-\mu^{2}\right)^{1 / 2}}
\end{aligned}
$$

The range of $\omega$ must include $\omega=1$, and so, we have

$$
\begin{equation*}
\tilde{\omega}=\frac{2}{1+\left(1-\mu^{2}\right)^{1 / 2}} \tag{2-10}
\end{equation*}
$$

if $\omega>\widetilde{\omega}, \lambda$ has two cojugate complex roots,

$$
\lambda=\frac{1}{2} \mu^{2} \omega^{2}-(\omega-1) \pm i \mu \omega\left\{(\omega-1)-\frac{1}{4} \mu^{2} \omega^{2}\right\}^{1 / 2}
$$

Thus

$$
|\lambda|=\omega-1
$$

Thus the minimum value of $\lambda$ is $\tilde{\lambda}=\tilde{\omega}-1$
where $\tilde{\omega}$ is given by $(2-10)$ and $\mu$ is the eigenvalue of $(I+U)$ in the range

$$
\begin{array}{ll} 
& 0<\mu \leqq P(I+U)<I, \\
\text { since } & g(\lambda)=P(I+U) \lambda^{I / 2}
\end{array}
$$

is the envelope of all the curves $g(\lambda)=\mu \lambda^{1 / 2}$, where

$$
0<\mu \leqq \rho(I+U)<I
$$

it follows that,

$$
\begin{equation*}
\min _{\omega} P\left(H_{\omega}\right)=P\left(H_{o p t}\right)=\omega_{o p t^{-1}} \tag{2-11}
\end{equation*}
$$

where $\omega_{\text {opt }}$ is given by,

$$
\begin{aligned}
\omega_{o p t} & =\frac{2}{I+\left(I-\mu_{o p t}^{2}\right)^{1 / 2}} \\
\mu_{o p t} & =P(I+U)
\end{aligned}
$$

Thus we found the value of $\boldsymbol{\omega}$, given by (2-11), which minimizes the maximum modulus eigenvalue of $\hat{H}_{\omega^{\circ}}$ Also since the point Jacobi method is convergent if $0<P(I+U)<I$ and so from equation (2-II) it follows that

$$
I<\omega_{o p t}<2
$$

and also from (2-10)

$$
0<P\left({\underset{\omega}{\omega \mathrm{pt}}}^{\mathrm{H}_{\mathrm{p}}}\right)<I
$$

## §3. SOME NOTES ABOUT THE ITERATIVE METHOD FOR

## SOLVING PARTIAL DIFFERENCE EQUATIONS

In the numerical solution by finite differences of boundary value problems involving partial differential equations, one is led to consider linear systems of high order of the form

$$
\sum_{j=1}^{n} a_{i j} u_{j}+d_{i}=0 \quad(i=1,2, \ldots, n) \quad(3-1)
$$

where $u_{1}, u_{2}, \ldots, u_{n}$ are unknown and where the real
numbers $a_{i j}$ and $d_{i}$ are known. The coefficients $a_{i j}$ satisfy the conditions
(a) $\left|a_{i i}\right| \geqq \sum_{j=1, j \neq i}^{n}\left|a_{i j}\right|$, and for some i the
(b) Given any two nonempty disjoint subsets $S$ and $T$ of $W$, the set of the first $n$ po- $(3-2)^{3}$ sitive integers such that $S \mathbb{T}=\mathbb{W}$, there exists $a_{i j} \neq 0$ such that $i \in S$ and $j \in \mathbb{T}$.

I tan be shown that the determinant of the matrix $A=\left(a_{i j}\right)$ does not vanish. Moreover, if the matrix $A^{*}=\left(a_{i j}^{*}\right)$ is symmetric, where $a_{i j}^{*}=\left.a_{i i} a_{i j}\right|_{a_{i i}} \mid$, (i,j $=1,2, \ldots, n$ ), then $A^{*}$ is positive definite. For if $\lambda$ is non positive real numbers, then the matrix * $-\lambda I$, where $I$ is the identity matrix, also satisfies ( $3-2$ ) and hence its determinant can not vanish.

There fore all eigenvalues of $A^{*}$ are positive, and * is positive definite. On the other hand if A. is positive definite then $a_{i i} \neq 0,(i=1,2, \ldots, n)$. An appropriate method for solving equations (3-1) numerically, is that of systematic iteration, which is better for computer. We shall consider linear systems such that either the matrix $\bar{A}$ satisfies conditions $(3-2)$ or such that the matrix $A^{*}$ is positive definite. In order to define the iterative methods it is necessary that $a_{i i} \neq 0(i=1,2, \ldots, n)$, we shall assume that $a_{i i}>0,(i=1,2, \ldots, n)$ also the matrix $A$ has properity (A): there exist two disjoint subsets $S$ and $T$ of $W$, the set of the first $n$ integers, such that $S \mathbb{U}=\mathbb{W}$ and if $a_{i j} \neq 0$ then $i=j$ or $i \in S$ and $j \in T$ or $i \in \mathbb{T}$ and $j \in S$. This is the Younge's condition for the matrix $A$.

A short summery will be given here for the solution of linear systems derived from boundary value-problems, the matrix of which satisfies (3-1) and has property (A).

An iterative method, which converge s fastes than the usual methods will given. We assume that the rows and columns of $A$ are arranged in the ordering $\sigma$.

$$
\begin{array}{r}
u_{i}^{m+1}=\omega\left\{\sum_{j=1}^{i-1} b_{i j} u_{j}^{m+1}+\sum_{j=i+1}^{n} b_{i j} u_{j}^{m}+c_{i}\right\}-(w-1) u_{i}^{m} \\
(m \geqq 0, \quad i=1,2, \ldots, n)(3-3
\end{array}
$$

where $u_{i}$ is arbitrary $(i=1,2, \ldots, n)$, and

$$
b_{i j}= \begin{cases}-a_{i j} / a_{i i} & (i \neq j)  \tag{3-4}\\ 0 & (i=j)\end{cases}
$$

and

$$
c_{i}=-d_{i} / a_{i i} \quad(i=1,2, \ldots, n)
$$

Equation (3-3) can be written in the form

$$
\begin{equation*}
u^{\mathrm{m}+1}=I_{\sigma, \omega}\left[u^{m}\right]+f \quad, m \geqq 0 \tag{3-5}
\end{equation*}
$$

where $u^{m}=\left(u_{1}^{m}, u_{2}^{m}, \ldots, u_{n}^{m}\right), f=\left(f_{1}, f_{2}, \ldots, f_{n}\right), x^{m}$ fixed, and $I_{\sigma, w}$ is a linear operator. denotes the ordering of the equations, and $\omega$ is the relaxation factor. This is the method of successive overrelaxation. As we show in $\} 2$. that if $A$ has property (A), then there exist certain orderings $\sigma$ such that for all $\omega$ a relation holds between the eigenvalues of $I_{\sigma_{0}, w}$ and the eigenvalues of the matrix $B=\left(b_{i j}\right)$ defined by (3-4). If $\mu$ denotes the spectral norm of $B$, i.e the maximum of modulii of the eigenvalues of $B$, then $I_{\sigma, \omega}$ converges if and only if $\mu<I$ (the Gauss-Seidel method). Conditions (3.2) imply $\mu<1$ 。

If $A$ is assumed to be symmetric and have property (A) then $\mu<I$ if and only if $A$ is positive definite. The optimum relaxation factor $\boldsymbol{\omega}_{\text {opt }}$ is given by,

$$
\begin{equation*}
\omega_{\text {opt }}^{2} \mu^{2}-4\left(\omega_{o p t}-1\right)=0 \tag{3-6}
\end{equation*}
$$

$$
\omega_{\mathrm{opt}}>2
$$

For more details and complete proves of the following theormes see [9]. Theorem 3.1

A matrix $A$ has property (A) if and only if there exists a vector $\gamma=\left(\gamma_{1}, \gamma_{2}, \ldots, \gamma_{n}\right)$ with integral components such that if $a_{i j} \neq 0$ and $i \neq j$ then $\left|\gamma_{i}-\gamma_{j}\right|=1$

## Theorem 3.2

Let $A$ be an $n x n$ matrix with property (A) and with a consistent ordering of rows and coulmns. If the alements of $A^{\prime}=\left(a_{i j}\right)$ and $A^{\prime \prime}=\left(a_{i j}^{\prime \prime}\right)$ are defined by

$$
a_{i j}= \begin{cases}a_{i j} & (i \leqq j) \\ a_{i j} & (i>j)\end{cases}
$$

$$
a_{i j}= \begin{cases}a_{i j} & (i=j) \\ 1 / 2 a_{i j} & (i \neq j)\end{cases}
$$

Then for all $\lambda$ we have

$$
\left|A^{\prime}\right|=\left|A^{\prime \prime}\right|
$$

Theorem 3.3
Let $A$ denote a matrix with property (A), and let $\sigma$ denote a consistent ordering. If $\boldsymbol{\omega} \neq 0$, and if $\lambda$ is a non-zero eigenvalue of $I_{\sigma}, \omega$ and if $\mu$ satisfies,

$$
\begin{equation*}
(\lambda+\omega-1)^{2}=\omega^{2} \mu^{2} \lambda \tag{3-7}
\end{equation*}
$$

Then $\mu_{\text {is an }}$ eigenvalue of $B$. On the other hand if $\mu$ is an eigenvalue of $B$, and if $\lambda$ satisfies ( $3-7$ ), then $\lambda$ is an eigenvalue of $I_{\sigma, \omega}$, where $B$ is given by $(3-4)$, and $I_{\sigma, w}$ is defined by $(3-5)$.
To prove this theorem we shall need the following Lemma, and corrollaries.

Lemma:
If $\mu$ is a $k$-fold non-zero eigenvalue of $B$, then $(-\mu)$ is a $k-f o l d$ eigenvalue of $B$.

Cor ollary 3.1
If $\mu$ is an eigenvalue of $B$, then $\mu^{2}$ is an eigenvalue of $I_{\sigma, I}$ (Gauss-Siedel), if $\lambda$ is a non zero eigenvalue of $I_{\sigma, I}$ and if $\mu^{2}=\lambda$, then $\mu_{\text {is }}$ an eigenvalue of $B$.

Cor ollary 3.2
If $A$ is symmetric, then the method of simultanous displacement converges if and only if $A$ is positive definite.

Cor ollary 3.3
If $A$ is symmetric, then there exists $\omega$ such that $I_{\sigma, \omega}$ converges if and only if $A$ is positive definite. Theorem 3.4

Let $\mu$ and $\bar{\lambda}(\omega)$ denote respectively the spectral norms of $B$ and $I_{\sigma, \omega}$. If $\omega_{\text {opt }}$ which satisfies

$$
\omega_{o p t}^{2} \mu^{2}-4\left(\omega_{o p t^{-1}}\right)=0
$$


where $\boldsymbol{\omega}_{\text {opt }}$, the optimum relaxation factor, then the rate of convergence of $I_{\sigma, \omega_{\text {opt }}}$ is given by,

$$
R\left(I_{\sigma, \omega_{o p t}}\right)=-2 \log \frac{\mu}{1+\left(1-\mu^{2}\right)^{1 / 2}}
$$

and for all real $\omega$ such that $\omega \neq \omega_{\text {opt }}$,

$$
R\left(I_{\sigma, \omega}\right)<R\left(I_{\sigma}, \omega_{\text {opt }}\right)
$$

§4. GARABEDIAN METHOD FOR THE ESTIMATION OF THE
RELAXATION FACTOR FOR SMALL MESH SIZE

Consider the Laplace difference equation for an unknown function $u$ of two independent variables in a region $D$ covered by a mesh with $h$ units spaced apart. We use the subscripts $p, q$ to the location of the node points, and superscript $n$ to indicate steps in the
relaxation process, so that the method of successive overrelaxation can be described by the equation. $4\left(u_{q, r}^{n+1}-u_{q, r}^{n}\right)=\omega\left(u_{q-1}^{n+1}, r+u_{q, r-1}^{n+1}+u_{q+1, r}^{n}+u_{q, r+1}^{n}-4 u_{q, r}^{n}\right)(4-1)$ where $\omega$ is the relaxation factor, we express $\omega$ in the form,

$$
\begin{equation*}
\omega=\frac{2}{1+C h} \tag{4-2}
\end{equation*}
$$

where C is any positive value, and constant, if we rearrange (4-1), we get,
$\frac{u_{q-1, r}^{n}+u_{q, r-1}^{n}+u_{q+1, r}^{n}+u_{q, r+1}^{n}-4 u_{q, r}^{n}}{h^{2}}=$
$\frac{u_{q, r}^{n+1}-u_{q, r}^{n}-u_{q-1}^{n+1}+u_{q-1, r}^{n}}{h^{2}}+\frac{u_{q, r}^{n+1}-u_{q, r}^{n}-u_{q, r-1}^{n+1}+u_{q, r-1}^{n}}{h^{2}}+$
$+2 c \frac{u_{q, r}^{n+1}-u_{q_{2}}^{n} r}{h}$
by refering the index $n$ as time variable, and that it indicate the location of new net points spaced at time intervals equall to the original mesh size $h$, it is known that $(4-3)$ is the difference analogue of the hyperbolic partial differential equation.

$$
u_{x x}+u_{y y}=u_{x t}+u_{y t}+2 c u_{t}
$$

where $u_{x x}$, $u_{y y}$ denotes differention with respect to x and y respectively.

Thus for small values of $h$ the convergence of the iterative scheme (4-1) can be investigated by an analysis of the decay of time-dependent terms in the solution of ( $4-4$ ).

The substitution $s=t+x / 2+y / 2$, makes (4-4) in a canonical form,

$$
\begin{equation*}
u_{x x}+u_{y y}-\frac{1}{2} u_{s s}-2 c u_{s}=0 \tag{4-5}
\end{equation*}
$$

For a fixed set of boundary conditions, the method of
separation of variables gives the representation,

$$
u=u_{0}(x, y)+\sum_{m=1}^{\infty}\left[a_{m} e^{-\left(p_{m} s\right)}+b_{m} e^{\left(-q_{m} s\right)}\right]_{U_{m}}(x, y) \quad(4-6)
$$

for the solution of (4-5), where $U_{0}$ is the steady-state solution, where $a_{m}, b_{m}$ are Fourier coefficients, where

$$
\begin{equation*}
p_{m}=2 C-\left(4 C^{2}-2 k_{m}^{2}\right)^{I / 2}, \quad q_{m}=2 C+\left(4 C-2 k_{m}^{2}\right)^{I / 2} \tag{4-7}
\end{equation*}
$$

where $U_{m}$ and $k_{m}^{2}$ are the eigenfunctions and eigenvalues of the equation,

$$
\begin{equation*}
\frac{2}{\nabla} U_{m}+k_{m}^{2} U_{m}=0 \tag{4-8}
\end{equation*}
$$

with homogenous boundary conditions,

$$
\begin{equation*}
\mathrm{p}=\operatorname{Re}\left[\mathrm{p}_{I}\right]=\operatorname{Re}\left[2 \mathrm{C}-\left(4 \mathrm{C}^{2}-2 k^{2}\right)^{1 / 2}\right] \tag{4-9}
\end{equation*}
$$

corresponding to the lowest eigenvalue $k_{I}^{2}$, governs the rate of convergence of the terms on the right in (4-6) with in creasing time t.

By (4-9) the choice of the positive constant $C$ which maximizes $p$ and hence yields the most rapied convergence is clearly $C=k_{i} / 2$, and if $A$ denotes the area of the region $D$, it can be shown that,

$$
\begin{equation*}
k_{1} A^{1 / 2} \geqq \pi^{1 / 2} \tag{4-10}
\end{equation*}
$$

where $k=2.405$ denotes the first root of the Beesel function of the first kind of order zero. Thus the good approximate formula for the relaxation $\omega$ is,

- 22 -
$\omega=\frac{2}{I+(\pi / 2 \mathrm{~A})^{I / 2} k h}=\frac{2}{1+3.014 \mathrm{~h} / \mathrm{A}^{I / 2}}$
This approach is given in the case of five-point Laplace difference equation, an approach to nine-point Laplace difference equation can also be given.


## CHAPTER II

## A NUMERICAL METHOD FOR SOLVING THE EQUATIONS

OF MOTION OF AN INCOMPRESSIBLE VISCOUS FLUID.

## INTRODUCTION

The equations of motion of an incompressible fluid are

$$
\begin{aligned}
& \frac{\partial u_{i}}{\partial t}+u_{j} \frac{\partial u_{i}}{\partial x_{j}}=-\frac{1}{\rho_{0}} \frac{\partial p}{\partial x_{i}}+\nu \nabla^{2} u_{i}+E_{i}, \\
& \frac{\partial u_{j}}{\partial x_{j}}=0 \quad, \quad \nabla^{2} \equiv \sum_{j} \frac{\partial^{2}}{\partial x_{j}^{2}}
\end{aligned}
$$

where $u_{i}$ are the velocity components, $p$ is the pressure, $\rho_{0}$ is the density, $E_{i}$ are the components of the external forces per unit mass, $\nu$ is the coefficient of the kinematic viscosity, $t$ is the time, $i, j=1,2,3 x_{i, j}$ denotes the space coordinates, the summation convention is used in the equations.

We begin by using the method of dimensionalless analysis, writing

$$
u_{i}^{\prime}=\frac{u_{i}}{U}, \quad x_{i}^{\prime}=\frac{x_{i}}{X}, \quad p^{\prime}=\left(\frac{X}{\rho_{O V} U}\right) p
$$

$$
E_{i}^{\prime}=\left(\frac{\nu U}{X^{2}}\right) E_{i}, \quad t^{\prime}=\left(\frac{\nu}{X^{2}}\right) t
$$

where $U$ is a reference velocity, and $X$ is a reference lenght, the equations become

$$
\begin{align*}
\frac{\partial u_{i}}{\partial t}+R u_{j} \frac{\partial u_{i}}{\partial x_{j}} & =-\frac{\partial p}{\partial x_{i}}+\nabla^{2} u_{i}+E_{i}  \tag{I}\\
\frac{\partial u_{j}}{\partial x_{j}} & =0 \tag{2}
\end{align*}
$$

where $R=\frac{U X}{\nu}$ is the Ryenolds number. We now try to introduce a finite difference method for solving these equations in a bounded region $D$, in either two or three dimentional space. The basic feature of this method Iies in the use of equations (1) and (2) rather than higher-order derived equations.

This makes it possible to solve the equations and to satisfy the imposed boundary conditions. We achive adequate computational efficiency, even in problems of three dimensions and space variables.

The princibles of the used method:
Equation (I) can be written in the form

$$
\begin{equation*}
\frac{\partial u_{i}}{\partial t}+\frac{\partial p}{\partial x_{i}}=\mathcal{F}_{i} u \tag{1}
\end{equation*}
$$

where $\mathcal{F}_{i} u$ depends on $u_{i}$ and $E_{i}$, but not on $p$,
equation (2) can be differentiated to give

$$
\begin{equation*}
\frac{\partial}{\partial x_{i}}\left(\frac{\partial u_{i}}{\partial t}\right)=0 \tag{2}
\end{equation*}
$$

The present method can be summerized as follows;
the time $t$ is discretized, and at every time step $\mathcal{F}_{i} u$ is evaluated, then it is decomposed into the sum of a vector with zero divergence and a vector with zero curl. The component with zero divergence is $\frac{\partial u_{i}}{\partial t}$ which can be used to obtain $u_{i}$ at the next time level, and the component with zero curl is $\frac{\partial p}{\partial x_{i}}$.
This decomposition exists and is uniqually determined whenever the initial value problem for the Navier-Stokes equations is well posed.

The existence and uniquness proofs for the solution
of these equations can be seen in[I].
Let $u_{i}, p$ denote only the solution of (I) but also its discrete approximation, and let $D u=0$ be a difference approximation to $\frac{\partial u_{j}}{\partial x_{j}}=0$.
It is assumed that at time $t=n \Delta t \quad$ a velocity and pressure fields $u_{i}^{n}, p^{n}$ are given such that $D u^{n}=0$. The method used is to evaluate $u_{i}^{n+1}, p^{n+1}$ from equation (I) so that $D u^{n+1}=0$. Let $T u_{i}=b u_{i}^{n+1}-B u_{i}$ approximate $\frac{\partial u_{i}}{\partial t}$, where $b$ is a constant and $B u_{i}$ is a suitable Linear Combination of $u_{i}^{n-j}, j \geqq 0 \cdot\left[e \cdot g \frac{\partial u^{n}}{\partial t}=\frac{u^{n+1}-u^{n-1}}{2 \Delta t}-\frac{u^{n+1}-2 u^{n}+2 u^{n-1}-u^{n-2}}{12 \Delta t}+O\left(\Delta t^{4}\right)\right]$.

An auxilairy field $u_{i}^{a u x}$ is first evaluated throught,

$$
\begin{equation*}
b u_{i}^{a u x}-B u_{i}=F_{i} u \tag{3}
\end{equation*}
$$

where $F_{i} u$ approximate $\mathcal{F}_{i} u \cdot u_{i}^{2 u x}$ differs from $u_{i}^{n+1}$ because the pressure term and equation (2) have not been taken into account. $u_{i}^{a u x}$ may be evaluated by an implicit scheme, ie $F_{i} u$ may depend on $u_{i}^{n}$, $u_{i}^{a u x}$ and intermediate fields, say $u_{i}^{*}, u_{i}^{* *}$.
b $u_{i}^{a u x}-B u_{i}$ now approximates $\mathcal{F}_{i}$ u within an error which may depend on $\Delta t$. Let $G_{i} p$ approximates $\frac{\partial p}{\partial x_{i}}$. To obtain $u_{i}^{n+1}$, $p^{n+1}$ it is necessary to perform the decomposition

$$
\begin{align*}
& F_{i} u=b u_{i}^{\operatorname{aux}}-B u_{i}=T u_{i}+G_{i} p^{n+1},  \tag{3}\\
& D(T u)=0
\end{align*}
$$

It is however, assumed that $D u^{n-j}=0, j \geqq 0$. Substituting the value of $T u_{i}$ into equation (3), we obtain

$$
\begin{equation*}
u_{i}^{a u x}=u_{i}^{n+1}+b^{-1} G_{i} p^{n+1} \tag{4}
\end{equation*}
$$

where $D u_{i}^{n+1}=0$, and $u_{i}^{n+1}$ satisfies the prescribed boundary conditions. Since $\mathrm{p}^{\mathrm{n}}$ is usually avaiable and is a good first guess for the values of $p^{n+1}$, the decomposition (4) is probably best done by iteration. For that purpose, we introduce the following iteration scheme:

$$
\begin{align*}
& u_{i}^{n+1, m+1}=u_{i}^{a u x}-b^{-1} G_{i}^{m} p \quad, m \geq 1  \tag{5a}\\
& p^{n+1, m+1}=p^{n+1, m}-\lambda D u^{n+1, m+1} \quad m \geq 1
\end{align*}
$$

where $\lambda$ is a parameter, $u_{i}^{n+1, m+1}$ and $p^{n+1, m+1}$ are successve approximations to $u_{i}^{n+1, m}$ and $p^{n+1, m}$ and $G_{i}^{m} p$ is a fundtron of $p^{n+1, m+1}$ and $p^{n+1, m}$ which converges to $G_{i} p^{n+1}$
as $\left|p^{n+1, m+1}-p^{n+1, m}\right| \rightarrow \underset{m \rightarrow \infty}{\operatorname{zero}}$.
We start by assuming that,

$$
\begin{equation*}
p^{n+1,1}=p^{n} \tag{5c}
\end{equation*}
$$

The iterations (5 a) are to be performed in the interior of $\mathscr{D}$, and the iterations ( 5 b ) in the interior of $\mathscr{D}$ and on its boundary.

It is clear that (5 a) tends to (4) if the iterations converge.
$G_{i}^{m} p$ is used instead of $G_{i} p$ in ( 5 a) to improve the rate of convergence of the iterations. A detailed dis= cussion will be given in a later section. The form of equation (5 b) was suggested by experience with the artifitial compressibility method [2], where for the perpose of finding steady state solutions of equations (1) and (2), p was related to $u_{i}$ by the aquation,

$$
\frac{\partial p}{\partial t}=\operatorname{const} \frac{\partial u_{i}}{\partial x_{j}}
$$

when for some $l$ and small predetermined constant $\mathcal{\varepsilon}$,

$$
\max _{\infty}\left|p^{n+1,1+1}-p^{n+1,1}\right|
$$

we set $u_{i}^{n+1}=u_{i}^{n+1}, t+1, \quad p^{n+1}=p^{n+1}, \mathbf{l}+1$ The iteration (5) ensure that equation (1) including the pressure term is satisfied inside $D$, and equation (2) is satisfied inside $\mathscr{D}$ and on its boundary.

Now we try to find specific schemes for evaluating $u_{i}^{2 u x}$ and specific representations for $D u, G_{i} p, G_{i}^{m} p$, many other schemes and representations can be used [7]. The method which will be presented is efficient, and suitable mainly for problems in which the boundary data are smooth and the domain has a simple shape. Evaluation of $u_{i}^{2 u x}$,

Schemes for evaluating $u_{i}^{a u x}$, defined by (3) will be presented here.

Equation (3) represents one step in time for the solution of the equation

$$
\frac{\partial u_{i}}{\partial t}=\mathcal{F}_{i} u
$$

We can use a combined DuFort-Frankel scheme, in which the time and first space derivatives were approximated by centered differences, and a second derivative such as $\frac{\partial^{2} u}{\partial x_{1}{ }^{2}}$ was replaced by

$$
\frac{1}{\Delta x_{1}{ }^{2}}\left(u_{q+1}^{n}+u_{q-1}^{n}-u_{q}^{n+1}-u_{q}^{n-1}\right), u_{q}^{n} \equiv u(q \Delta x, n \Delta t)
$$

This scheme is sutable only when an asymptotic steady solution is sought. It is inaccurate when real time dependence is studied, unless $\Delta t$ is small.

Our reason for studying this scheme is that, the
DuFort-Frankel scheme is explicit and unconditionally stable; it is a natural scheme to use when the nonlinear terms in (1) are differenced in "Conservation-Law"
form, i.e $\frac{\partial\left(u_{i} u_{j}\right)}{\partial x_{j}}$ rather than $u_{j} \frac{\partial u_{j}}{\partial x_{i}}$, it is found in problems in which the viscosity is not small, it is preferable to use "non-conservative" difference scheme for non-linear terms, and avoid the DuFort-Frankel one. The equation cam be approximated in many ways. But we shall use schemes which are implicit, and accurate to $O(\Delta t)+O\left(\Delta x^{2}\right)$.

Implicit schemes were used because explicit ones requiere, in three space dimenssions that

$$
\Delta t<\frac{1}{6} \Delta x^{2}
$$

which is restrictive condition [2]. Also implicit schemes of accuracy higher than $O(\Delta t)$, require the solution of non-linear equations at every time-step, and make it necessary to evaluate $u_{i}^{a u x}$ and $u_{i}^{n+1}$ simultaneously rather than in succession.

Two schemes will be presented, for both of them,

$$
T u_{i} \equiv\left(u_{i}^{n+1}-u_{i}^{n}\right) / \Delta t ; \quad\left(b \stackrel{-1}{\equiv} \Delta t, B u_{i} \equiv u_{i}^{n} / \Delta t\right)
$$

(A) In two-dimensional problems, we use a BeacemanRachford analogue formula [7]. The implicit form of equation (I) can be written in the form (neglecting the pressure term),

$$
\begin{aligned}
& \exp \left[-\frac{1}{2} \Delta t\left(-R u_{1} D_{1}+D_{1}^{2}\right)\right] \cdot \exp \left[-\frac{1}{2} \Delta t\left(R u_{2} D_{2}+D_{2}^{2}\right)\right] U_{i}(q, r)= \\
& \exp \left[\frac{1}{2} \Delta t\left(-R u_{1} D_{1}+D_{1}^{2}\right)\right] \cdot \exp \left[\frac{1}{2} \Delta t\left(-R u_{2} D_{2}+D_{2}^{2}\right)\right] u_{i}^{n}(q, r)+E_{i}(q, r)
\end{aligned}
$$

$$
\begin{aligned}
& \text { i.e } \\
& \qquad\left[1-\frac{1}{2} \Delta t\left(-R u_{1} D_{1}+D_{1}^{2}\right)\right]\left[I-\frac{1}{2} \Delta t\left(-R u_{2} D_{2}+D_{2}^{2}\right)\right] u_{i}(q, r)= \\
& {\left[I+\frac{1}{2} \Delta t\left(-R u_{1} D_{1}+D_{1}^{2}\right)\right]\left[1+\frac{1}{2} \Delta t\left(-R u_{2} D_{2}+D_{2}^{2}\right)\right] u_{i(q r)}^{n}+} \\
& \quad+E_{i(q, r)}
\end{aligned}
$$

which can be split into two forms, if an intermediate value $\stackrel{*}{u}_{i}^{n+1}=u_{i}^{n+}=\stackrel{*}{u}_{i}$ is introduced, retaining only the second order terms.

$$
\begin{aligned}
& {\left[I-\frac{I}{2} \Delta t\left(-R u_{1} D_{1}+D_{1}^{2}\right)\right]{ }_{u_{i}}^{*}(q, r)=\left[1+\frac{1}{2} \Delta t\left(-R u_{2} D_{2}+D_{2}^{2}\right)\right] u_{i}^{n}(q, r)^{+}} \\
& +\frac{I}{2} E_{i}(q, r) \\
& {\left[I-\frac{I}{2} \Delta t\left(-R u_{2} D_{2}+D_{2}^{2}\right)\right]{ }_{u} q_{i}(q, r)=\left[1+\frac{1}{2} \Delta t\left(-R u_{1} D_{1}+D_{1}^{2}\right)\right] u_{i(q, r)}^{*}} \\
& +\frac{I}{2} E_{i(q, r)}
\end{aligned}
$$

which gives

$$
\begin{aligned}
& \stackrel{u}{*}_{i(q, r)}^{*}=u_{i(q, r)}^{n}-R \frac{\Delta t}{\Delta \Delta x_{1}} u_{1}^{n}(q, r)\left(u_{i}^{*}(q+1, r)-u_{i}^{*}(q-1, r)\right)- \\
& R \frac{\Delta t}{4 \Delta x_{2}} u_{2(q, r)}^{n}\left(u_{i}^{n}(q, r+1)-u_{i}^{n}(q, r-1)\right)+\frac{\Delta t}{2 \Delta x_{1}^{2}}\left(u_{i}^{*}(q+1, r)+\right. \\
& \left.+\stackrel{u}{u}_{i(q-1, r)}^{*}-2 u_{i}^{*}(q, r)\right)+\frac{\Delta t}{2 \Delta x_{2}^{2}}\left(u_{i}^{n}(q, r+1)+u_{i}^{n}(q, r-1)-\right. \\
& \left.-2 u_{i(q, r)}^{n}\right)+\frac{\Delta t}{2} E_{i} \\
& u_{i(q, r)}^{\operatorname{aux}}=u_{i(q, r)}^{*}-R \frac{\Delta t}{4 \Delta x_{1}} u_{I(q, r)}^{*}\left(u_{i}^{*}(q+1, r)-u_{i}^{*}(q-1, r)\right) \\
& -R \frac{\Delta t}{4 \Delta x_{2}} \stackrel{u}{u}_{2(q, r)}^{*}\left(u_{i(q, r+1)}^{a u x}-u_{i(q, r-1)}^{a u x}\right)+\frac{\Delta t}{2 \Delta x_{1}^{2}}\left(u_{i(q+1, r)}^{*}+u_{(q-1, r)}^{*}-2 u_{i(q, r)}^{*}\right) \\
& +\frac{\Delta t}{2 \Delta x_{2}^{2}}\left(u_{i(q, r+1)}^{a_{u x}}+u_{i(a, r-1)}^{a_{u 1 x}}-2 u_{i(q, r)}^{a u x}\right)+\frac{\Delta t}{2} E_{i} . \quad \text { (db) }
\end{aligned}
$$

where $D_{i}=\frac{\partial}{\partial x_{i}}, D_{i}^{2}=\frac{\partial^{2}}{\partial x_{i}^{2}}, i=1,2$

$$
D_{i} u_{m}^{n}=\frac{u_{m+1}^{n}-u_{m-1}^{n}}{2 \Delta x_{i}}, \quad D_{i}^{2} u_{m}^{n}=\frac{u_{m+1}^{n}-2 u_{m}^{n}+u_{m-1}^{n}}{\Delta x^{2}}
$$

where $\vec{u}_{i}^{*}$ are intermediate fields, $u_{i}(q, r) \equiv u_{i}\left(q \Delta x_{i}, r \Delta x_{2}\right)$
(B) In two-dimensional and three-dimensional problems
another scheme suggested by Samaraskii[3]will presented
in the form,

$$
\begin{aligned}
& \ddot{u}_{i}(q, r, s)=u_{i}^{n}(q, r, s)-R \frac{\Delta t}{2 \Delta x_{1}} u_{1}^{n}(q, r, s)\left(\stackrel{u}{u}_{i}(q+1, r, s)-\right. \\
& \left.-u_{i}^{*}(q-1 r, s)\right)+\frac{\Delta t}{\Delta x_{2}^{2}}\left(u_{i}^{*}(q+1, r, s)+\ddot{u}_{i(q-1, r, s)}^{*}-2{ }^{*} u_{i}(q, r, s,\right.
\end{aligned}
$$

$\stackrel{u}{u}_{i(q, r, s)}^{* *}=\stackrel{u}{u}_{i}^{*}(q, r, s)-R \frac{\Delta t}{2 \Delta x_{2}}{\underset{u}{u}}_{2}^{*}(q, r, s)\left({ }^{*}{ }_{u}^{*}(q, r+1, s)-\right.$
$\left.-\stackrel{u}{u}_{i}^{*}(q, r-1, s)\right)+\frac{\Delta t}{\Delta x_{2}{ }^{2}}\left(\dot{u}_{i}^{* *}(q, r+1, s)+\stackrel{u}{u}_{i}^{*}(q, r-1, s)-\right.$
$\left.-2 u_{i}^{* *}(q, r, s)\right)$
(7 b)
$u_{i(q, r, s)}^{u_{x}}=u_{i(q, r, s)}^{* *}-R_{\frac{\Delta t}{2 \Delta x_{3}}}^{u_{i}^{*}} u_{3(q, r, s)}^{*} u_{i(q, r, s+1)}^{q_{i}}-$
$\left.-u_{i(q, r, s-1)}^{u_{1}}\right)+\frac{\Delta t}{\Delta x_{j}^{2}}\left(u_{i}^{\operatorname{qux}}(q, r, s+1)+u_{i(q, r, s-1)^{2}}\right.$
$\left.-2 u_{i(q, r, s)}^{\operatorname{aux}}\right)+\Delta t E_{i(q, r, s)}$
where $u_{i}(q, r, s) \equiv u_{i}\left(q \Delta x_{1}, r \Delta x_{2}, s \Delta x_{3}\right)$

$$
\begin{aligned}
& \mathbb{E}_{i(q, r, s)} \equiv E_{i}\left(q \Delta x_{1}, r \Delta x_{2}, s \Delta x_{3}\right) \\
& \text { and } u_{i}^{*}, \text { U }_{i}^{*} \text { are auxiliary fields. }
\end{aligned}
$$

In symbolic form equations (6) can be written in the form

$$
\begin{aligned}
& \stackrel{u}{i}_{i}^{*}(q, r)=u_{i}^{n}(q, r)-\frac{R}{2} \Delta t u_{i}^{n}(q, r) \frac{\partial u_{i}^{*}(q, r)}{\partial x_{1}}-R \frac{\Delta t}{2} u_{2(q, r)}^{n} .
\end{aligned}
$$

$$
\begin{align*}
& u_{i(q, r)}^{\operatorname{aux}}=\stackrel{u}{u}_{i(q, r)}^{*}-R \frac{\Delta t}{2} \stackrel{*}{u}_{I(q, r)} \frac{\partial u_{i(q, r)}^{*}}{\partial x_{1}}-R \frac{\Delta t}{2} . \\
& \cdot \stackrel{*}{u}_{2(q, r)} \frac{\partial u_{i(q, r)}^{a u x}}{\partial x_{2}}+\frac{\Delta t}{2} \cdot \frac{\partial u_{i(q, r)}^{a u x}}{\partial x_{2}^{2}}+\frac{\Delta t}{2} E_{i} \tag{8~b}
\end{align*}
$$

ie

$$
\begin{align*}
& \left(I+\frac{R}{2} \Delta t u_{I}^{n}(q, r) \frac{\partial}{\partial x_{1}}-\frac{\Delta t}{2} \frac{\partial 2}{\partial x_{1}^{2}}\right) \stackrel{*}{u}_{i}(q, r)= \\
= & \left(I-\frac{R}{2} \Delta t u_{2}^{n}(q, r) \frac{\partial}{\partial x_{2}}+\frac{\Delta t}{2} \frac{\partial^{2}}{\partial x_{2}^{2}}\right) u_{i(q, r)}^{n}+\frac{\Delta t}{2} E_{i} \\
& \left(I-\Delta t Q_{1}\right) u_{i}^{*}=\left(I-\Delta t Q_{2}\right) u_{i}^{n}+\frac{\Delta t}{2} E_{i}
\end{align*}
$$

where

$$
\begin{aligned}
& Q_{1}=\frac{1}{2} \frac{\partial^{2}}{\partial x_{1}^{2}}-\frac{R}{2} u_{1(q, r)}^{n} \frac{\partial}{\partial x_{1}} \\
& Q_{2}=\frac{R}{2} u_{2(q, r)}^{n} \frac{\partial}{\partial x_{2}}+\frac{1}{2} \frac{\partial^{2}}{\partial x_{2}^{2}}
\end{aligned}
$$

also,

$$
\begin{align*}
& \left(I+\frac{R}{2} \Delta t \stackrel{*}{u}_{2(q, r)}^{*} \frac{\partial}{\partial x_{2}}+\frac{\Delta t}{2} \frac{\partial^{2}}{\partial x_{2}^{2}}\right) u_{i}^{a u x}= \\
& \left(I-R \frac{\Delta t}{2} \stackrel{u}{u}_{I(q, r)}^{*} \frac{\partial}{\partial x_{I}}+\frac{\Delta t}{2} \frac{\partial^{2}}{\partial x_{1}^{2}}\right) \stackrel{u}{u}_{i}^{*}+\frac{\Delta t}{2} \mathbb{E}_{i}  \tag{9a}\\
& \left(I-\Delta t Q_{2}^{*}\right) u_{i}^{q u x}=\left(I-\Delta t Q_{1}^{*}\right) \stackrel{u}{u}_{i}^{*}+\frac{\Delta t}{2} \mathbb{E}_{i} \tag{9~b}
\end{align*}
$$

where

$$
\begin{aligned}
& Q_{1}^{*}=\frac{R}{2} u_{1}^{*}(q, r) \frac{\partial}{\partial x_{1}}-\frac{1}{2} \frac{\partial^{2}}{\partial x_{1}^{2}} \\
& Q_{2}^{*}=\frac{1}{2} \frac{\partial^{2}}{\partial x_{2}^{2}}-\frac{R}{2} \dot{u}_{2(q, r)}^{*} \frac{\partial}{\partial x_{2}}
\end{aligned}
$$

where $Q_{1}, Q_{2},{ }^{*}, \stackrel{*}{Q}_{2}$, involves differentiation with respect to variables $x_{1}, x_{2}$, and $I$ the identity operator.

$$
\begin{align*}
& u_{i}^{*}=\left[\left(I-\Delta t Q_{2}\right) u_{i}^{n}+\frac{\Delta t}{2} E_{i}\right]\left(I+\Delta t Q_{1}\right)+O\left(\Delta t^{2}\right) \\
& u_{i}^{*}=\left(I-\Delta t Q_{2}+\Delta t Q_{1}\right) u_{i}^{n}+\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right) \tag{6a}
\end{align*}
$$

$u_{i}^{\text {aux }}=\left[\left(I+\Delta t Q_{1}^{*}\right) \stackrel{u_{i}^{*}}{u_{i}}+\frac{\Delta t}{2} E_{i}\right]\left(I+\Delta t Q_{2}^{*}\right)=\left(I+\Delta t Q_{I}^{*}+\right.$
$\left.+\Delta t Q_{2}^{*}-\Delta t Q_{1}-\Delta t Q_{2}\right) u_{i}^{n}+\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right)$
$u_{i}^{a u x}=\left(I+\Delta t Q_{1}^{*}+\Delta t Q_{2}^{*}\right) \stackrel{u}{u}_{i}^{*}+\frac{\Delta t}{2} E_{i}+o\left(\Delta t^{2}\right)$
$u_{i}^{n+1}=u_{i}^{Q_{i}}-\Delta t G_{i} p^{n+1}=\left(I+\Delta t Q_{1}^{*}+\Delta t Q_{2}^{*}\right)\left\{\left(I-\Delta t Q_{2}+\right.\right.$ $\left.\left.+\Delta t Q_{1}\right) u_{i}^{n}-\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right)\right\}-\Delta t G_{i} p^{n+1}$

$$
\begin{aligned}
& u_{i}^{n+1}=\left(I-\Delta t Q_{2}+\Delta t Q_{1}+\Delta t Q_{1}^{*}+\Delta t \stackrel{*}{Q}_{2}^{*}\right) u_{i}^{n}+\frac{\Delta t}{2} \mathbb{E}_{i}+ \\
& +O\left(\Delta t^{2}\right)-\Delta t G_{i} p^{n+1}
\end{aligned}
$$

we can set at the boundary

$$
\begin{gathered}
{\stackrel{*}{u_{i}}=\left(I-\Delta t Q_{1}^{*}-\Delta t Q_{2}^{*}\right) u_{i}^{n+1}-\Delta t E_{i}+\Delta t \overline{G_{i} p}}^{\text {where } \overline{G_{i} p}=G_{i} p+0(\Delta t)}
\end{gathered}
$$

at the boundary the normal component of $G_{i} p$ is approximated by one-sided differences while it is not necessary in the interior of $\mathscr{D}$.
ie

$$
\begin{aligned}
\stackrel{*}{u_{i}} & =u_{i}^{n+1}-\Delta t \stackrel{\rightharpoonup}{Q}_{1}^{*} u_{i}^{n+1}-\Delta t \stackrel{Q}{Q}_{2}^{*} u_{i}^{n+1}-\Delta t E_{i}+\Delta t \overline{G_{i} p} \\
u_{i}^{u_{i} u x} & =u_{i}^{n+1}+\Delta t G_{i} p
\end{aligned}
$$

But, $\stackrel{u}{u}_{i}=\left(I-\Delta t Q_{2}+\Delta t Q_{I}\right) u_{i}^{n}+\frac{\Delta t}{2} \mathbb{E}_{i}+O\left(\Delta t^{2}\right)$

$$
\begin{aligned}
& u_{i}^{*}=\left[I-\Delta t\left(\frac{R}{2} u_{2(q, r)}^{n} \frac{\partial}{\partial x_{2}}+\frac{1}{2} \frac{\partial^{2}}{\partial x_{2}^{2}}\right)+\right. \\
& \left.+\Delta t\left(\frac{z}{2} \frac{\partial^{2}}{\partial x_{1}{ }^{2}}-\frac{R}{2} u^{n} I(q, r) \frac{\partial}{\partial x_{1}}\right)\right] u_{i}^{n}+\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right) \\
& \stackrel{u}{i}_{i}^{*}=u_{i(q, r)}^{n} \Delta t\left\{\frac{R}{4 \Delta x_{2}} u_{2(q, r)}^{n}\left(u_{i}^{n}(q, r+1)-u_{i}^{n}(q, r-1)\right)+\right. \\
& \left.\left.=\frac{1}{2 \Delta x_{2}{ }^{2}}\left(u_{i(q, r+1}^{n}\right)+u_{i(q, r-1)}^{n}-2 u_{i(q, r)}^{n}\right)\right\}+\Delta t\left\{\frac{1}{2 \Delta x_{1}{ }^{2}} .\right. \\
& \left(u_{i}^{n}(q+1, r)+u_{i}^{n}(q-1, r)-2 u_{i}^{n}(q, r)\right)-\frac{R}{4 \Delta x_{1}} u^{n} I(q, r) \cdot
\end{aligned}
$$

$$
\begin{align*}
& \text { - } \left.\left.\left(u_{i(q+1, r)}^{n}-u_{i(q-1, r)}^{n}\right)\right\}\right]+\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right)  \tag{10}\\
& \text { also, } \\
& u_{i}^{\text {aux }}=\left[\left(I+\Delta t Q_{I}^{*}+\Delta t Q_{2}^{*}-\Delta t Q_{I}-\Delta t Q_{2}\right) u_{i}^{n}+\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right)\right] \\
& =\left[I+\Delta t\left(\frac{R}{2} u_{I}^{*} \frac{\partial}{\partial x_{I}}-\frac{I}{2} \frac{\partial^{2}}{\partial x_{I}^{2}}\right)+\Delta t\left(\frac{I}{2} \frac{\partial^{2}}{\partial x_{2}^{2}}-\frac{R}{2} u_{2}^{*} \frac{\partial}{\partial x_{2}}\right)-\right. \\
& \left.-\Delta t\left(\frac{1}{2} \frac{\partial^{2}}{\partial x_{1}{ }^{2}}-\frac{R}{2} u_{1}^{n} \frac{\partial}{\partial x_{1}}\right)-\Delta t\left(\frac{R}{2} u_{2}^{n} \frac{\partial}{\partial x_{2}}+\frac{1}{2} \frac{\partial^{2}}{\partial x_{2}{ }^{2}}\right)\right] u_{i}^{n}+ \\
& +\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right) \\
& \text { ide } \\
& u_{i}^{a u x}=\left[u_{i}^{n}+\Delta t\left\{\frac{R}{4 \Delta x_{I}} \stackrel{u}{u}_{I}^{*}\left(u_{i(q+1, r)}^{n}-u_{i(q-1, r)}^{n}\right)-\frac{I}{4 \Delta x_{1}{ }^{2}} .\right.\right. \\
& \text { - }\left(\dot{u}_{i}^{n}(q+1, r)+u_{i}^{n}(q-1, r)-2 u_{i}^{n}(q, r)\right\}+\Delta t\left\{\frac{1}{4 \Delta x_{1}{ }^{2}}\right. \text {. } \\
& \text { - }\left(u_{i(q, r+1)}^{n}+u_{i(q, r-1)}^{n}-2 u_{i(q, r)}^{n}\right)-\frac{R}{4 \Delta x_{2}} u_{2}^{*}\left(u_{i}^{n}(q, r+1)-\right. \\
& \left.\left.-u_{i(q, r-1}^{n}\right)\right\}-\Delta t \quad \frac{1}{4 \Delta x_{1}^{2}}\left(u_{i(q+1, r)}^{n}+u_{i(q-1, r)}^{n}-2 u_{i(q, r)}^{n}\right) \\
& \left.-\frac{R}{4 \Delta x_{1}} u_{1(q, r)}^{n}\left(u_{i(q+1, r)}^{n}-u_{i(q-1, r)}^{n}\right)\right\}-\Delta t\left\{\frac{R}{4 \Delta x_{2}} u_{2(q, r)}^{n} \cdots\right. \\
& \text { - } \left.\left(u_{i(q, r+1}^{n}\right)-u_{i(q, r-1)}^{n}\right)+\frac{1}{4 \Delta x_{2}^{2}}\left(u_{i(q, r+1)}^{n}+u_{i(q, r-1)}^{n}-\right. \\
& \left.\left.-2 u_{i(q, r)}^{n}\right)\right\}+\frac{\Delta t}{2} E_{i}+O\left(\Delta t^{2}\right)  \tag{11}\\
& \text { also, }
\end{align*}
$$

A similar expressions for scheme (B) can be written in symbolic form, as follows:

$$
\begin{align*}
& \left.\left(I-\Delta t Q_{1}\right)\right)_{i}^{*}=u_{i}^{n} \\
& \left(I-\Delta t Q_{2}\right) u_{i}^{* *}=u_{i}^{*}  \tag{12}\\
& \left(I-t Q_{3}\right) u_{I}^{a u x}=\stackrel{u}{u}_{i}^{* *}+\Delta t E_{i}
\end{align*}
$$

Where $I$ is the identity operator, and $Q_{p}$ represents differentiations with respeet to $\mathrm{x}_{\mathrm{l}}$ only. It is clear that scheme (6) is accurate to $O\left(\Delta t^{2}\right)+O\left(\Delta x^{2}\right)$ in both cases when $R=0$ and $R \neq 0$.

If both schemes are to be used in a problem in wich the velocities are known at the boundary, values of ${\underset{u}{i}}_{*}^{*}, u_{i}^{* *}$, $u_{i}^{a u x}$ at the boundary have to be given in advance so that the several implicit operators can be inverted. In the case of scheme (12), we have,
$u_{i}^{n+1}=\left(I+\Delta t Q_{1}+\Delta t Q_{2}+\Delta t Q_{3}\right) u_{i}^{n}+\Delta t E_{i}-\Delta t G_{i} p^{n}+O\left(\Delta t^{2}\right)$
$\stackrel{u_{i}}{*}=\left(I+\Delta t Q_{I}\right) u_{i}^{n}+O\left(\Delta t^{2}\right)$
$\stackrel{* *}{u_{i}}=\left(I+\Delta t Q_{I}+\Delta t Q_{2}\right) u_{i}^{n}+O\left(\Delta t^{2}\right)$
$u_{i}^{\text {aux }}=\left(I+\Delta t Q_{I}+\Delta t Q_{2}+\Delta t Q_{3}\right) u_{i}^{n}+\Delta t E_{i}+O\left(\Delta t^{2}\right)$
The scheme will be accurate to $O(\Delta t)$ at the boundary if, $u_{i}^{*}=u_{i}^{n+1}-\Delta t Q_{2} u_{i}^{n+1}-\Delta t Q_{3} u_{i}^{n+1}-\Delta t E_{i}+\Delta t G_{i} p^{n}$

$$
\begin{aligned}
& u_{i}^{* *}=u_{i}^{n+1}-\Delta t Q_{3} u_{i}^{n+1}+\Delta t \overline{G_{i} p} \\
& u_{i}^{\text {aux }}=u_{i}^{n+1}+\Delta t G_{i} p^{n}
\end{aligned}
$$

where $\qquad$

$$
G_{i} p^{n}=G_{i} p^{n}+o(\Delta t)
$$

It is clear that more accurate expressions for the auxiliary fields at the boundaries can be used but it needs great programming effort on the computer. In case of negligible viscosity, i.e. $\nu=0$, another schemes will used, i.e explicit schemes which accurate to $O\left(\Delta t^{2}\right)+O\left(\Delta x^{2}\right)$, and stable when $\Delta t=O(\Delta x)$. A scheme we suggested for this case can be given in the explicit form,

$$
\begin{aligned}
& u_{i(q, r)}^{n+1}=\exp \left[-\Delta t\left(u_{1} \frac{\partial u_{i}}{\partial x_{1}}+u_{2} \frac{\partial u_{i}}{\partial x_{2}}\right)\right] u_{i(q, r)}^{n}+\mathbb{E}_{i(q, r)} \\
= & {\left[1-\Delta t\left(u_{1} \frac{\partial u_{i}}{\partial x_{1}}+u_{2} \frac{\partial u_{i}}{\partial x_{2}}+\cdots\right] u_{i(q, r)}^{n}+\mathbb{E}_{i f q, r)}=\right.} \\
= & u_{i(q, r)}^{n}-\Delta t\left\{u_{1}^{n}(q, r) \frac{u_{i}^{n}(q+1, r)-u_{i}^{n}(q-1, r)}{2 \Delta x_{1}}\right. \\
= & \left.u_{2(q, r)}^{n} \frac{u_{i(q, r+1)}^{n}-u_{i-q, r-1)}^{n}}{2 \Delta x_{2}}\right\}+\mathbb{E}_{i(q, r)}
\end{aligned}
$$

retained only the first order terms, which can be solved to give $u_{i}^{n+1}$. Such problem can be discussed later. The rest of this work will show how we can derive $D$, $G_{i}^{m}$, and the choice of $\lambda$, used in ( 5 a) and ( 5 b ), so we need some facts about the DuFort - Frankel scheme for heat equation, and its relation to the relaxation
method for solving the Laplace equation [7]. Consider the equation,

$$
\begin{equation*}
-\nabla^{2} u= \pm, \quad \nabla^{2}=\frac{\partial^{2}}{\partial x_{1}{ }^{2}}+\frac{\partial^{2}}{\partial x_{2}{ }^{2}} \tag{13}
\end{equation*}
$$

in some domain $D$, rectangle for example. $u$ is assumed known on the boundary of $\mathfrak{D}$, it can be approximate to the equation,

$$
\begin{equation*}
-\operatorname{Iu}=f \tag{14}
\end{equation*}
$$

it is clear that $I$ is the five point approximation to the Laplacian, $u$ and $f$ are now m-component vectors.

$m$ is the number of internal nodes of the resulting difference equation. For simplicity we assume that the mesh wedthes in the $\mathrm{x}_{1}, \mathrm{x}_{2}$ directions are equal, i.e.

$$
\Delta x_{1}=\Delta x_{2}=\Delta x \text {, so the operator I is represented }
$$ by an m $x$ m matrix $A$. The matrix $A$ can be written in the form,

$$
A=A^{\prime}-E-E^{\prime}
$$

where Á is diagonal, and $\mathbb{E}, E^{\prime}$ respectively upper and Iower triangular matrices. The convergent relaxation scheme for solving (14) is given by,

$$
\left(A^{\prime}-\omega \mathbb{E}\right) u^{n+1}=\left\{(1-\omega) A^{\prime}+\omega_{E}^{\prime}\right\} u^{n}+\omega_{f}
$$

and hence $u^{n+1}=\left(A^{\prime}-\omega E\right)^{-1}\left[(1-\omega) A^{\prime}+\omega E^{\prime}\right] u^{n}+\boldsymbol{\omega}\left(A^{\prime}-\omega E\right)^{-1} I^{\prime}$ where $\omega$ is the relaxation factor, $0<\omega<2$, and $u^{n}$ are the successive iterates. [5], [7].
It is known that there is optimal relaxation factor $\boldsymbol{\omega}_{\text {opt }}$ depends on the fact that Aasatisfies "Young, s condition (A)" [9]. i.e there exists a permutation matrix $P$ such that,

$$
\begin{equation*}
P^{-1} A P=\Lambda-\mathbb{N} \tag{15}
\end{equation*}
$$

where $\Lambda$ is diagonal, and $N$ has the normal form,

$$
\left[\begin{array}{ll}
0 & G \\
G & 0
\end{array}\right]
$$

The zero submatrices here are squere, under this condition $\omega_{o p t}$ can be determined. The matrix $A$ depends on the order in which the components of $u^{n+1}$ are computed from $u^{n}$. The changing of that order is equivalent to transforming $A$ into $P^{-1} A P$, where $P$ is a permutation matrix.

The solution (17) is consider to be the steady solution of,

$$
\begin{equation*}
\frac{\partial u}{\partial \tau}=\nabla^{2} u+£ \tag{16}
\end{equation*}
$$

The latter equation can be approximated by the DuFortFrankel scheme,
$u_{q, r}^{n+1}-u_{q, r}^{n-1}=\frac{2 \Delta \tau}{\Delta x^{2}}\left(u_{q+1, r}^{n}+u_{q-1, r}^{n}+u_{q, r+1}^{n}+u_{q, r-1}^{n}-2 u_{q, r}^{n+1}-\right.$
$\left.-2 u_{q, r}^{n-1}\right)+2 \Delta \tau f$
where

$$
u_{q, r}^{n} \equiv u\left(q \Delta x_{1}, r \Delta x_{2}, n \Delta \tau\right)
$$

which approximates (13), where $\Delta \tau=O(\Delta x)$, we obtain, $\left(1+4 \frac{\Delta \tau}{\Delta x^{2}}\right) u_{q, r}^{n+1}-\left(1-4 \frac{\Delta \tau}{\Delta x^{2}}\right) u_{q, r}^{n-1}=2 \frac{\Delta \tau}{\Delta x^{2}}\left(u_{q+1, r}^{n}+\right.$ $\left.+u_{q-1, r}^{n}+u_{q, r+1}^{n}+u_{q, r-1}^{n}\right)+2 \boldsymbol{\Delta} \tau f$
where

$$
u_{q, r}^{n} \equiv \frac{1}{2}\left(u_{q, r}^{n+1}+u_{q, r}^{n-1}\right)
$$

Clearly $u_{q, r}^{n}$ does not appear in (17) so the calculation splits into two independent calculations on inteqwined meshes, one of which can be omitted then we can write,

$$
U^{n+1}=\binom{u^{2 n}}{u^{2 n+1}} \quad, \quad\left(U^{n+1} \text { has } m \text { components }\right)
$$

when we write,

$$
\begin{equation*}
\omega=\frac{8 \Delta \tau / \Delta x^{2}}{1+4 \Delta \tau / \Delta x^{2}} \tag{18}
\end{equation*}
$$

We see that the iteration (17) reduces to an iteration of the form (15), where the new components of $U^{n+1}$ are calculated in an order such that $A$ has the normal form $(15)^{\prime}$. [This is clear since, the difference equation of equation (14) can be written in the form,

$$
\frac{1}{\Delta x^{2}}\left(u_{q, r+1}^{n}+u_{q+1, r}^{n}+u_{q, r-1}^{n}+u_{q-1, r}^{n}-4 u_{q, r}^{n}\right)+f=0
$$ we can write,

$$
u_{q, r}^{n}=\frac{1}{2}\left(u_{q, r}^{n+1}+u_{q, r}^{n-1}\right) \text {, hence the equation takes }
$$ the form,

$$
\begin{aligned}
& \quad \frac{1}{\Delta x^{2}}\left(u_{q+1, r}^{n}+u_{q-1, r}^{n}+u_{q, r+1}^{n}+u_{q, r-1}^{n}-\right. \\
& \left.\left.-2 u_{q, r}^{n+1}-2 u_{q, r}^{n-1}\right)+f=0\right]
\end{aligned}
$$

Then it, ${ }^{\text {is }}$ clear that the DuFort - Franke scheme appears to be a particular ordering of the over-relaxation method whose existence is equivalent to Young's condition (A).

The best value of $\Delta \tau$, i.e $\Delta \tau$ opt can be determined from $\omega_{\text {opt }}$ and equation (18), clarly $\Delta \tau_{\text {opt }}=O(\Delta x)$, then for $\Delta \tau=\Delta \tau_{\text {opt }}$ the DuFort - Frankel scheme approximate also the equation,

$$
\frac{\partial u}{\partial \tau}=\nabla^{2} u-2\left(\frac{\Delta \tau}{\Delta x}\right)^{2} \frac{\partial^{2} u}{\partial \tau^{2}}+f
$$

see [4].
These remarks can be generalized to problems of more than two space variables. Also it will be noted that, we can approximate equation (16) by explicit method

$$
\begin{equation*}
u_{q, r}^{n+1}-u_{q, r}^{n}=\frac{\Delta \tau}{\Delta x^{2}}\left(u_{q+1, r}^{n}+u_{q-1, r}^{n}+\frac{u_{q, r+1}^{n}}{n}+u_{q, r-1}^{n}-4 u_{q, r}^{n}\right)+\Delta \tau f \tag{19}
\end{equation*}
$$

and used as an iteration procedure for solving (14), but the iteration converges only when $\Delta \tau / \Delta x^{2}<I / 4$, and
converges very slow [4].

## The representation of $D_{2},{ }_{i}, G_{i}^{m}$ and the iteration

 procedure for determining $u_{i}^{n+1}, p^{n+1}$.For simplicity we shall assume that the domain $\mathscr{D}$ is two-dimensional and rectangular, and the velocities are know at the boundary. Extension three - dimensional problems is possible, also domains of other shapes can be treated by the help of interpolation procedures. Firstly we define $D$. Let $\beta$ denote the boundary of $\operatorname{Dand} C$ the set of mesh nodes with a neighbor in $\beta$. Ind- $\beta$ we approximate the equation of continuaty by the centered differences, ie.

$$
D u=\frac{1}{2 \Delta x_{1}}\left(u_{1(q+1, r)}-u_{1(q-1, r)}\right)+\frac{1}{2 \Delta x_{2}}\left(u_{2(q, r+1)}-u_{2(q, r-1)}\right)=0(20)
$$

At the points of $\beta$ we use second-order one-sided differences, so that $D u$ is accurate to $O\left(\Delta x^{2}\right)$ everywhere. On the boundary line $x_{2}=0$, we have,

 $i=I, 2$ and $p^{n+1, m}$ we shall evaluate simultanausly
$u_{1}^{n+1, m+1}(q+1, r), u_{2(q, r+1}^{n+1, m+1}$ and $p_{q, r}^{n+1, m+1}$ using the formla

$$
p_{q, r}^{n+1, m+1}=p_{q, r}^{n+1, m}-\lambda D u^{n+1, m+1}
$$

with similar expressions at the other boundaries. Clearly equation (20) states that the total flow of the fluid into a rectangle of sides $2 \Delta x_{1}, 2 \Delta x_{2}$ is zero, while equation (2l) does not have this elementry interpolation.


Also we define $G_{i} p$ at every point of $\mathcal{D}_{-} \beta$ by,

$$
\begin{aligned}
& G_{1} p=\frac{1}{2 \Delta x_{1}}\left(p_{q+1, r}-p_{q-1, r}\right) \\
& G_{2} p=\frac{1}{2 \Delta x_{2}}\left(p_{q, r+1}-p_{q, r-1}\right)
\end{aligned}
$$

where

$$
p_{q, r} \equiv p\left(q \Delta x_{1}, r \Delta x_{2}\right)
$$

It is clear that $\frac{\partial P}{\partial x}$ is approximated by centered differenes. One can use other forms for $G_{i} p$ and Du. Our purpose now is to perform the decomposition (4). $u_{i}^{n+I}$ is given on the boundary $\beta$, $u_{i}^{a u x}$ is given in $D_{-} \beta$, also $\mathrm{p}^{\mathrm{n}+1}$ is to be found in $D$ (including the boundary) and $u^{n+1}$ in $D-\beta$, so that in $D-\beta$

$$
u_{i}^{q u x}=u_{i}^{n+1}+\Delta t G_{i} p
$$

and in $D(i n c l u d i n g$ the boundary)

$$
D u^{n+1}=0
$$

This must be done using the iterations (5), until now the form of $G{ }_{i}^{m} p$ is not specified. At a point ( $q, r$ ) in D- $\beta-C$, i.e. far from the boundary, we cun substitute equation (5 a) into equation (5 b) and obtain,

$$
\begin{equation*}
p^{n+1, m+1}-p^{n+1, m}=-\lambda D u^{2 u x}+\Delta t D G^{m} p \tag{22}
\end{equation*}
$$

An analogae to this method was used by Harlow and Welch [9], as follows:
Let $D u=0$ approximate $\frac{\partial u_{j}}{\partial x_{j}}=0$, and $G_{i} p$ approximate $\frac{\partial P}{\partial x_{i}}$. It is assumed that at time $t=n \Delta t$ velocity fields $u_{i}^{n}$ are given, satisfying $D u^{n}=0$, then equation (2) can be approximated by

$$
\begin{equation*}
u_{i}^{n+1}=u_{i}^{n}+\Delta t I u_{i}^{n}-\Delta t Q_{i} u^{n}-\Delta t Q_{i} p^{n}+\Delta t E_{i} \tag{23}
\end{equation*}
$$

wher $\operatorname{Iu}$ approximates $\nabla^{2} u$, and $Q_{i} u$ approximates $\frac{\partial u_{i} u_{j}}{\partial x_{j}}$. Performing the operator $D$ on the previous equation, assuming

$$
\begin{array}{r}
D u^{n+1}=0 \text {, we have } \\
I p^{n}=-\frac{D u^{n}}{\Delta t}+D I u^{n}-D Q u^{n}+D E_{i} \tag{23}
\end{array}
$$

wher $I \prime p \equiv D G p$ approximates $\nabla^{2} p$. This equation is a difference analogue of the equation

$$
\begin{equation*}
\nabla^{2} p=-\frac{\partial}{\partial x_{i}} \frac{\partial}{\partial x_{j}} u_{i} u_{j}+\frac{\partial E_{j}}{\partial x_{j}} \tag{24}
\end{equation*}
$$

which can be obtained from equation (I) by taking its divergence. In view of the definitions of $D, G_{i}^{m}$, and $u_{i}^{a u x}$, equation (22) is an iteration procedure for solving an analogue of equation (23). In this sense the method used is related to Harlow and Welch like a predictor-corrector method, wheareas Harlow and Welch first determine $p^{n}$ so that $D u^{n+1}=0$, a guess will made at the values at $u_{i}^{n+1}, p^{n+1}$, and then correct them until the condition $D u^{n+I}=0$ is satisfied. It is clear that at the points of $\beta$ or $C \quad$ it is posible to substitute ( 5 a) into ( 5 b) because at the boundary $u_{i}^{n+1}$ is prescribed, $u^{n+1, m+1}=u^{n+1}$ for allm, ( 5 a) does not hold and therfore (22) is not true. Near the boundary the iterations (5) provide boundary data for (23) and ensure that the constraint of incompressibility is satisfied. We proceed as fallows:
We chose $G_{i}^{m} p$ and $\lambda$ such that (22) is rapidly converging iteration for solving (23); $G_{i}^{m} p$ at the boundary are chosen so that the iteration (5) converges everywhere.

Let ( $q, r$ ) again be a node inD $-\beta-C \cdot u_{i}^{n+1, m+1}$
and $p^{n+1, m}$ are assumed known. We shall evaluate simultanously $p_{q, r}^{n+1, m+1}$ and the velocity components involved in the equation $D u^{n+1}=0$ at $(q, r)$, i.e $\left.u_{1}^{n+1, m+1} q_{1}^{+1}, r\right), u_{2(q, r+1)}^{n+1, m+1}$

These velocity components depend on the value of $p$ at ( $q, r$ ) and on the values of $p$ at the other points: The value of $p$ at ( $q, r$ ) can taken to be,

$$
\frac{1}{2}\left(p_{q, r}^{n+1}, m+1+p_{q, r}^{n+1, m}\right)
$$

while at the other points we use $p^{n+1, m}$. This leads to the following formula,

$$
\begin{equation*}
\mathrm{p}_{\mathrm{q}, \mathrm{r}}^{\mathrm{n}+1, m+1}=\mathrm{p}_{\mathrm{q}, \mathrm{r}}^{\mathrm{n}+1, m}-\lambda D u^{\mathrm{n}+1, m+1} \tag{25a}
\end{equation*}
$$

where Du is given by (20).

$$
\begin{align*}
& u_{1(q+1, r)}^{n+1, m+1}=u_{(q+1, r)}^{\text {aux }}-\frac{\Delta t}{2 \Delta x_{1}}\left(P_{q+2, r}^{n+1, m}-\frac{1}{2}\left(P_{q, r}^{n+1, m+1}+P_{q, r}^{n+1, m}\right)\right)  \tag{25~b}\\
& u_{1(q-1, r)}^{n+1, m+1}={\underset{i}{(q-1, r)}}_{\text {aux }}^{u}-\frac{\Delta t}{2 \Delta x_{1}}\left(\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)-p_{q-2, r}^{n+1, m}\right)  \tag{25c}\\
& {\underset{2}{u}(q, r+1)}_{u^{n+1, m+1}}=u_{2(q, r+1)}^{a u x}-\frac{\Delta t}{2 \Delta x_{2}}\left(p_{q, r+2}^{n+1, m}-\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)\right)  \tag{25d}\\
& u_{2(q, r-1)}^{n+1, m+1}={\underset{2}{(q, r-1)}}_{u_{c c x}}^{u^{c}}-\frac{\Delta t}{2 \Delta x_{2}}\left(\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)-p_{q, r-2}^{n+1, m}\right) \tag{25e}
\end{align*}
$$

It is clear that $G_{i}^{m} p \rightarrow G_{i} p\left[\right.$ since, $u_{I}^{n+1, m+1}(q-1, s) \rightarrow u_{I}^{n+1}(q-1, r)$
as $p^{n+1, m+1} \rightarrow p^{n+1}$. ${ }^{m}$ The first equation gives,

$$
\begin{aligned}
& U_{1(q-1, r)}^{n+1, m+1}=u_{1(q-1, r)}^{a_{1}}-\frac{\Delta t}{2 \Delta x_{1}}\left(\frac{1}{2}\left(p_{q, r}^{n+1, m}+p_{q, r}^{n+1, m}\right)-p_{q-2, r}^{n+1, m}\right. \\
& u_{1(q-1, r)}^{n+1}=u_{1(q-1, r)}^{u_{1}}-\frac{\Delta t}{2 \Delta x_{1}}\left(p_{q, r}^{n+1, m}-p_{q-2, r}^{n+1, m}\right)
\end{aligned}
$$

$\left.u_{1(4-1, r)}^{n+1}=u_{1(q-1, r)}^{a u x}-\Delta t \frac{\partial P_{q_{1}, r}^{n+1}}{\partial x_{1}}=u_{1(q-1, r)}^{a u x}-\Delta t G_{1} p^{n+1}\right]$.
and similar expressions for the other equations. In $C$ and $\beta$ these formulae have to be modified. Consider again the boundary line $\mathrm{x}_{2}=0$, assume the velocities are prescribed at the boundary ie $u_{i}^{n+1}(q, I)$ are given, $i=1$, 2. There are several ways of including that information in the iteration (5). The consistent way would to set.

$$
u_{i(q, I)}^{u_{x}}=u_{i(q, I)}^{n+1}+\Delta t G_{i} p^{n}
$$

and

$$
u_{i(q, I}^{n+1}, m+1, u_{i}^{n+1}(q, I)
$$

for the sake of simplicity, we chose an in inconsistent way of treating the boundary, we set $\left.u_{i(q, I)}^{a u x}=u_{i}^{n+1}, m, I\right)=$ $=u_{i(q, 1}^{n+1}, m+1 \quad u_{i(q, I)}^{n+1}$. This does not affect the values of $u_{i}^{n+1}$, it introduce an additional error of $O(\Delta t)$ into the computed ppessure term. Equations (25) can be solred for $p_{q, r}^{n+1, m+1}$ as follows:

$$
p_{q, r}^{n+1, m+1}=p_{q, r}^{n+1, m}-\lambda\left[\frac{1}{2 \Delta x_{1}}\left(u_{1(q+1, r)}^{n+1, m+1}-u_{i(q-1, r)}^{n+1, m+1}\right)+\frac{1}{2 \Delta x_{2}}\left(\begin{array}{c}
u_{2}^{n+1, m+1} \\
2(q, r+1) \\
-u_{2(q, r-1)}^{n+1, m+1}
\end{array}\right)\right]
$$

$$
=p_{q, r}^{n+1, m}-\lambda\left[\frac{1}{2 \Delta x_{1}}\left\{\begin{array}{l}
\text { aux } \\
u_{(q+1, r)}-\frac{\Delta t}{2 \Delta x_{1}}\left(p_{q+2, r}^{n+1, m}-\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)\right), ~\left({ }_{q}\right)
\end{array}\right)\right.
$$

$$
\left.-\left(u_{1(q-1, r)}^{\text {and }}-\frac{\Delta t}{2 \Delta x_{1}}\left(\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)-p_{q-2, r}^{n+1, m}\right)\right)\right\}+
$$

$$
\begin{aligned}
& \frac{1}{2 \Delta x_{2}}\left\{\begin{array}{l}
u_{2}^{a u x} \\
2(q, r+1)
\end{array}-\frac{\Delta t}{2 \Delta x_{2}}\left(p_{q, r+2}^{n+1, m}-\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)\right)-\right. \\
& \quad\left(\begin{array}{l}
\left.\left.u_{2(q, r-1)}^{a u x}-\frac{\Delta t}{2 \Delta x_{2}}\left(\frac{1}{2}\left(p_{q, r}^{n+1, m+1}+p_{q, r}^{n+1, m}\right)-p_{q, r-2}^{n+1, m}\right)\right\}\right]
\end{array}, ~\right.
\end{aligned}
$$

collecting the similar terms we find that;

$$
\begin{aligned}
& p_{q, r}^{n+1, m+1}\left(1+\alpha_{1}+\alpha_{2}\right)=\left(1-\alpha_{1}-\alpha_{2}\right) p_{q, r}^{n+1, m}-\lambda D u^{a u x}+\alpha_{1}\left(p_{q+2, r}^{n+1, m}+p_{q-2, r}^{n+1, m}\right)+\alpha_{2}\left(p_{q, r+2}^{n+1, m}+p_{q, r-2}^{n+1, m}\right) \\
& \quad \text { ie. } \\
& p_{q, r}^{n+1, m+1}=\left(1+\alpha_{1}+\alpha_{2}\right)^{-1}\left[\left.\left(1-\alpha_{1}-\alpha_{2}\right)\right|_{q, r} ^{n+1, m}-\lambda D u^{a u x}+\alpha_{1}\left(p_{q+2, r}^{n+1, m}+p_{q-2, r}^{n+1, m}\right)+\alpha_{2}\left(p_{q, r+2}^{n+1, m}+p_{q, r-2}^{n+1, m}\right)\right]
\end{aligned}
$$

where

$$
\begin{gathered}
\alpha_{i}=\lambda \Delta t / 4 \Delta x_{i}^{2}, i=1,2 \text {, also, } \\
\left.D u^{\text {aux }}\right|_{q, r}=\frac{1}{2 \Delta x_{1}}\left(u_{1(q+1, r)}^{\text {aux }}-u_{1(q-1, r)}^{\text {aux }}\right)+\frac{1}{2 \Delta x_{2}}\left(u_{2(q, r+1)}^{\text {aux }}-\frac{u_{2}^{a u x}}{(q, r-1)}\right)
\end{gathered}
$$

This can be seen to be a DuFort - Frankel relaxation scheme for the solution of (23). The $\Delta \tau$ of the proceeding equations (17) is replaced by $\lambda \frac{\Delta t}{2}$. It is clear that corresponding to $\Delta \tau$ opt or $\omega_{\text {opt }}$, we find $\lambda_{\text {opt }}$. If $p$ were known on $\beta$ and $c$, convergence of the iterations (26 a) would fallow and $\boldsymbol{\lambda}=\boldsymbol{\lambda}_{\text {opt }}$ would lead to fastest convergence.
In $\beta$ and $C$ formulae (25) are modified by the use of the values of $u_{i}^{n+1}$ at the boundary.

$D u^{n+1, m+1}=\frac{1}{2 \Delta x_{1}}\left(u_{1(q+1,2)}^{n+1, m+1}-u_{1(9-1,2)}^{n+1, m+1}\right)+\frac{1}{2 \Delta x_{2}}\left(\begin{array}{l}n+1, m+1 \\ u_{2}(q, 3)\end{array} u_{2(q, 1)}^{n+1, m+1}\right)$
$u_{i}^{n+1, m+1}$ takes the form,
$u_{1(q+1,2)}^{n+1, m+1}=u_{1(q+1,2)}^{a u x}-\frac{\Delta t}{2 \Delta x_{1}}\left(p_{q+2,2}^{n+1, m}-\frac{1}{2}\left(p_{q, 2}^{n+1, m+1}+p_{q, 2}^{n+1, m}\right)\right)$
$u_{1(q-1,2)}^{n+1, m+1}={\underset{1}{\text { aux }}}_{u_{q-1,2)}}^{\text {and }}-\frac{\Delta t}{2 \Delta x_{1}}\left(\frac{1}{2}\left(p_{q, 2}^{n+1, m+1}+p_{q, 2}^{n+1, m}\right)-p_{q-2,2}^{n+1, m}\right), q>2$
$u_{\substack{1(q-i, 2) \\ n+1, m+1}}^{n+1, m+1}=u_{1(q-1,2)}^{a u x} \quad, q=2$
$u_{2(q, 3)}^{n+1, m+1}={\underset{2}{2}(q, 3)}_{a u x}^{u}-\frac{\Delta t}{2 \Delta x_{2}}\left(p_{q, 4}^{n+1, m}-\frac{1}{2}\left(p_{q, 2}^{n+1, m+1}+p_{q, 2}^{n+1, m}\right)\right)$
$u_{2(9,1)}^{n+1, m+1}=u_{2(9,1)}^{\operatorname{aux}}$
Subestituting equations (28) into (27) we have, using
the equation,

$$
\begin{aligned}
p_{q, 2}^{n+1, m+1} & =p_{q, 2}^{n+1, m}-\frac{\lambda}{2 \Delta x_{1}} u_{1(q+1,2)}^{n+1, m+1}+\frac{\lambda}{2 \Delta x_{1}(q-1,2)} u_{(1, m+1}^{n+\frac{\lambda}{2 \Delta x_{2}^{2}(q, 3)}} u_{2+1, m+1}^{n}+\frac{\lambda}{2 \Delta x_{2}} u_{2(q, 1)}^{n+1, m+1} \\
u_{q, 2}^{n+1, m+1} & =p_{q, 2}^{n+1, m}-\frac{\lambda}{2 \Delta x_{1}}\left\{u_{(q+1,2)}^{\text {un }}-\frac{\Delta t}{2 \Delta x_{1}}\left(p_{q+2,2}^{n+1, m}-\frac{1}{2}\left(p_{q, 2}^{n+1, m+1}+p_{q, 2}^{n+1, m}\right)\right)\right\} \\
& +\frac{\lambda}{2 \Delta x_{1}}\left\{\begin{array}{l}
\left.u_{1(q-1,2)}^{\text {aux }}-\frac{\Delta t}{2 \Delta x_{1}}\left(\frac{1}{2}\left(p_{q, 2}^{n+1, m+1}+p_{q, 2}^{n+1, m}\right)-p_{q-2,2}^{n+1, m}\right)\right\}-
\end{array}\right.
\end{aligned}
$$

Similarly the equations for $p_{q, I}^{n+1}, m+1$ on the boundary $\beta$ ( eng the line $X_{2}=0$ ) are given as flows,

$$
\left.D u\right|_{q, 1} ^{n+1, m+1}=\frac{2}{\Delta x_{2}}\left[\left(\begin{array}{l}
\left(u^{n+1, m+1}\right. \\
2(q, 2)
\end{array}{\underset{2}{2}(9,1)}_{n+1, m+1}^{u^{2}}-\frac{1}{4}\left(u_{2(q, 3)}^{n+1, m+1}-u_{2(q, 1)}^{n+1, m+1}\right)\right]+\frac{1}{2 \Delta x_{1}}\left(u_{1(q+1,1)}^{n+1, m+1}-u_{1(q-1,1)}^{n+1, m+1}\right)\right.
$$

where,

$$
u_{2(q, 1)}=u_{2(q, 1)}^{u_{1}}
$$

$$
u_{1(4+1,1)}^{u_{1}}=u_{1(9+1,1)}^{a u x}
$$

$$
\begin{aligned}
& n+1, m+1 \\
& u_{1}(q-1,1)=a_{1}(q-1,1)
\end{aligned}
$$

hence,

$$
\begin{aligned}
P_{q, 1}^{n+1, m+1} & =P_{q_{1} 1}^{n+1, m}-\frac{2 \lambda}{\Delta x_{2}} u_{2(q, 2)}^{n+1, m+1}+\frac{2 \lambda}{\Delta x_{2}} u_{2(q, 1)}^{n+1, m+1}+\frac{2 \lambda}{4 \Delta x_{2}} u_{2(q, 3)}^{n+1, m+1}-\frac{2 \lambda}{4 \Delta x_{2} 2(q, 1)} u^{n+1, m+1} \\
& -\frac{\lambda}{2 \Delta x_{1}} u_{1(q+1,1)}^{n+1, m+1}+\frac{\lambda}{2 \Delta x_{1}} u_{1(q-1,1)}^{n+1, m+1}
\end{aligned}
$$

$$
\begin{aligned}
& \underset{\substack{u_{2(1,2)}^{n+1} \\
n+1, m+1}}{\substack{n+1, m+1 \\
u_{2}(q, 12) \\
u_{2(q, 2)}}} \quad 2 \Delta x_{2}, q=1 \\
& u^{n+1, m+1} \quad \text { aux }
\end{aligned}
$$

$$
\begin{aligned}
& -\frac{\lambda}{2 \Delta x_{2}}\left\{{\underset{\sim}{2(q, 3)}}_{\text {aux }}^{u_{2}}-\frac{\Delta t}{2 \Delta x_{2}}\left(p_{q, 4}^{n+1, m}-\frac{1}{2}\left(p_{q, 2}^{n+1, m+1}+p_{q, 2}^{n+1, m}\right)\right)\right\}+\frac{\lambda}{2 \Delta x_{2}} u_{2(q, 1)}^{\text {aux }}, q>2 \\
& \text { ide. } \\
& \int_{q, 2}^{n+1, m+1}=\left(1+\alpha_{1}+\frac{1}{2} \alpha_{2}\right)^{-1}\left[\left(1-\alpha_{1}-\frac{1}{2} \alpha_{2}\right) p_{q, 2}^{n+1, m}-\left.\lambda D u\right|_{q, 2} ^{a u x}+\alpha_{1}\left(p_{q+2,2}^{n+1, m}+p_{q-2,2}^{n+1, m}\right)\right. \\
& \left.+\alpha_{2} p_{q, 4}^{n+1, m}\right], q>2 \\
& P_{4,2}^{n+1, m+1}=\left(1+\frac{\alpha_{1}}{2}+\frac{\alpha_{2}}{2}\right)^{-1}\left[\left(1-\frac{\alpha_{1}}{2}-\frac{\alpha_{2}}{2}\right) p_{q, 2}^{n+1, m}-\lambda D u^{a u x}+\alpha_{1} p_{q+2,2}^{n+1, m}+\alpha_{2} p_{q, 4}^{n+1, m}\right], q=2 \\
& \left.D u^{\text {where }}\right|_{q, 2}=\frac{1}{2 \Delta x_{1}}\left(\begin{array}{c}
a u x \\
\left.u_{1(q+1,2)}-u_{1(q-1,2)}^{a u x}\right)+\frac{1}{2 \Delta x_{2}}\left(\begin{array}{c}
a u x \\
u_{2} \\
(q p)
\end{array}-\underset{2(q, 1)}{a u x}\right) ~
\end{array}\right.
\end{aligned}
$$

i.e
$p_{q, 1}^{n+1, m+1}=p_{q, 1}^{n+1, m}-\frac{2 \lambda}{\Delta x_{2}}\left\{\begin{array}{l}u^{\text {aux }} \\ 2(q, 2)\end{array}-\frac{\Delta t}{2 \Delta x_{2}}\left(p_{q, 3}^{n+1, m}-\frac{1}{2}\left(p_{q, 1}^{n+1, m+1}+p_{q, 1}^{n+1, m}\right)\right)\right\}$

$$
+\frac{2 \lambda}{\Delta x_{2}} u_{2(q, 1)}^{\text {aux }}+\frac{1}{4} \frac{2 \lambda}{\Delta x_{2}}\left\{u_{2(q, 3)}^{a u x}-\frac{\Delta t}{2 \Delta x_{2}}\left(p_{q, 4}^{n+1, m}-p_{q, 2}^{n+1, m}\right)\right\}
$$

$$
-\frac{2 \lambda}{4 \Delta x_{2}} u_{2(q, 1)}^{\text {ac(x }}-\frac{\lambda}{2 \Delta x_{1}} u_{1(q+1,1)}^{\text {aux }}+\frac{\lambda}{2 \Delta x_{1}} u_{1(q-1,1)}^{\text {aux }}, q>1
$$

$p_{q+1, m+1}^{n+1}=\left(1+2 \alpha_{2}\right)^{-1}\left[\left(1-2 \alpha_{2}\right) p_{q, 1}^{n+1, m}-\lambda D u_{q, 1}^{a u x}+4 \alpha_{2}\left(p_{q, 3}^{n+1, m}-\frac{1}{4}\left(p_{q, 4}^{n+1, m} p_{q, 2}^{n+1, m}\right)\right)\right], q \geqslant 1$
$q_{q+1, m+1}^{n+1}=p_{q, 1}^{n+1, m}-\lambda D u^{a q x_{1}} \quad$ c)
other equations canntbe derived on the other boundaries, e.g $p_{I, r}$ i.e the Iine $X_{I}=0$ etc, by similar expresions, $D u_{1,1}^{a u x}=\frac{2}{\Delta x_{2}}\left[\left(u_{2(9,2)}^{a u x}-\frac{u_{2}(9,1)}{a u x}\right)-\frac{1}{4}\left(u_{2(9,3)}^{u_{1}}-u_{2(9,1)}^{u^{a u x}}\right)\right]$ $+\frac{1}{2 \Delta x_{1}}\left(u_{(q+1,1)}^{a u x}-u_{(q-1,1)}^{a u k}\right)$
it is clear that we consider, $u_{i}^{2 u x}$ at the boundaxy is interpreted as $u_{i}^{n+1}$.
The whole iteration system i.e. equations (26 a), $(26 \mathrm{~b}),(26 \mathrm{c})$, converges for $a 11 \lambda>0$ and converges fastest when $\lambda \sim \lambda_{\text {opt }}$.
Because of our representation of $D u=0$, which expresses the balance of mass in a rectangle of sides $2 \Delta x_{i}$, $i=1,2$. The pressure iterations split into to calculations on inter度wined meshes, coupled at the boundary. The most efficient ordering for performing the iterations are such that resulting over-all scheme is a DuFort-Frankel scheme for each one of the intertwined meshes. The iterations are to be done until for some k,

$$
\max _{q, r}\left|p_{q, r}^{n+1, k+1}-p_{q, r}^{n+1, k}\right| \leqq \varepsilon
$$

where $\boldsymbol{\varepsilon}$ is a given small number.
The new velocities $u_{i}^{n+1}, i=1,2$ are to be evaluated using (25 b), (25 c), (25 d), (25 e). This must be done only after $p^{n+1, m}$ are converged. It is also better to evaluating $D u^{a u x}$ at the beging of each iteration. There are two advantages for this iteration procedure (1) $D u^{n+1}$ can be made as small as one desire independently of the error in $D u^{n}$ (2) we could then use latest iterate $p^{n+1, k+1}$ to evaluate $u_{i}^{n+1}$ through formula such as,

$$
u_{i}^{n+1}=u_{i}-G_{i} p^{n+1, k+1}
$$

where $G_{i} p$ approximates $\frac{\partial P}{\partial x_{i}}$. if $\Delta t=O\left(\Delta x^{2}\right)$, when $\mathrm{p}^{\mathrm{n}+1, k+1}$ and $\mathrm{p}^{\mathrm{n}+1, k}$ differ by less than $\boldsymbol{\varepsilon}$, $D u^{n+1}=0(\boldsymbol{\varepsilon} / \lambda)$. Also a gein in accuracy appears, which can use to relax the convergence messure for iterations. This gain in accuracy is due to the fact that $u_{i}^{n+1}$ are evaluated by using an appropriate combination of $p^{n+1, k}$ and $p^{n+1, k+1}$, rather than $p^{n+1, k+1}$. The problem of stability and convergence will be supported by numerical evidence.

## Solution of a Test Problem.

Our method can be applied to a simple-two-dimensional test problem. $D$ is the square $0 \leqq x_{i} \leqq \pi, i=1,2$. The external forces $\mathrm{E}_{1}, \mathrm{E}_{2}$ assumed to be zero. ie. $E_{1}=E_{2}=0$

The boundary data are,
$\left.u_{1}\right|_{x_{1}=0}=-\sin x_{2} e^{-2 t},\left.\quad u_{1}\right|_{x_{2}=0}=0$
$\left.u_{7}\right|_{x_{1}=\pi}=\sin x_{2} e^{-2 t},\left.\quad u_{7}\right|_{x_{2}}=\pi=0$
$\left.u_{2}\right|_{x_{1}=0} \quad,\left.\quad u_{2}\right|_{x_{2}=0}=\sin x_{1} e^{-2 t}$
$\left.u_{2}\right|_{x_{1}=\pi},\left.u_{2}\right|_{x_{2}=\pi}=-\sin x_{1} e^{-2 t}$
also the initial data are,
$\left.u_{1}\right|_{t=0}=-\cos x_{1} \sin x_{2},\left.\quad u_{2}\right|_{t=0}=\sin x_{1} \operatorname{eos} x_{2}$
The exact solution of the problem is.
$u_{1}=-\cos x_{1} \sin x_{2} e^{-2 t}, u_{2}=\sin x_{1} \cos x_{2} e^{-2 t}$

$$
p=-R \frac{I}{4}\left(\cos 2 x_{1}+\cos 2 x_{2}\right) e^{-4 t}
$$

where $R$ is the Reynolds number,

We first evaluate $\lambda_{\text {opt }}$ for the equation,

- $\operatorname{Iu}=\mathrm{f}$
 known on the boundary then,

$$
\omega_{\text {opt }}=\frac{2}{1+\left(1-\alpha^{2}\right)^{1 / 2}}
$$

where $\alpha=\frac{1}{2}\left(\cos 2 \Delta x_{1}+\cos 2 \Delta x_{2}\right)$ is the largest eigenvalue of the associated Jacobi matrix [5], [7]. weput $q=\frac{\lambda_{0 p t}}{2}\left(\frac{\Delta t}{\Delta x_{1}^{2}}+\frac{\Delta t}{\Delta x_{2}^{2}}\right)$
equation (18) can be written in the form

$$
\omega_{o p t}=\frac{8 q}{1+4 q}
$$

then

$$
q=\frac{1}{\left(1-\alpha^{2}\right)^{1 / 2}}
$$

and

$$
\lambda_{o p t}=\frac{4}{\left(\Delta t / \Delta x_{1}^{2}+\Delta t / \Delta x_{2}^{2}\right)} \frac{1}{\left(1-\alpha^{2}\right)^{1 / 2}}
$$

if we assume $\Delta x_{1}=\Delta x_{2}=\Delta x$, then,

$$
\lambda_{o p t}=\frac{2 \Delta x^{2}}{\Delta t \sin (2 \Delta x)}
$$

## CHAPTER III

We record here the list of symbols, Algol 60 listing, and logic flow chart for the two-dimensional viscous flow test problem given at the end of chapter II.
A. List of Symbols
 velocity •
$P 1[Q, R]$
$P 2[Q, R]$
$\mathrm{UI}[\mathrm{Q}, \mathrm{R}]$

U2 $[Q, R]$
$\operatorname{UIN}[Q, R]$
$\mathrm{U} 2 \mathbb{N}[\mathrm{Q}, \mathrm{R}]$
$D A[Q, R]$

IA
1
$u_{2}^{n}(q, r), y-c o m p o n e n t$ of velocity at time step $n$.
$p_{(q, r)}^{n+1, m}$, the pressure before iteration.
$p_{(q, r)}^{n+1, m+1}$, the pressure after iteration.
 time step n.
$u_{1}^{n+1, m+1}$ the computed $x$-component of velocity.
$\left.u_{2}^{n+1, m+1}, r\right)$ the computed $y-c o m p o n e n t$ of velocity.
$D u^{a u x}(q, r)$
$\lambda_{\mathrm{opt}}$ 。

| AI | $\alpha$ |
| :---: | :---: |
| DT | time step $\Delta$ t. |
| Dx | mesh interval $\Delta x$ 。 |
| EPS | $\varepsilon$ |
| MAX | $\begin{aligned} & \max _{q, r}\left\|p_{q, r}^{n+1, m+1}-p_{q, r}^{n+1, m}\right\| \text { at each } \\ & \text { time step. } \end{aligned}$ |
| UTER | absolute difference between the exact and computed value of $u_{I}(q, r)$. |
| U2ER | absolute difference between the exact and computed value of $u_{2}(q, r)$. |
| PER | absolute difference between the exact and computed value of $p$. |
| N | number of time steps. |
| M | number of iterations. |
| Q | $x$-coordinate |
| R | $y$-coordinate |
| T | time |

The Algol listing (Section C) and flow chart (Section B) follow.



Flow clart

Schemes $A$ and $B$, were used for the solution of the test problem, i.e. formulae (6 a), ( 6 b ), (7a), (7b), (7c) were used to evaluate $u_{i}^{a u x}$. $\boldsymbol{\varepsilon}$ is the convergence criterion. In tables $I$, II, $n$ is the number of time steps; $e\left(u_{i}\right)$, $i=$ $=1,2$, are the maxima over $\mathcal{D}$ of the difference between the exact and the computed solution $u_{i}$, $e(p)$ in the tables represents the maximum over the grid of the differences between the exact pressure at time $n \Delta t$ and the computed $p^{n}$ divided by R. The accuracy of the scheme is to be judged by the smallness of $e\left(u_{i}\right) \cdot m$ is the number of iterations.

## Table I



## Table II

| scheme $B ;$ | $\Delta x=\pi / 39 ;$ | $\Delta t=\Delta x^{2} ;$ | $\varepsilon=\Delta x^{2} ;$ | $R=20$ |
| :---: | :--- | :--- | :--- | :---: |
| $n$ | $e\left(u_{1}\right)$ | $e\left(u_{2}\right)$ | $e(p)$ | $m$ |
| 1 | $3 \cdot 9 \times 10^{-3}$ | $4 \cdot 4 \times 10^{-3}$ | 0.0404 | 16 |
| 3 | $5 \cdot 9 \times 10^{-3}$ | $6 \cdot 0 \times 10^{-3}$ | 0.0466 | 11 |
| 5 | $8.5 \times 10^{-3}$ | $6 \cdot 7 \times 10^{-3}$ | 0.0505 | 10 |
| 7 | $1.0 \times 10^{-2}$ | $7 \cdot 4 \times 10^{-3}$ | 0.0551 | 10 |
| 9 | $1.1 \times 10^{-2}$ | $7 \cdot 9 \times 10^{-3}$ | 0.599 | 10 |
| 20 | $1.0 \times 10^{-2}$ | $7.8 \times 10^{-3}$ | 0.0839 | 10 |

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