# F. BREZZI <br> L. D. MARINI <br> On the numerical solution of plate bending problems by hybrid methods 

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# ON THE NUMERICAL SOLUTION OF PLATE BENDING PROBLEMS BY HYBRID METHODS 

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Summary. - We study the convergence of the "assumed stresses hybrid method» of Ptan and Tong for plate bending problems We also give the error bound for a large class of approximations

## INTRODUCTION

Let us consider the «model problem » of an homogeneous isotropic thin plate clamped along the entire boundary and acted by an uniformely distributed load $p$. It is well known that, if $\Omega$ is the portion of the $\left(x_{1}, x_{2}\right)$ plane occupied by the plate, the transversal displacement $w\left(x_{1}, x_{2}\right)$ of the plate is solution of the following boundary value problem :

$$
\left\{\begin{array}{lll}
\Delta^{2} w=p & \text { in } & \Omega  \tag{P}\\
w=\frac{\partial w}{\partial n}=0 & \text { on } & \partial \Omega
\end{array}\right.
$$

where $\partial \Omega$ is the boundary of $\Omega$ and $n$ is the normal outward direction to $\partial \Omega$. The problem $(P)$ is a classical one, and has been studied from a theoretical and numerical point of view by many authors for a long time (see e.g. [17], [18], [44], but the literature on this subject is quite large). The use of finite element methods has recently contributed new developments to the numerical approach
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to the problem ; presently many different types of approximations by means of finite element methods are used : conforming, non conforming, equilibrium, hybrid, and mixed (see e.g. [9], [11], [19], [33], [42], [46] and the bibliography of these papers). We shall treat here the «assumed stresses hybrid method» due to Pian and Tong (see e.g. [32]). Although this method has been used for quite a few years, the first proof of convergence has been given only in the last year in [7] and only for some particular choice of discretisations. In the present paper we construct a very large family of discretisations, including the classical ones, depending on three parameters ( $m=$ degree of the stress field inside each element, $r=$ degree of the displacements at the interelement boundaries, and $s=$ degree of the normal derivatives of the displacements at the interelement boundaries) and we give sufficient conditions on the value of the parameters in order to have convergence (Theorem 3.7.). A bound for the error is also given for each choice of the parameters. Another result of some interest deals with the problem of the search for a « particular solution " $f$ of the equation

$$
f_{i j / i j}=p \quad \text { in each element }
$$

which is needed in order to apply the method. We show that the practical computation of each term of the discrete problem which contains $f$ can be reduced to the computation of integrals of the known function $p$, times some suitable known polynomials; therefore a knowledge of $f$ is not really needed. At the end of the paper we also give some results which have been obtained in the numerical experiments performed by the authors using the Honeywell 6030 of the «Centro di Caicoli Numerici dell'Università di Pavia».

The scheme of the paper is the following.
In paragraph 1 we give a general idea of the assumed stresses hybrid method of Pian and Tong : for any given decomposition $\mathcal{C}_{h}$ of $\Omega$ into convex polygonal subdomains (for the sake of simplicity $\Omega$ is supposed to be a convex polygon), problem $(P)$ is transformed into a saddle point problem $\left(P^{\prime}\right)$ in which two spaces appear : a space of stresses $V\left(\mathcal{G}_{h}\right)$ defined independently in each element and a space (of Lagrangian multipliers) of displacements $W\left(\mathcal{C}_{h}\right)$ defined (essentially) with their first derivatives at the interelement boundaries. A theorem of existence and uniqueness of problem ( $P^{\prime}$ ) is given and its solution is related to the solution of $(P)$.

In paragraph 2 we give, at first, an abstract theorem of convergence (Theorem 2.1) for a general discretisation of $\left(P^{\prime}\right)$ by means of finite dimensional spaces $V_{h} \subset V\left(\mathcal{G}_{h}\right)$ and $W_{h} \subset W\left(\mathcal{G}_{h}\right)$; the convergence is proved if $V_{h}$ and $W_{h}$ satisfy an «abstract» hypothesis, H1. After that we introduce a general family of spaces $V_{h}$ and $W_{h}$; sufficient conditions in order that $V_{h}$ and $W_{h}$, constructed in the indicated manner, satisfy H 1 are given in Theorem 3.7. The last part of paragraph 2 and the whole of paragraph 3 deal with the different steps which lead to the proof of Theorem 3.7. Others different
sufficient conditions can be obtained by means of the « intermediate steps », i.e. Theorem 2.2, Theorem 3.1 and Lemma 3.1.

In paragraph 4 we study the behaviour (in $|h|$ ) of the error bounds given by Theorem 2.1, in the case in which $V_{h}$ and $W_{h}$ are constructed in the indicated manner. We underline in particular Lemma 4.1, first proved by L. Tartar, which is a generalisation of the Bramble-Hilbert lemma and which can be useful in many other situations. A new proof to this lemma is also reported.

In paragraph 5 we study the discrete problem from the computational point of view. We indicate a procedure which avoids the difficulties connected with a knowledge of the particular solution $\underline{f}$; we show that the linear system of equations which is obtained is «equivalent» (in the sense that we have the same degrees of freedom, the same topological matrix, and, with our procedure on $\underline{f}$, the difficulties for the computation of the «known vector» are of the same type) to a «displacement method», conforming or non conforming, but in a more general context. We can find in this way the «hybrid analogues» of the classical «displacement» approaches. At the end of paragraph 5 we give finally the results obtained in some numerical experiments; in particular, we find (see also e.g. [31]) that different choices of the discrete stresses $V_{h}$ do not change the «structure» of the final matrix but they can affect the precision.

## 1. THE HYBRID APPROACH TO THE PROBLEM

Let us consider the problem

$$
\left\{\begin{array}{lll}
\Delta^{2} w=p\left(x_{1}, x_{2}\right) & \text { in } & \Omega  \tag{1.1}\\
w=\frac{\partial w}{\partial n}=0 & \text { on } & \Gamma=\partial \Omega
\end{array}\right.
$$

where $\Omega$ is a convex polygon in the $\left(x_{1}, x_{2}\right)$ plane, $p\left(x_{1}, x_{2}\right)$ is an element of $L^{2}(\Omega)$ and $n$ is the direction of the outward normal to $\Gamma$. It is well known that the solution $w\left(x_{1}, x_{2}\right)$ of (1.1) can be regarded as the displacement along the $x_{1}$-axis of an elastic uniform isotropic plate which is clamped along the entire boundary $\Gamma$ and subjected to an uniformly distributed load $p\left(x_{1}, x_{2}\right)$. Let us introduce the space of stresses

$$
\begin{equation*}
\mathcal{S}=\left\{\underline{v} \mid v_{i j} \in L^{2}(\Omega)(i, j=1,2), v_{12}=v_{21}\right\} \tag{1.2}
\end{equation*}
$$

with the scalar product

$$
\begin{equation*}
[\underline{v}, \underline{u}]=\int_{\Omega} v_{i j} u_{i j} \mathrm{~d} x \tag{1.3}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.
where (here and in the following) the convention of summation of repeated indices is used. The norm of an element $\underline{v}$ in $S$ will be noted by

$$
\begin{equation*}
\|\underline{v}\|_{0}=[\underline{v}, \underline{v}]^{1 / 2} \tag{1.4}
\end{equation*}
$$

We define also the space

$$
\begin{equation*}
S=\left\{\underline{v} \mid \underline{v} \in \mathbf{S}, v_{\imath \jmath / \imath \jmath} \in L^{2}(\Omega)\right\} \tag{1.5}
\end{equation*}
$$

where (here and in the following) the classical notation $g_{j t}=\frac{\partial g}{\partial x_{t}}$ is used.
From the principle of minımum complementary energy we have that if $w$ is the solution of (1.1) then the tensor $\underline{\sigma}$ with components given by

$$
\begin{equation*}
\sigma_{\imath \jmath}=w_{i t J}, \tag{1.6}
\end{equation*}
$$

minimizes the functional

$$
\begin{equation*}
J(\underline{v})=\frac{1}{2} \int_{\Omega} v_{i j} v_{i j} \mathrm{~d} x=\frac{1}{2}\|\underline{v}\|_{0}^{2} \tag{1.7}
\end{equation*}
$$

over the manifold

$$
\begin{equation*}
S_{p}=\left\{\underline{v} \mid \underline{v} \in S, v_{\imath \jmath / \imath \jmath}=p \text { in } \Omega\right\} . \tag{1.8}
\end{equation*}
$$

Let us consider now a decomposition $\mathscr{C}_{h}$ of $\Omega$ into convex subdomains; to the decomposition $\mathfrak{C}_{h}$ we associate the space

$$
\begin{equation*}
\mathcal{V}=\left\{\underline{v} \mid \underline{v} \in S, v_{\mathrm{t} / \mathrm{t} \mathrm{~J}} \in L^{2}(K) \text { for each } K \in \mathscr{C}_{h}\right\} \tag{1.9}
\end{equation*}
$$

with the norm

$$
\begin{equation*}
\|\underline{v}\|_{\vartheta}^{2}=\|\underline{v}\|_{0}^{2}+\sum_{K \in \mathscr{C}_{h}}\left\|v_{l_{j / t}}\right\|_{L^{2}(K)}^{2} \tag{1.10}
\end{equation*}
$$

and we define the continuous bilnear form on $\vartheta \times H_{0}^{2}(\Omega)\left({ }^{1}\right)$ :

$$
\begin{equation*}
b(\underline{v}, \varphi)=\sum_{K \in \mathcal{O}_{h}} \int_{K}\left(v_{\imath \jmath} \varphi_{/ \imath \jmath}-v_{\imath \jmath / \imath \jmath} \varphi\right) \mathrm{d} x . \tag{1.11}
\end{equation*}
$$

It is easy to verify that if $\underline{v} \in \mathcal{V}$ then

$$
\begin{equation*}
b(\underline{v}, \varphi)=0 \quad \forall \varphi \in H_{0}^{2}(\Omega) \tag{1.12}
\end{equation*}
$$

Iff $\underline{v} \in S$. So, introducing the manifold

$$
\begin{equation*}
V_{p}\left(\mathcal{C}_{h}\right)=\left\{\underline{v} \mid \underline{v} \in \mathbb{S}, v_{\imath \jmath / \imath \jmath}=p \text { in each } K \in \mathscr{C}_{h}\right\} \tag{1.13}
\end{equation*}
$$

we have that $S_{p}$ can be presented as

$$
\begin{equation*}
S_{p}=\left\{\underline{v} \mid \underline{v} \in V_{p}\left(\mathscr{C}_{h}\right), b(\underline{v}, \varphi)=0 \quad \forall \varphi \in H_{0}^{2}(\Omega)\right\} \tag{1.14}
\end{equation*}
$$

[^0]The central idea of the stress hybrid method of Pian and Tong [33] is now to minimize $J(\underline{v})$ over $V_{p}\left(\mathcal{G}_{h}\right)$ requiring (1.12) to be satisfied by the method of Lagrangian multipliers. We are therefore led to consider a problem of the type

$$
\left\{\begin{array}{l}
\text { Find }(\underline{\sigma}, \psi) \text { in } V_{p}\left(\mathcal{C}_{h}\right) \times H_{0}^{2}(\Omega), \text { saddle point of }  \tag{1.15}\\
\mathcal{L}(\underline{v}, \varphi)=J(\underline{v})-b(\underline{v}, \varphi) .
\end{array}\right.
$$

By Green's formula, it can be shown that $b(\underline{v}, \varphi)$ depends only on the values of $\underline{v}, \varphi$ and of their first derivatives along the interelement boundaries. So, no uniqueness for $\psi$ in the problem (1.15) can be expected. It is reasonable, then, to «restrict» the space of Lagrangian multipliers by considering the space

$$
\begin{equation*}
W\left(\mathscr{G}_{h}\right)=\left\{\varphi \mid \varphi \in H_{0}^{2}(\Omega), \Delta^{2} \varphi=0 \text { in each } K \in \mathscr{C}_{h}\right\} \tag{1.16}
\end{equation*}
$$

with the norm

$$
\begin{equation*}
\|\varphi\|_{W}^{2}=|\varphi|_{2, \Omega}^{2}=\int_{\Omega} \varphi_{i j} \varphi_{l i j} \mathrm{~d} x . \tag{1.17}
\end{equation*}
$$

Our problem becomes now :

$$
\left\{\begin{array}{l}
\text { Find }(\underline{\sigma}, \psi) \text { in } V_{p}\left(\mathcal{C}_{h}\right) \times W\left(\mathcal{C}_{h}\right) \text { saddle point of : }  \tag{1.18}\\
\mathcal{L}(\underline{v}, \varphi)=J(\underline{v})-b(\underline{v}, \varphi) .
\end{array}\right.
$$

Problem (1.18) can be linearized by considering an element $\underline{f}$ of $V_{p}\left(\mathscr{C}_{h}\right)$ and introducing the space

$$
\begin{equation*}
V\left(\mathcal{C}_{h}\right)=V_{p}\left(\mathcal{C}_{h}\right)-\{\underline{f}\}=\left\{\underline{v} \mid \underline{v} \in \mathbb{S}, v_{i j / i j}=0 \text { in each } K \in \mathcal{C}_{h}\right\} \tag{1.19}
\end{equation*}
$$

Since $V\left(\mathcal{G}_{h}\right)$ is a closed subspace of $S$, it will be equipped with the norm

$$
\begin{equation*}
\|\underline{v}\|_{V}=\|\underline{v}\|_{0} \tag{1.20}
\end{equation*}
$$

It will also be convenient for the following to associate to every $\varphi$ in $W\left(\mathcal{C}_{h}\right)$ a stress tensor $M(\varphi)$ in $V\left(\mathcal{C}_{h}\right)$ defined by

$$
\begin{equation*}
M(\varphi)_{i j}=\varphi_{/ i j} \quad(i, j=1,2) \tag{1.21}
\end{equation*}
$$

we have therefore for all $\varphi$ in $W\left(\mathcal{C}_{h}\right)$ :

$$
\begin{equation*}
\|\varphi\|_{W}=\|M(\varphi)\|_{V}=\|M(\varphi)\|_{0} \tag{1.22}
\end{equation*}
$$

It can be easily verified that (1.18) is equivalent to the following problem :

$$
\begin{cases}\text { find }(u, \psi) \text { in } V\left(\mathcal{G}_{h}\right) \times W\left(\mathfrak{C}_{h}\right) \text { such that : }  \tag{1.23}\\ {[\underline{u}, \underline{v}]+[\underline{f}, \underline{v}]-b(\underline{v}, \psi)=0} & \forall \underline{v} \in V\left(\mathfrak{C}_{h}\right) \\ b(\underline{u}, \varphi)+b(\underline{f}, \varphi)=0 & \forall \varphi \in W\left(\mathcal{C}_{h}\right),\end{cases}
$$

in the sense that $(\underline{\sigma}, \psi)$ is a solution of $(1.18)$ iff $(\underline{u}, \psi)=(\underline{\sigma}-\underline{f}, \psi)$ is a solution of (1.23). We shall give now a theorem that characterizes the relations between the solution of (1.23) and the solution of (1.1).
$n^{\circ}$ décembre 1975, R-3.

Theorem 1.1. - Problem (1.23) has a unique solution ( $\underline{u}, \psi$ ) which is related to the solution $w$ of (1.1) by

$$
\left\{\begin{array}{l}
u_{i j}+f_{i j}=w_{i j j} \quad(i, j=1,2),  \tag{1.24}\\
\psi=w \text { on } \Sigma=\bigcup_{K \in \mathcal{C}_{h}} \partial K \\
\psi_{/ i}=w_{l i} \text { on } \Sigma \quad(i=1,2)
\end{array}\right.
$$

Proof. - Let $(\underline{u}, \psi)$ be a couple in $V\left(\mathcal{C}_{h}\right) \times W\left(\mathcal{G}_{h}\right)$ satisfying conditions (1.24). We have then for all $\underline{v}$ in $V\left(\mathscr{C}_{h}\right)$ :

$$
\begin{align*}
& {[\underline{u}, \underline{v}]+[\underline{f}, \underline{v}]-b(\underline{v}, \psi)=[\underline{u}+\underline{f}, \underline{v}]-b(\underline{v}, \psi)=}  \tag{1.25}\\
& \\
& =\int_{\Omega} w_{/ i j} v_{i j} \mathrm{~d} x-\int_{\Omega} v_{i j} \psi_{/ i j} \mathrm{~d} x=0 .
\end{align*}
$$

Moreover, for all $\varphi$ in $W\left(\mathcal{G}_{h}\right)$,

$$
\begin{equation*}
b(\underline{u}, \varphi)+b(\underline{f}, \varphi)=b(\underline{u}+\underline{f}, \varphi)=0 \tag{1.26}
\end{equation*}
$$

since $\underline{u}+\underline{f}$ belongs to $S$ and $\varphi \in H_{0}^{2}(\Omega)$.
Therefore $(\underline{u}, \psi)$ is a solution of (1.23); let now $\left(\underline{u}^{*}, \psi^{*}\right)$ be another solution of (1.23); with classical arguments we have immediately that $\underline{u}^{*}=\underline{u}$, and therefore

$$
\begin{equation*}
b\left(\underline{v}, \psi-\psi^{*}\right)=0 \quad \forall \underline{v} \in V\left(\mathcal{C}_{h}\right) . \tag{1.27}
\end{equation*}
$$

Hence by taking $\underline{v}=M\left(\psi-\psi^{*}\right)$ in (1.27) we have

$$
\begin{equation*}
b\left(M\left(\psi-\psi^{*}\right), \psi-\psi^{*}\right)=\int_{\Omega}\left(\psi-\psi^{*}\right)_{i j}\left(\psi-\psi^{*}\right)_{i j j} \mathrm{~d} x=0 \tag{1.28}
\end{equation*}
$$

and then $\psi=\psi^{*}$. So (1.23) has a unique solution and the proof is complete.
Remark 1.1. - Existence and uniqueness of the solution of (1.23) follows immediately also from the abstract results of Brezzi [8]. In fact it is sufficient to observe that for all $\varphi$ in $W\left(\mathcal{C}_{h}\right)(\varphi \neq 0)$ :

$$
\operatorname{Sup}_{v \in V\left(\zeta_{h}\right)-\{0\}}\|\underline{v}\|_{0}^{-1} b(\underline{v}, \varphi) \geqslant\|M(\varphi)\|_{0}^{-1} b(M(\varphi), \varphi)=\|M(\varphi)\|_{0}=\|\varphi\|_{W}
$$

## 2. NUMERICAL APPROXIMATION

Let us now consider a sequence $\left\{\mathcal{C}_{h}\right\}_{h}$ of decompositions of $\Omega$ into convex subdomains, and let, for any decomposition $\mathcal{C}_{h}, V_{h}$ and $W_{h}$ be closed subspace of $V\left(\mathscr{C}_{h}\right)$ and $W\left(\mathcal{G}_{h}\right)$ respectively. We consider the approximate problem :

$$
\begin{cases}\text { find }\left(\underline{u}_{h}, \psi_{h}\right) \text { in } V_{h} \times W_{h} \text { such that : }  \tag{2.1}\\ {\left[\underline{u}_{h}+\underline{f}, \underline{v}_{h}\right]-b\left(\underline{v}_{h}, \psi_{h}\right)=0} & \forall \underline{v}_{h} \in V_{h}, \\ b\left(\underline{u}_{h}+\underline{f}, \varphi_{h}\right)=0 & \forall \varphi_{h} \in W_{h} .\end{cases}
$$

Suppose that $V_{h}$ and $W_{h}$ satisfy the following hypothesis :
H1. - There exists a positive constant $\gamma>0$, independent of the decomposition such that :

$$
\begin{equation*}
\operatorname{Sup}_{v \in V_{h}-\{0\}}\|\underline{v}\|_{0}^{-1} b(\underline{v}, \varphi) \geqslant \gamma\|\varphi\|_{W} \quad \forall \varphi \in W_{h} \tag{2.2}
\end{equation*}
$$

Then, from the abstract theory of [8], we get the following result.
Theorem 2.1. - If H 1 is satisfied, (2.1) has a unique solution $\left(\underline{u}_{h}, \psi_{h}\right)$; moreover if $(\underline{u}, \psi)$ is the solution of (1.23) then :

$$
\begin{equation*}
\left\|\underline{u}-\underline{u}_{h}\right\|_{0}+\left\|\psi-\psi_{h}\right\|_{W} \leqslant c\left(\operatorname{Inf}_{\underline{v} \in V_{h}}\|\underline{u}-\underline{v}\|_{0}+\operatorname{Inf}_{\varphi \in W_{h}}\|\psi-\varphi\|_{W}\right) \tag{2.3}
\end{equation*}
$$

where $c$ is a constant independent of the decomposition.
We shall give in the following some general examples of spaces $V_{h}, W_{h}$ which satisfy H1. First of all, suppose that the sequence $\left\{\mathcal{C}_{h}\right\}_{h}$ verifies the following conditions.
c1) There exists a convex polygon $\hat{K}$ such that for every $\mathcal{G}_{h}$ and for every $K$ in $\mathcal{C}_{h}$ we can find a transform $F$ which maps $\hat{K}$ onto $K$, of the type :

$$
\begin{equation*}
x=F \hat{x}=B \hat{x}+\xi \tag{2.4}
\end{equation*}
$$

where $B$ is a $2 \times 2$ non singular matrix and $\xi$ is a vector in $\mathbb{R}^{2}$.
c2) There exists two positive constants $\sigma_{1}, \sigma_{2}$, independent of $\mathscr{C}_{h}$ and of $K$, such that :

$$
\begin{equation*}
\|B\| \leqslant \sigma_{1}|h|,\left\|B^{-1}\right\| \leqslant \sigma_{2}|h|^{-1} \tag{2.5}
\end{equation*}
$$

where $|h|$ is defined, for each decomposition $\mathfrak{C}_{h}$, as

$$
\begin{equation*}
|h|=\max _{K \in \mathcal{O}_{h}}\{\text { diameter of } K\} . \tag{2.6}
\end{equation*}
$$

Let now $\hat{V}$ be a finite dimensional space of smooth symmetric tensors $\underline{v}$, defined on $\hat{K}$ and self-equilibrating, in the sense that :

$$
\begin{equation*}
v_{i j / i j}=0 \text { in } \hat{K} ; \tag{2.7}
\end{equation*}
$$

for any given $K$ in $\mathfrak{G}_{h}$ we define the space

$$
\begin{equation*}
V(\hat{V}, K)=\left\{\underline{v} \mid \exists \underline{\hat{\hat{v}}} \in \hat{V} ; v_{i j}=\left(\hat{v}_{r s} o F\right) b_{i r} b_{j s}\right\} \tag{2.8}
\end{equation*}
$$

where $F$ is the transform of the type (2.4) which maps $\hat{K}$ on $K$ and $b_{l m}$ are the coefficients of the matrix $B$. We have obviously from (2.7) and (2.8) that, for every $\underline{v}$ in $V(\hat{V}, K)$,

$$
\begin{equation*}
v_{i j / i j}=0 \text { on } K . \tag{2.9}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.

Therefore we can define

$$
\begin{equation*}
V\left(\hat{V}, \mathcal{G}_{h}\right)=\left\{\underline{v} \mid \underline{v} \in \mathcal{S}, v_{\left.\right|_{K}} \in V(\hat{V}, K) \forall K \in \mathcal{C}_{h}\right\} \tag{2.10}
\end{equation*}
$$

and we obtain that $V\left(\hat{V}, \mathcal{G}_{h}\right)$ is a closed subspace of $V\left(\mathcal{G}_{h}\right)$.
We define now, for every decomposition $\mathfrak{C}_{h}$ and for every pair $(r, s)$ of integers such that $r \geqslant 3, s \geqslant 1$, the space $W\left(r, s, \mathcal{G}_{h}\right)$ as

$$
\begin{equation*}
W\left(r, s, \mathcal{C}_{h}\right)=\left\{\varphi \mid \varphi \in W\left(\mathcal{C}_{h}\right), \varphi_{\mid \partial K} \in P_{r}(\partial K), \frac{\partial \varphi}{\partial n_{\mid \partial K}} \in P_{s}(\partial K) \forall K \in \mathcal{C}_{h}\right\} \tag{2.11}
\end{equation*}
$$

where $P_{m}(\partial K)$ indicates, for every integer $m$, the space of functions defined on $\partial K$ which are polynomials of degree $\leqslant m$ on each side (and not necessarly continuous). We observe that the condition $« \varphi \in W\left(\mathcal{G}_{h}\right)$ » implies that the functions $\varphi, \varphi_{/ 1}, \varphi_{/ 2}$ are continuous, and this justifies the requirements $r \geqslant 3$, $s \geqslant 1$.

We want now to prove some sufficient conditions on $\hat{V}, r, s$ such that the spaces $V\left(\hat{V}, \mathcal{G}_{h}\right)$ and $W\left(r, s, \mathcal{G}_{h}\right)$ verify condition H1. For the sake of simplicity we will examine, as possible choice for $\hat{K}$, only the classical cases $\hat{K}=$ triangle (and then each $K$ will be a triangle) and $\hat{K}=$ unit square (and then each $K$ will be a parallelogram).

First of all, we remark that each transform $F$ of the type (2.4) can be decomposed in a finite number of ways as the product of :
i) a transform of type :

$$
\begin{align*}
& \bar{x}_{1}=\beta_{11} x_{1} \\
& \bar{x}_{2}=\beta_{21} x_{1}+x_{2} \tag{2.12}
\end{align*}
$$

ii) a contraction :

$$
\begin{align*}
& \overline{\bar{x}}_{1}=\rho \bar{x}_{1}  \tag{2.13}\\
& \overline{\bar{x}}_{2}=\rho x_{2}
\end{align*}
$$

iii) a rigid displacement :

$$
\begin{equation*}
\left(x_{1}, x_{2}\right)=R\left(\bar{x}_{1}, \overline{\bar{x}}_{2}\right) \tag{2.14}
\end{equation*}
$$

Moreover, condition c2) guarantees that there exist four positive constants $\rho_{1}, \rho_{2}, \rho_{3}, \rho_{4}$, independents of $K$ and of $\mathcal{C}_{h}$ such that :

$$
\begin{gather*}
\rho_{1}|h| \leqslant \rho \leqslant \rho_{2}|h|,  \tag{2.15}\\
\beta_{11} \geqslant \rho_{3},  \tag{2.16}\\
\beta_{11}^{2}+\beta_{21}^{2} \leqslant \rho_{4} . \tag{2.17}
\end{gather*}
$$

From the decomposition i), ii), iii) we obtain that each function $\varphi$ such that $\Delta^{2} \varphi=0$ in $K$ is transformed by $F$ into a function

$$
\begin{equation*}
\hat{\varphi}=\varphi \circ F \tag{2.18}
\end{equation*}
$$

such that

$$
\begin{equation*}
A \hat{\varphi}=0 \text { in } \hat{K} \tag{2.19}
\end{equation*}
$$

where $A$ is a fourth order elliptic operator given by

$$
\begin{equation*}
A=\left[\left(\frac{\partial}{\partial x_{1}}-\beta_{21} \frac{\partial}{\partial x_{2}}\right)^{2}+\left(\beta_{11} \frac{\partial}{\partial x_{2}}\right)^{2}\right]^{2} . \tag{2.20}
\end{equation*}
$$

We will now prove some technical lemmas. For this, first of all, we define the operator

$$
\begin{equation*}
\vec{\gamma}: H^{2}(\hat{K}) \rightarrow\left(L^{2}(\partial \hat{K})\right)^{3} \tag{2.21}
\end{equation*}
$$

by

$$
\begin{equation*}
\vec{\gamma}(\varphi)=\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)_{\partial \hat{\mathrm{K}}} \tag{2.22}
\end{equation*}
$$

and the space $T(\partial \hat{K})=\vec{\gamma}\left(H^{2}(\hat{K})\right)\left({ }^{1}\right)$. We define also, for each $A$ of type (2.20) the operator $G_{A}: T(\partial \hat{K}) \rightarrow H^{2}(\hat{K})$ as follows :

$$
\psi=G_{A} \vec{\xi} \Leftrightarrow\left\{\begin{array}{l}
A \psi=0 \text { in } \hat{K}  \tag{2.23}\\
\vec{\gamma} \psi=\vec{\xi} .
\end{array}\right.
$$

Finally we consider the space

$$
\begin{equation*}
\hat{L}=\left\{\underline{v} \mid v_{i j} \in L^{2}(\hat{K})(i, j=1,2), v_{12}=v_{21}\right\} \tag{2.24}
\end{equation*}
$$

with the norm

$$
\begin{equation*}
\|\underline{v}\|_{\hat{L}}^{2}=\|\underline{v}\|_{\Lambda}^{2}=\sum_{i, j=1}^{2}\left\|v_{i j}\right\|_{L^{2}(\hat{K})}^{2}, \tag{2.25}
\end{equation*}
$$

and the operator $\hat{M}$ from $H^{2}(\hat{K})$ into $\hat{L}$ defined as

$$
\begin{equation*}
\underline{v}=\hat{M} \varphi \Leftrightarrow v_{i j}=\varphi_{/ i j} . \tag{2.26}
\end{equation*}
$$

Lemma 2.1. - Let $\left\{A_{\lambda}\right\}_{\lambda \in I}$ be a family of operators of type :

$$
\begin{equation*}
A_{\lambda}=\left(\left(\frac{\partial}{\partial x_{1}}-\beta_{21}(\lambda) \frac{\partial}{\partial x_{2}}\right)^{2}+\left(\beta_{11}(\lambda) \frac{\partial}{\partial x_{2}}\right)^{2}\right)^{2}, \tag{2.27}
\end{equation*}
$$

(1) $T(\partial \hat{K})$ will be equipped with the (natural) norm :

$$
\|\vec{\xi}\|_{T(\hat{C} \hat{K})}=\operatorname{Inf}_{\substack{\varphi \in H^{2}(K) \\ \overrightarrow{\mathcal{F}_{\varphi}=\tilde{\xi}}}}\|\varphi\|_{H^{2}(K)}
$$

with the conditions

$$
\begin{array}{ll}
\beta_{11}(\lambda) \geqslant \bar{\beta} & \forall \lambda \in I \\
\beta_{11}^{2}+\beta_{21}^{2} \leqslant \overline{\bar{\beta}} & \forall \lambda \in I \tag{229}
\end{array}
$$

Then there exists a constant $\alpha$, independent of $\lambda$, such that, for all $\vec{\xi} \in T(\partial \hat{K})$, we have

$$
\begin{equation*}
\left\|\hat{M} G_{A_{\lambda}} \vec{\xi}\right\|_{\Lambda} \leqslant \alpha\left\|\hat{M} G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda} \tag{230}
\end{equation*}
$$

Proof - Let, for all $\vec{\xi}$ in $T(\partial \hat{K})$,

$$
\begin{equation*}
\zeta=G_{\Delta^{2}} \vec{\xi}, \quad \psi=G_{A_{\lambda}} \vec{\xi}, \quad \bar{\psi}=\zeta-\psi \tag{231}
\end{equation*}
$$

We have immediately that

$$
\left\{\begin{array}{l}
A_{\lambda} \bar{\psi}=A_{\lambda} \zeta  \tag{232}\\
\vec{\gamma} \bar{\psi}=0
\end{array}\right.
$$

Therefore, if $a_{\lambda}(u, v)$ is the continuous bilinear form associated with the operator $A_{\lambda}$, we have that $\bar{\psi}$ is solution of the problem

$$
\begin{equation*}
a_{\lambda}(\bar{\psi}, \varphi)=a_{\lambda}(\zeta, \varphi) \quad \forall \varphi \in H_{0}^{2}(\hat{K}) \tag{233}
\end{equation*}
$$

Moreover, from (2 28), (2 29), settıng

$$
\begin{equation*}
|\varphi|_{2, K}=\|\hat{M} \varphi\|_{\Lambda}, \quad \varphi \in H^{2}(\hat{K}) \tag{234}
\end{equation*}
$$

we get immediately that there exist two positive constants $\alpha_{1}, \alpha_{2}$, independent of $\lambda$, such that

$$
\begin{array}{cc}
a_{\lambda}(\varphi, \varphi) \geqslant \alpha_{1}|\varphi|_{2 \hat{K}}^{2} & \forall \varphi \in H_{0}^{2}(\hat{K}) \\
a_{\lambda}\left(\varphi_{1}, \varphi_{2}\right) \leqslant \alpha_{2}\left|\varphi_{1}\right|_{2 \hat{K}}\left|\varphi_{2}\right|_{2, \hat{K}} & \forall \varphi_{1}, \varphi_{2} \in H^{2}(\hat{K}) \tag{236}
\end{array}
$$

So from (233), (235), (236) we get

$$
\begin{equation*}
|\psi|_{2, \hat{K}} \leqslant \alpha_{1}^{-1} \alpha_{2}|\zeta|_{2 \hat{K}} \tag{237}
\end{equation*}
$$

and from (2 31), (2 37)

$$
\begin{equation*}
|\bar{\psi}|_{2, \hat{K}} \leqslant\left(1+\alpha_{1}^{-1} \alpha_{2}\right)|\zeta|_{2 \hat{K}} \tag{238}
\end{equation*}
$$

We have then, for all $\vec{\xi} \in T(\partial \hat{K})$, that

$$
\begin{equation*}
\left\|M G_{A_{\lambda}} \vec{\xi}\right\|_{\Lambda} \leqslant\left(1+\alpha_{1}^{-1} \alpha_{2}\right)\left\|M G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda} \tag{239}
\end{equation*}
$$

which proves (230) with $\alpha=\left(1+\alpha_{1}^{-1} \alpha_{2}\right)$
Lemma 2.2. - Let $\left\{A_{\lambda}\right\}_{\lambda \in \Lambda}$ be a famıly of operators which satisfy (27), (2 28), (2 29), let $\hat{w}$ be a closed cone belonging to a finite dimensional subspace of $T(\partial \hat{K})$, with the property that

$$
\begin{equation*}
\forall \vec{\xi} \in \hat{w}, \forall \vec{\xi}_{1} \in \mathcal{N}, \quad \vec{\xi}+\vec{\xi}_{1} \in \hat{w} \tag{240}
\end{equation*}
$$

where $\mathcal{N}$ is the space of traces of polynomials of degree $\leqslant 1 \mathrm{~m} \hat{K}$, that is :

$$
\begin{equation*}
\mathcal{N}=\vec{\gamma}\left(P_{1}(\hat{K})\right)=\left\{\vec{\xi} \mid \vec{\xi} \in T(\partial \hat{K}), \hat{M} G_{\Delta^{2}} \vec{\xi}=0\right\} . \tag{2.41}
\end{equation*}
$$

Let finally $\hat{V}$ be a finte dimensional subspace of $\hat{L}$ of smooth self-equilibrating tensors, and suppose that the following condition holds:
$\mathbf{H} 1(\hat{V}, \hat{W})\left\{\begin{array}{l}\vec{\forall} \vec{\xi} \in \hat{W}, \\ \text { if } \int_{\hat{\mathbf{x}}}\left(G_{\Delta \Delta} \vec{\xi}\right)_{i, t} v_{l J} \mathrm{~d} \hat{x}=0 \quad \forall \underline{v} \in \hat{V} \\ \text { then }\left\|\hat{M} G_{\Delta{ }^{2}} \vec{\xi}\right\|_{\Lambda}=0 .\end{array}\right.$
Then there exists a positive constant $\hat{\alpha}$, depending only on $\hat{V}$ and $\hat{w}$ such that .

$$
\begin{equation*}
\operatorname{Sup}_{\underline{v} \in \hat{V}-\{0\}}\|\underline{v}\|_{\Lambda}^{-1} \int_{\hat{K}}\left(G_{A_{\lambda}} \vec{\xi}\right)_{/ t J} v_{v y} \mathrm{~d} \hat{x} \geqslant \hat{\alpha}\left\|\hat{M} G_{A_{\lambda}} \vec{\xi}\right\|_{\Lambda}, \tag{2.42}
\end{equation*}
$$

for all $\vec{\xi}$ in $\hat{w}$.
Proof. - First of all we remark that the value of

$$
\begin{equation*}
\int_{\hat{\boldsymbol{K}}}\left(G_{A_{\lambda}} \vec{\xi}\right)_{l y} v_{y y} \mathrm{~d} \hat{x} \tag{2.43}
\end{equation*}
$$

is independent of $\lambda$; in fact, for all $\varphi$ in $H^{2}(\hat{K})$ and for all smooth self equil1brating tensor $\underline{v}$ in $\hat{L}$ we have Green's formula :

$$
\begin{equation*}
\int_{\hat{K}} \varphi_{/ l J} v_{1 J} \mathrm{~d} x=\int_{\partial \hat{\mathrm{K}}}\left(\varphi_{l / 2} v_{l} \hat{n}_{J}-\varphi v_{\imath_{J / J}} \hat{n}_{I}\right) \mathrm{d} \hat{l}, \tag{2.44}
\end{equation*}
$$

where $\hat{\underline{h}}$ is the outward normal direction to $\partial \hat{K}$ and $\mathrm{d} \hat{l}$ is the elementary part of $\hat{\partial} \hat{K}$. Therefore the value of (2.43) depends only on $\vec{\xi}$ and $\underline{v}$ and we can set :

$$
\begin{equation*}
\beta(\underline{v}, \vec{\xi})=\int_{\hat{K}}\left(G_{\Delta \Sigma} \vec{\xi}\right)_{l u v} v_{\imath J} \mathrm{~d} \hat{x} . \tag{2.45}
\end{equation*}
$$

Hence, from lemma 2.1, it is sufficient to show that there exists a positive constant $\bar{\alpha}$ such that, for all $\vec{\xi}$ in $\hat{w}$, we have :

$$
\begin{equation*}
\operatorname{Sup}_{\underline{v} \in \bar{V}-\{0\}}\|\underline{v}\|_{\Lambda}^{-1} \beta(\underline{v}, \vec{\xi}) \geqslant \bar{\alpha}\left\|\hat{M} G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda} . \tag{2.46}
\end{equation*}
$$

For this, we remark that, for all $\vec{\xi}$ in $\hat{w}$ such that $\left\|\hat{M} G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda} \neq 0$, we have from H1 $(\hat{V}, \hat{w})$ that :

$$
\begin{equation*}
\operatorname{Sup}_{\underline{v} \in \hat{V}-\{0\}}\|\underline{v}\|_{\Lambda}^{-1} \beta(\underline{v}, \vec{\xi})=\alpha(\vec{\xi})>0 . \tag{2,47}
\end{equation*}
$$

$\mathbf{n}^{\circ}$ decembre 1975, R-3

Let us now define the sets

$$
\begin{equation*}
\hat{S}=\left\{\vec{\xi} \mid \vec{\xi} \in \hat{w},\left\|\hat{M} G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda}=1\right\} \tag{2.48}
\end{equation*}
$$

and

$$
\begin{equation*}
\hat{D}=\{\vec{\xi} \mid \vec{\xi} \in \hat{S}, \vec{\xi} \text { is orthogonal to } \mathcal{N}\} ; \tag{2.49}
\end{equation*}
$$

we note that $\hat{D}$ is a compact set, since $\hat{w}$ is closed. We also note that, from (2.47) and from the continuity of $\beta(\underline{v}, \vec{\xi}), \alpha(\vec{\xi})$ is continuous, and therefore we have :

$$
\begin{equation*}
\inf _{\vec{\xi} \in \hat{D}} \alpha(\vec{\xi}) \geqslant \bar{\alpha}>0 . \tag{2.50}
\end{equation*}
$$

Let now $\vec{\xi}$ be an element of $\hat{w}$ and let

$$
\begin{equation*}
\rho=\left\|\hat{M} G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda} . \tag{2.51}
\end{equation*}
$$

If $\rho=0$ then (2.46) holds; if $\rho \neq 0$ we have that

$$
\begin{equation*}
\vec{\xi}^{\prime}=\frac{\vec{\xi}}{\rho} \in \hat{S} \tag{2.52}
\end{equation*}
$$

since $\hat{\mathcal{W}}$ is a cone. Let now $\vec{\xi}_{1}$ be the projection of $\vec{\xi}^{\prime}$ on $\mathcal{N}$; we have that $\vec{\xi}^{\prime}-\vec{\xi}_{1} \in \hat{D}$ from (2.40), and moreover :

$$
\begin{align*}
\operatorname{Sup}_{\underline{v} \in \hat{V}-\{0\}}\|\underline{v}\|_{\Lambda}^{-1} \beta(\underline{v}, \vec{\xi})= & \rho \operatorname{Sup}_{\underline{v} \in \hat{V}-\{0\}}\|\underline{v}\|_{\Lambda}^{-1} \beta\left(v, \vec{\xi}^{\prime}\right)=\rho  \tag{2.53}\\
& \operatorname{Sup}_{\underline{v} \in \hat{V}-\{0\}}\|\underline{v}\|_{\Lambda}^{-1} \beta\left(\underline{v}, \vec{\xi}^{\prime}-\vec{\xi}_{1}\right) \geqslant \rho \bar{\alpha}=\bar{\alpha}\left\|\hat{M} G_{\Delta^{2}} \vec{\xi}\right\|_{\Lambda} .
\end{align*}
$$

Therefore (2.46) holds and the proof is complete.
The following theorem gives us a connection between hypothesis $\mathrm{H} 1(\hat{V}, \hat{w})$ and hypothesis H1.

Theorem 2.2. - Let $V\left(\hat{V}, \mathfrak{C}_{h}\right)$ and $W\left(r, s, \mathfrak{G}_{h}\right)$ be constructed as previously stated. For all $\mathfrak{C}_{h}$ and for all $K$ in $\mathfrak{G}_{h}$ we construct the spaces :

$$
\begin{align*}
& W\left(K, \mathfrak{C}_{h}\right)=\left\{\hat{\varphi} \mid \exists \varphi \in W\left(r, s, \mathcal{C}_{h}\right), \hat{\varphi}=\left(\varphi_{\mid K}\right) \circ F\right\},  \tag{2.54}\\
& \mathcal{W}\left(K, \mathcal{C}_{h}\right)=\vec{\gamma}\left(W\left(K, \mathcal{C}_{h}\right)\right), \tag{2.55}
\end{align*}
$$

and we define :

$$
\begin{equation*}
\hat{\mathcal{W}}=\text { closure in } T(\partial \hat{K}) \text { of } \bigcup_{\mathcal{C}_{h}} \bigcup_{K \in \mathcal{F}_{h}} w\left(K, \mathcal{C}_{h}\right) . \tag{2.56}
\end{equation*}
$$

Then, if $\mathrm{H} 1(\hat{V}, \hat{w})$ holds, the spaces $V\left(\hat{V}, \mathcal{G}_{h}\right)$ and $W\left(r, s, \mathcal{G}_{h}\right)$ verify hypothesis H1.

Proof. - Given a decomposition $\mathcal{G}_{h}$, let $\varphi \in W\left(r, s, \mathcal{C}_{h}\right)$; for every $K$ in $\mathcal{G}_{h}$ we define

$$
\begin{equation*}
\hat{\varphi}=\left(\varphi_{\mid K}\right) \circ F ; \tag{2.57}
\end{equation*}
$$

we have obviously $\hat{\varphi} \in W\left(K, \mathcal{G}_{h}\right)$; moreover we have

$$
\begin{equation*}
A_{\lambda} \hat{\varphi}=0 \tag{2.58}
\end{equation*}
$$

for some operator $A_{\lambda}$ of the type (2.27)-(2.29). Hence

$$
\begin{equation*}
\vec{\gamma} \hat{\varphi} \in w\left(K, \mathscr{C}_{h}\right) \subseteq \hat{w} \tag{2.59}
\end{equation*}
$$

and

$$
\begin{equation*}
\hat{\varphi}=G_{A_{\lambda}}(\vec{\gamma} \hat{\varphi}) . \tag{2.60}
\end{equation*}
$$

Then from lemma 2.2 there exists a $\underline{\hat{v}}$ in $\hat{V}$ such that

$$
\left\{\begin{array}{l}
\int_{\hat{\hat{K}}} \hat{\varphi}_{i j} \hat{v}_{i j} \mathrm{~d} \hat{x} \geqslant \hat{\alpha}|\hat{\varphi}|_{2, \hat{K}}^{2},  \tag{2.61}\\
\|\underline{\hat{v}}\|_{\Lambda}^{2} \leqslant|\hat{\varphi}|_{2, \hat{K}} .
\end{array}\right.
$$

We define now an element $\underline{v}=\underline{v}(K)$ in $V(\hat{V}, K)$ by

$$
\begin{equation*}
v_{i j}=\left(\hat{v}_{r s} \circ F\right) b_{i r} b_{j s} \tag{2.62}
\end{equation*}
$$

and we obtain

$$
\begin{align*}
& \int_{K} \varphi_{/ i j} v_{i j} \mathrm{~d} x=\int_{\hat{K}} \hat{\varphi}_{/ i} \hat{v}_{i j}|\operatorname{det} B| \mathrm{d} \hat{x} \geqslant \hat{\alpha}|\hat{\varphi}|_{2, \hat{K}}^{2}|\operatorname{det} B|
\end{aligned} \begin{aligned}
& \geqslant \hat{\alpha}|\varphi|_{2, K}^{2}\left\|B^{-1}\right\|^{-4},  \tag{2.63}\\
& \begin{aligned}
& \int_{K} v_{i j} v_{i j} \mathrm{~d} x= \int_{\hat{K}} \hat{v}_{r s} \hat{v}_{l m} b_{i r} b_{j s} b_{i l} b_{j m}|\operatorname{det} B| \mathrm{d} \hat{x} \leqslant \\
& \leqslant\|\hat{\hat{v}}\|_{\Lambda}^{2}\|B\|^{4}|\operatorname{det} B| \leqslant|\hat{\varphi}|_{2, \hat{K}}^{2}\|B\|^{4}|\operatorname{det} B| \leqslant \\
& \leqslant|\varphi|_{2, K}^{2}\|B\|^{8} .
\end{aligned}
\end{align*}
$$

Let now $\underline{v}$ be the element of $V\left(\hat{V}, \mathfrak{G}_{h}\right)$ defined by :

$$
\begin{equation*}
\underline{v}_{\mid K}=\underline{v}(K) \quad \forall K \in \mathscr{C}_{h} ; \tag{2.65}
\end{equation*}
$$

we have :

$$
\begin{align*}
& b(\underline{v}, \varphi)=\sum_{K \in \mathcal{O}_{h}} \int_{K} \varphi_{l i j} v_{i j} \mathrm{~d} x \geqslant \sum_{K \in \mathcal{O}_{h}} \hat{\alpha}|\varphi|_{2, K}^{2}\left\|B^{-1}\right\|^{-4} \geqslant \hat{\alpha}|\varphi|_{2, \Omega}^{2} \sigma_{2}^{-4}|h|^{4}  \tag{2.66}\\
& \|\underline{v}\|_{0}^{2}=\sum_{K \in \mathcal{C}_{h}} \int_{K} v_{i j} v_{i j} \mathrm{~d} x \leqslant \sum_{K \in \mathcal{O}_{h}}|\varphi|_{2, K}^{2}\|B\|^{8} \leqslant|\varphi|_{2, \Omega}^{2} \sigma_{1}^{8}|h|^{8} .
\end{align*}
$$

Hence :

$$
\begin{equation*}
\frac{b(\underline{v}, \varphi)}{\|\underline{v}\|_{0}} \geqslant \hat{\alpha} \sigma_{2}^{-4} \sigma_{1}^{4}\|\varphi\|_{W} \tag{2.68}
\end{equation*}
$$

and therefore H 1 holds with $\gamma=\hat{\alpha} \sigma_{2}^{-4} \sigma_{1}^{4}$.
$\mathrm{n}^{\circ}$ décembre 1975, R-3.

We characterize now the closed cone $\hat{w}$, or, more generally, some closed cone containing $\hat{w}$, for different choices of the parameters $r, s$ and of $\hat{K}$.

Proposition 2.1. - Let $\hat{K}$ be a triangle or the unit square. Then if $s \geqslant r-1$, $\hat{w}$ is contained in the space $T(r, s, K)$ of triplets $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ which are continuous on $\partial K$ and such that :

$$
\begin{align*}
& \varphi \in P_{r}(\partial \hat{K})  \tag{2.69}\\
& \varphi_{/ i} \in P_{s}(\partial \hat{K}) \quad i=1,2 \tag{2.70}
\end{align*}
$$

The proof is immediate.
Proposition 2.2. - Let $\hat{K}$ be a triangle and suppose that $s \leqslant r-1$. Then $\hat{w}$ is contained in the cone TC $(r, s, \vartheta)$ of triplets of functions $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ that are continuous on $\partial \hat{K}$ and such that the following conditions are satisfied :
i) $\varphi \in P_{r}(\partial \hat{K})$,
ii) $\varphi_{/ i} \in P_{r-1}(\partial \hat{K}) \quad(i=1,2)$,
iii) there exist three directions $\underline{\tilde{n}}^{(1)}, \underline{\tilde{n}}^{(2)}, \underline{\tilde{n}}^{(3)}$ such that :

$$
\left\{\begin{array}{l}
\frac{\partial \varphi}{\tilde{\tilde{n}}^{(i)}} \text { is of degree } \leqslant \text { s on } L_{i},  \tag{2.71}\\
\left|\frac{\tilde{n}^{(i)}}{} \cdot \underline{\tilde{n}}^{(j)}\right| \leqslant \vartheta<1 \quad i \neq j,
\end{array}\right.
$$

where $L_{1}, L_{2}, L_{3}$ are the edges of $\hat{K}$, and $\vartheta$ depends on the constants $\sigma_{1}$ and $\sigma_{2}$ which appear in (2.5).

The proof is immediate, if one consideres $\underline{\tilde{n}}^{(1)}, \tilde{\tilde{n}}^{(2)}, \tilde{\tilde{n}}^{(3)}$ as the images by $F$ of the normal directions $\underline{n}^{(1)}, \underline{n}^{(2)}, \underline{n}^{(3)}$ to the edges of $K$ (of course $\tilde{n}^{(i)}$ depend on $\varphi$ ) and observing that, since $F$ transforms parallel directions into parallel directions, we have $\left|\tilde{\underline{\tilde{n}}}^{(i)} \cdot \tilde{\underline{n}}^{(j)}\right|<1$, and then, if $K$ is « not too flat» (condition c 2 )), we have : $\left|\tilde{\tilde{n}}^{(i)} \cdot \tilde{n}^{(j)}\right| \leqslant \vartheta<1$ for $i \neq j$.

Proposition 2.3. - Let $\hat{K}$ be the unit square and suppose that $s<r-1$. Then $\hat{w}$ is contained in the cone $\operatorname{TCS}(r, s, \vartheta)$ of triplets of functions $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ that are continuous on $\partial \hat{K}$ and satisfy the following conditions :
i) $\varphi \in P_{r}(\partial \hat{K})$,
ii) $\quad \varphi_{/ i} \in P_{r-1}(\partial \hat{K}) \quad(i=1,2)$,
iii') there exist two directions $\underline{\tilde{n}}^{(1)}, \underline{\tilde{n}}^{(2)}$ such that :

$$
\begin{aligned}
& \frac{\partial \varphi}{\partial \tilde{\tilde{n}}^{(1)}} \text { is of degree } \leqslant s \text { on the sides } x_{1}=\mathrm{const}, \\
& \frac{\partial \varphi}{\partial \tilde{n}^{(2)}} \text { is of degree } \leqslant s \text { on the sides } x_{2}=\mathrm{const}, \\
& \left|\underline{\tilde{n}}^{(1)} \cdot \tilde{\tilde{n}}^{(2)}\right| \leqslant \vartheta<1
\end{aligned}
$$

where $\vartheta$ depends on the constants $\sigma_{1}$ and $\sigma_{2}$ which appear in (2.5).
In the following paragraph we shall give some examples of choices of $\hat{V}, r, s, \hat{K}$ such that H1 $(\hat{V}, \hat{w})$ is satisfied.

## 3. EXAMPLES

First of all, we introduce the Southwell functions $U, V$ which are associated with a self-equilibrating stress field $v$ by the formulas

$$
\begin{equation*}
U_{/ 2}=v_{11}, \quad-\frac{1}{2}\left(U_{/ 1}+V_{/ 2}\right)=v_{12}=v_{21}, \quad V_{11}=v_{22}, \tag{3.1}
\end{equation*}
$$

and we consider, for $\varphi$ and $\underline{v}$ sufficiently smooth, the formula :

$$
\begin{equation*}
\beta(\underline{v}, \varphi)=\int_{\hat{\mathbf{K}}} \varphi_{/ i j} v_{i j} \mathrm{~d} \hat{x} ; \tag{3.2}
\end{equation*}
$$

by substituting (3.1) in (3.2) we have :

$$
\begin{equation*}
\beta(\underline{v}, \varphi)=\int_{\hat{K}}\left(\varphi_{/ 11} U_{/ 2}-\varphi_{/ 12} U_{/ 1}-\varphi_{/ 21} V_{/ 2}+\varphi_{/ 22} V_{/ 1}\right) \mathrm{d} \hat{x} \tag{3.3}
\end{equation*}
$$

On the other hand, it is well known that if $f$ and $g$ are sufficiently smooth functions (say, if $f$ and $g$ belong to $H^{2}(\hat{K})$ ) then we have :

$$
\begin{equation*}
\int_{\hat{K}}\left(f_{/ 1} g_{/ 2}-f_{/ 2} g_{/ 1}\right) \mathrm{d} \hat{x}=\int_{\partial \hat{K}} f \frac{\partial g}{\partial s} \mathrm{~d} \hat{l}=-\int_{\partial \hat{K}} \frac{\partial f}{\partial s} g \mathrm{~d} \hat{l} \tag{3.4}
\end{equation*}
$$

where $s$ is the tangent direction to $\partial \hat{K}$.
Therefore from (3.3) and (3.4) we have :

$$
\begin{equation*}
\beta(\underline{v}, \varphi)=\int_{\partial \hat{K}}\left(\varphi_{/ 1} \frac{\partial U}{\partial s}-\varphi_{/ 2} \frac{\partial V}{\partial s}\right) \mathrm{d} \hat{l} . \tag{3.5}
\end{equation*}
$$

We note now that the couple $(U, V)$ associated with a given stress field $\underline{v}$ is not unique, but if $(U, V)$ and $(\bar{U}, \bar{V})$ are associated to the same $\underline{v}$, then

$$
\begin{align*}
& U=\bar{U}+\alpha x_{1}+\beta  \tag{3.6}\\
& V=\bar{V}-\alpha x_{2}+\gamma \tag{3.7}
\end{align*}
$$

with $\alpha, \beta, \gamma$ constants ; therefore

$$
\begin{equation*}
\int_{\partial \hat{K}}\left(\varphi_{/ 1} \frac{\partial U}{\partial s}-\varphi_{/ 2} \frac{\partial V}{\partial s}\right) \mathrm{d} \hat{l}=\int_{\partial \hat{K}}\left(\varphi_{/ 1} \frac{\partial \bar{U}}{\partial s}-\varphi_{/ 2} \frac{\partial \bar{V}}{\partial s}\right) \mathrm{d} \hat{l} \tag{3.8}
\end{equation*}
$$

for each $\varphi$ in $H^{2}(\hat{K})$.
$\mathrm{n}^{\circ}$ décembre 1975, R-3.

We consider now, for each space $\hat{V}$ of smooth symmetrical and selfequilibrating tensors in $\hat{K}$, the space $S F(\hat{V})$ of the traces on $\partial \hat{K}$ of all Southwell functions which are associated to some element $\underline{v}$ in $\hat{V}$. From the previous considerations we obtain immediately the following theorem.

Theorem 3.1. - Let $\hat{V}$ be a space of smooth symmetric self-equilibrating tensors in $\hat{K}$ and let $\hat{w}$ be a closed cone in $T(\partial \hat{K})$. Then $\mathrm{H} 1(\hat{V}, \hat{w})$ is satisfied iff the following condition holds :

$$
C(\hat{V}, \hat{w})\left\{\begin{array}{l}
I f \vec{\xi}=\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right) \in \hat{W} \text { and } \\
\int_{\partial \hat{K}} \varphi_{/ 2} \frac{\partial Z_{i}}{\partial s} \mathrm{~d} \hat{l}=0 \quad \forall\left(Z_{1}, Z_{2}\right) \in S F(\hat{V}) \\
\text { then } \varphi_{/ 2} \text { is constant }(i=1,2) .
\end{array}\right.
$$

The proof is immediate since condition

$$
\begin{equation*}
\varphi_{/ i}=\text { constant } \quad(i=1,2) \tag{3.9}
\end{equation*}
$$

is equivalent to

$$
\begin{equation*}
\left\|\hat{M} G_{\Delta^{2}}\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)\right\|_{\Lambda}=0 \tag{3.10}
\end{equation*}
$$

for all $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ in $T(\partial \hat{K})$.
We will now prove some technical lemmas.
Lemma 3.1. - Let Kbe a (convex) polygon with edges $L_{1}, L_{2}, \ldots, L_{k}$. Let, for any integer $m \geqslant 1$,

$$
\tilde{P}_{m}(\partial K)
$$

be the set of all continuous functions defined on $\partial K$ which are polynomials of degree $\leqslant m$ on each side of $K$. Then if $\varphi$ is an element of $\tilde{P}_{m}(\partial K)$ such that

$$
\begin{equation*}
\int_{\partial K} \varphi \frac{\partial p}{\partial s} \mathrm{~d} l=0 \quad \forall p \in \widetilde{P}_{m}(\partial K) \tag{3.11}
\end{equation*}
$$

we get the existence of $(k+1)$ constants $c, c_{1}, \ldots, c_{k}$ such that :

$$
\begin{equation*}
\varphi=c_{i} l_{i}^{(m)}+c \quad \text { on } \quad L_{i} \quad(i=1, \ldots, k) \tag{3.12}
\end{equation*}
$$

where $l_{l}^{(m)}$ is the (normalized) Legendre polynomial of degree $m$ on $L_{i}$.
Proof. - First of all we remark that if $p_{1}, p_{2}, \ldots, p_{k}$ are $k$ polynomials of degree $\leqslant m-1$ defined respectively on $L_{1}, \ldots, L_{k}$ and such that

$$
\begin{equation*}
\sum_{\imath=1}^{k} \int_{L_{i}} p_{\imath} \mathrm{d} l_{i}=0 \tag{3.13}
\end{equation*}
$$

then there exists a $p \in \widetilde{P}_{m}(\partial K)$ such that :

$$
\begin{equation*}
\frac{\partial p}{\partial s}=p_{\imath} \text { on } L_{\imath} \quad(i=1, \ldots, k) \tag{3.14}
\end{equation*}
$$

which is quite easy to verify. We remark also that, for any given $\varphi$ in $P_{m}(\partial K)$ it is possible to find $(k+1)$ constants $c, c_{1}, \ldots, c_{k}$ and $k$ polynomials $p_{1}, \ldots, p_{k}$ of degree $\leqslant m-1$, such that :

$$
\left\{\begin{array}{l}
\varphi=c_{\imath} l_{\imath}^{(m)}+p_{\imath}+c \text { on } L_{\imath}, \quad(i=1, \ldots, k)  \tag{3.15}\\
\sum_{i=1}^{k} \int_{L_{i}} p_{\imath} \mathrm{d} l_{\imath}=0
\end{array}\right.
$$

Let now $p$ be an element of $\tilde{P}_{m}(\partial K)$ such that

$$
\begin{equation*}
\frac{\partial p}{\partial s}=p_{\imath} \text { on } L_{\imath} \quad(i=1, \ldots, k) \tag{3.16}
\end{equation*}
$$

from (3.11) we get

$$
\begin{equation*}
\int_{\partial K} \varphi \frac{\partial p}{\partial s} \mathrm{~d} l=0 \tag{3.17}
\end{equation*}
$$

and from (3.15), (3.16).

$$
\begin{equation*}
\int_{\partial K} \varphi \frac{\partial p}{\partial s} \mathrm{~d} l=\sum_{t=1}^{k} \int_{L_{\imath}}\left(c_{t} l_{t}^{(m)}+p_{t}\right) p_{\imath} \mathrm{d} l_{t}=\sum_{t=1}^{k} \int_{L_{t}} p_{t}^{2} \mathrm{~d} l_{t} \tag{3.18}
\end{equation*}
$$

Therefore, from (3.17) (3.18), we get $p_{\imath}=0(i=1, \ldots, k)$ and (3.12) is proved.

Lemma 3.2. - In the same hypotheses of lemma 3.1, if $k$ and $m$ are odd, $\varphi$ is constant.

The proof follows immediately from lemma 3.1 and from the antisymmetry of the Legendre polynomials of odd degree.

Lemma 3.3. - Let $K$ be as in lemma 3.1. If $\varphi$ is an element of $\widetilde{P}_{m}(\partial K)$ such that

$$
\begin{equation*}
\int_{\partial K} \varphi \frac{\partial p}{\partial s} \mathrm{~d} l=0 \quad \forall p \in \tilde{P}_{m+1}(\partial K) \tag{3.19}
\end{equation*}
$$

then $\varphi$ is constant.
Proof. - Let us set

$$
\begin{equation*}
\bar{c}=\int_{\partial K} \varphi \mathrm{~d} l, \tag{3.20}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.
obviously :

$$
\begin{equation*}
\sum_{\imath=1}^{k} \int_{L_{t}} p_{\imath} \mathrm{d} l_{\imath}=0 \tag{3.22}
\end{equation*}
$$

Let now $p$ be an element of $\widetilde{P}_{m+1}(\partial K)$ such that

$$
\begin{equation*}
\frac{\partial p}{\partial s}=p_{t} \text { on } L_{t} \quad(i=1, \ldots, k) \tag{3.23}
\end{equation*}
$$

we have from (3.19), (3.21), (3.23)

$$
\begin{equation*}
\int_{\partial K} \frac{\partial p}{\partial s} \mathrm{~d} l=\sum_{i=1}^{k} \int_{L_{i}}\left(p_{i}+\bar{c}\right) p_{\imath} \mathrm{d} l_{i}=\sum_{i=1}^{k} \int_{L_{t}} p_{\imath}^{2} \mathrm{~d} l_{i}=0 . \tag{3.24}
\end{equation*}
$$

So from (3.21), (3.24) we have $\varphi=\bar{c}$.
From lemmas 3.1, 3.2, 3.3 we obtain immediately the following theorems.
Theorem 3.2. - Let $\hat{K}$ be a triangle or the unit square. If $s=r-1$, and $S F(\hat{V})$ contains $\tilde{P}_{s+1}(\partial \hat{K}) \times \tilde{P}_{s+1}(\partial \hat{K})$, taking $\hat{W}=T(r, s, \hat{K})$, condition $C(\hat{V}, \hat{w})$ is satisfied.

Theorem 3.3. - Let $\hat{K}$ be a triangle ; if $s=r-1$, if $S F(\hat{V})$ contains

$$
\widetilde{P}_{s}(\partial \hat{K}) \times \widetilde{P}_{s}(\partial \hat{K})
$$

and if $s$ is odd, taking $\hat{w}=T(r, s, \hat{K})$, condition $C(\hat{V}, \hat{w})$ is satisfied.
We shall now study the case $s<r-1$. For this, suppose first of all that $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ is an element of $T C(r, s, \vartheta)$ such that

$$
\begin{equation*}
\int_{\partial \hat{K}} \varphi_{/ 2} \frac{\partial p}{\partial s} \mathrm{~d} \hat{l}=0 \quad \forall p \in \widetilde{P}_{r-1}(\partial \hat{K}) \quad(i=1,2) \tag{3.25}
\end{equation*}
$$

Then if $(r-1)$ is odd, we immediately get from lemma 3.2 that $\varphi_{/ 2}=$ constant $(i=1,2)$. If $(r-1)$ is even we observe that from (3.25) and lemma 3.1 we get

$$
\begin{array}{ll}
\varphi_{/ 1}=c_{\imath} l_{\imath}^{(r-1)}+c \text { on } L_{\imath} & (i=1,2,3), \\
\varphi_{/ 2}=\bar{c}_{\imath} l_{\imath}^{(r-1)}+\bar{c} \text { on } L_{\imath} & (i=1,2,3) . \tag{3.27}
\end{array}
$$

We also note that, if $(r-1)$ is even, the value of $l_{t}^{(r-1)}$ at the boundaries of $L_{i}$ must be equal, for each $i=1,2,3$. Therefore if, for the sake of simplicity, $\hat{K}$ is an equilateral triangle, we get

$$
\begin{align*}
& c_{1}=c_{2}=c_{3},  \tag{3.28}\\
& \bar{c}_{1}=\bar{c}_{2}=\bar{c}_{3} . \tag{3.29}
\end{align*}
$$

By requiring now that $\frac{\partial \varphi}{\partial \tilde{n}^{(i)}}$ is of degree $\leqslant s$ on $L_{i}(i=1,2,3)$, we have :

$$
\begin{align*}
\operatorname{deg} \frac{\partial \varphi}{\partial \tilde{n}^{(i)}} & =\operatorname{deg}\left(\tilde{n}_{1}^{(i)} \varphi_{/ 1}+\tilde{n}_{2}^{(i)} \varphi_{/ 2}\right)  \tag{3.30}\\
& =\operatorname{deg}\left(\tilde{n}_{1}^{(i)}\left(c_{1} l_{i}^{(r-1)}+c\right)+\tilde{n}_{2}^{(i)}\left(\tilde{c}_{1} l_{i}^{(r-1)}+\tilde{c}\right)\right) \leqslant s<r-1
\end{align*}
$$

therefore

$$
\begin{equation*}
\tilde{n}_{1}^{(i)} c_{1}+\tilde{n}_{2}^{(i)} \bar{c}_{1}=0 \quad(i=1,2,3) \tag{3.31}
\end{equation*}
$$

and, from (2.71), conditions (3.31) imply

$$
\begin{equation*}
c_{1}=\bar{c}_{1}=0 \tag{3.32}
\end{equation*}
$$

and then $\varphi_{/ i}$ is constant $(i=1,2)$.
We have proved the following theorem.
Theorem 3.4. - Let $\hat{K}$ be an (equilateral) triangle; if $s<r-1$ and $S F(\hat{V})$ contains $\tilde{P}_{r-1}(\partial \hat{K}) \times \tilde{P}_{r-1}(\partial \hat{K})$, taking $\hat{\mathcal{W}}=T C(r, s, \vartheta)$, condition $C(\hat{V}, \hat{w})$ is satisfied.

In a similar manner we can also prove the analogous result for $K=$ unit square (the only difference will be, when ( $r-1$ ) is odd, that conditions (3.28), (3.29) become : $c_{1}=-c_{2}=c_{3}=-c$ and $c_{1}=-c_{2}=c_{3}=-c$ ), and the following theorem can be stated.

Theorem 3.5. - Let $K$ be the unit square; if $s<r-1$ and $S F(V)$ contains $P_{r-1}(\partial K) \times P_{r-1}(\partial K)$, taking $w=T C S(r, s, \vartheta)$, condition $C(V, w)$ is satisfied.

We shall now consider the last case, i.e. $s>r-1$. Suppose that $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ is an element of $T(r, s, \hat{K})$ (with $s>r-1$ and $\hat{K}$, for sake of simplicity, equal to the unit square) such that :

$$
\begin{equation*}
\int_{\partial \hat{K}} \varphi_{/ i} \frac{\partial p}{\partial s} \mathrm{~d} \hat{l}=0 \quad \forall p \in \tilde{P}_{s}(\partial \hat{K}) . \tag{3.33}
\end{equation*}
$$

We have from lemma 3.1 that :

$$
\begin{array}{lll}
\varphi_{/ 1}=c_{i} i_{i}^{l s}+c & \text { on } & L_{i} \\
\varphi_{/ 2}=\bar{c}_{i}^{l(s)}+\bar{c} & \text { on } & L_{i} \tag{3.35}
\end{array}
$$

Moreover if, for instance, $L_{1}=\left\{0 \leqslant x_{1} \leqslant 1, x_{2}=0\right\}, c_{1}$ must be equal to zero, since $\varphi_{/ 1}$ coincides, on $L_{1}$, with the derivative of $\varphi$ which is of degree $\leqslant r-1<s$. From the continuity of $\varphi_{/ 1}$ we have then $c_{2}=c_{3}=c_{4}=0$, and $\varphi_{/ 1}$ is constant. In a similar manner we prove that $\varphi_{i 2}$ is constant. The same arguments can be used for the case $\hat{K}=$ triangle

$$
\left\{x_{1} \geqslant 0 ; x_{2} \geqslant 0 ; x_{1}+x_{2} \leqslant 1\right\} .
$$

Therefore we can state the following theorem.
$\mathrm{n}^{\circ}$ décembre 1975, R-3.

Theorem 3.6. - If $\hat{K}$ is the unit square or the triangle

$$
\left\{x_{1} \geqslant 0 ; x_{2} \geqslant 0 ; x_{1}+x_{1} \leqslant 1\right\}
$$

if $s>r-1$ and if $S F(\hat{V})$ contains $\tilde{P}_{s}(\partial \hat{K}) \times \widetilde{P}_{s}(\partial \hat{K})$, taking $\hat{\mathcal{W}}=T(r, s, \hat{K})$, condition $C(\hat{V}, \hat{w})$ is satisfied.

We remark now that if, for all integer $m \geqslant 0$, we set :

$$
\begin{equation*}
P_{m}(\hat{K})=\{\text { polynomials of degree } \leqslant m \text { on } \hat{K}\} \tag{3.36}
\end{equation*}
$$

(3.37) $Q_{m}(\hat{K})=\{$ polynomials of degree $\leqslant m$ in each variable on $\hat{K}\}$,

$$
\begin{array}{r}
S(m, P, \hat{K})=\left\{\underline{v} \mid v_{12}=v_{21}, v_{i j} \in P_{m}(\hat{K}), v_{i j / i j}=0 \text { in } \hat{K}\right\} \\
S(m, Q, \hat{K})=\left\{\underline{v} \mid v_{11}=U_{/ 2}, v_{12}=v_{21}=-\frac{1}{2}\left(U_{/ 1}+V_{l 2}\right),\right.  \tag{3.39}\\
\left.\left.v_{22}=V_{l 1},(U, V) \in \underset{m+1}{Q(\hat{K}}\right)\right\}
\end{array}
$$

then sufficient conditions in order to have that $S F(\hat{V})$ contains

$$
\tilde{P}_{k}(\partial \hat{K}) \times \widetilde{P}_{k}(\partial \hat{K}),
$$

are respectively :
$\hat{V}$ contains $S(k-1, P, \hat{K})$ if $\hat{K}$ is a triangle, $\hat{V}$ contains $S(k-1, Q, \hat{K})$ if $\hat{K}$ is the unit square.

All preceeding results can be summarized in the following theorem.
Theorem 3.7. - Let $\left\{\mathcal{G}_{h}\right\}_{h}$ be a sequence of decompositions which satisfies C 1$)$ and C 2$)$; let, for each decomposition $. \mathcal{G}_{h}, V_{h}=V\left(\hat{V}, \mathcal{G}_{h}\right)$ and $W_{h}=W\left(r, s, \mathfrak{G}_{h}\right)$ be defined starting from $\hat{V}$ and $r, s$ (respectively) as in (2.8), (2.10), (2.11); the following conditions are sufficient in order that $V_{h}$ and $W_{h}$ satisfy hypothesis H1.

1) $\hat{K}$ is a triangle, $\hat{V}$ contains $S(m, P, \hat{K})$ and $m, r, s$ verify the following conditions :

$$
\begin{array}{lll}
m \geqslant r-2 & \text { if } & r-1>s \\
m \geqslant r-2 & \text { if } & r-1=s \text { and } s \text { is odd } \\
m \geqslant r-1 & \text { if } & r-1=s \text { and } s \text { is even } \\
m \geqslant s-1 & \text { if } & r-1<s
\end{array}
$$

2) $\hat{K}$ is the unit square, $\hat{V}$ contains $S(m, Q, \hat{K})$ and $m, r, s$ verify the following conditions :

$$
\begin{array}{lll}
m \geqslant r-2 & \text { if } & r-1>s \\
m \geqslant r-1 & \text { if } & r-1=s \\
m \geqslant s-1 & \text { if } & r-1<s
\end{array}
$$

where $S(m, P, \hat{K})$ and $S(m, Q, \hat{K})$ are defined in (3.36)-(3.39).

Remark. - Using lemma 3.1 it is possible to find many other sufficient conditions such that hypothesis H1 is satisfied. For example if we require that the functions $\varphi$ of $\hat{w} \subseteq T(r, s, \hat{K})$, for $s=r-1$ and $\hat{K}=$ unit square, are such that « $\varphi_{/ 12}$ is continuous at each vertex», we can prove that if $S F(\hat{V})$ contains $\widetilde{P}_{s}(\partial \hat{K}) \times \widetilde{P}_{s}(\partial \hat{K})$ and $s$ is odd then $C(\hat{V}, \hat{w})$ holds.

## 4. ERROR BOUNDS

We want now to study, in the cases of the examples treated before, the behaviour in $|h|$ of the error bounds:

$$
\begin{equation*}
\operatorname{Inf}_{\underline{v} \in V_{h}}\|\underline{u}-\underline{v}\|_{0} \tag{4.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\operatorname{Inf}_{\varphi \in W_{h}}\|\psi-\varphi\|_{W} \tag{4.2}
\end{equation*}
$$

We shall prove, as an upper bound for (4.1), the following result.
Theorem 4.1. - If $V_{h}=V\left(\hat{V}, \mathcal{G}_{h}\right)$ and $\hat{V}$ contains the space $S(m, P, \hat{K})$, then, if $u_{i j \mid K} \in H^{m+1}(K)(i, j=1,2)$ for all $K$ in $\mathcal{C}_{h}$, we have :

$$
\begin{equation*}
\operatorname{Inf}_{\underline{v} \in V_{h}}\|\underline{u}-\underline{v}\|_{0} \leqslant C|h|^{m+1}\|\underline{u}\|_{m+1, ⿻_{i}} \tag{4.3}
\end{equation*}
$$

where $C$ is a constant independent of $u$ and $\mathcal{G}_{h}$, and $\|\underline{u}\|_{m+1, \mathcal{F}_{h}}$ is given by

$$
\begin{equation*}
\|\underline{u}\|_{m+1, \mathscr{C}_{h}}^{2}=\sum_{i, j=1}^{2} \sum_{K \in \mathscr{G}_{h}}\left\|u_{i j}\right\|_{H^{m+1}(K)}^{2} \tag{4.4}
\end{equation*}
$$

The proof of theorem 4.1 is based on an abstract lemma (cf. [43]) of the Bramble-Hilbert type which will be reported here in the most general form since it can be useful in many other situations. We present here a different proof of the same result.

Lemma 4.1. - Let $E$ be a Banach space and let $E_{0}, E_{1}, F$ three normed linear spaces; let moreover $A_{0}, A_{1}, L$ be linear continuous operators from $E$ into $E_{0}, E_{1}, F$ respectively. If :

$$
\begin{align*}
& \|g\|_{E} \simeq\left\|A_{0} g\right\|_{E_{0}}+\left\|A_{1} g\right\|_{E_{1}},  \tag{4.5}\\
& L g=0 \quad A_{1} g=0  \tag{4.6}\\
& A_{0} \text { is compact } \tag{4.7}
\end{align*}
$$

then there exists a constant $c$ such that

$$
\begin{equation*}
\|L g\|_{F} \leqslant c\left\|A_{1} g\right\|_{E_{1}} \quad \forall g \in E \tag{4.8}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.

Proof. - First of all we remark that from (4.5) and (4.7) we have that

$$
\begin{equation*}
P=\left\{g \mid g \in E, A_{1} g=0\right\}=\operatorname{ker}\left(A_{1}\right) \tag{4.9}
\end{equation*}
$$

is a finite dimensional subspace of $E$. In fact if $g_{n}$ is a sequence in $P$ such that (4.10)

$$
g_{n} \rightharpoonup 0 \text { in } E
$$

we have from (4.7) that

$$
\begin{equation*}
A_{0} g \rightarrow 0 \text { in } E_{0} \tag{4.11}
\end{equation*}
$$

and therefore from (4.5)

$$
\begin{equation*}
g_{n} \rightarrow 0 \text { in } E . \tag{4.12}
\end{equation*}
$$

The proof of the lemma will be now given in two steps. We shall first prove that, defining, for all $g$ in $E$,

$$
\begin{equation*}
Q(g)=\operatorname{Inf}_{p \in P}\|g-p\|_{E} \tag{4.13}
\end{equation*}
$$

we get that there exists a constant $c_{1}$, such that

$$
\begin{equation*}
Q(g) \leqslant c_{1}\left\|A_{1} g\right\|_{E_{1}} \quad \forall g \in E . \tag{4.14}
\end{equation*}
$$

Secondly we shall prove that there exists a constant $c_{2}$ such that :

$$
\begin{equation*}
\|L g\| \leqslant c_{2} Q(g) \quad \forall g \in E . \tag{4.15}
\end{equation*}
$$

In order to prove (4.14) we suppose, by contradiction, that there exists a sequence $g_{n}$ in $E$ such that

$$
\begin{align*}
& \left\|A_{1} g_{n}\right\|_{E_{1}} \rightarrow 0,  \tag{4.16}\\
& Q\left(g_{n}\right)=1 . \tag{4.17}
\end{align*}
$$

Then, since $P$ is finite dimensional, there exists a sequence $\widetilde{g}_{n}=g_{n}-p_{n}$ defined by

$$
\begin{equation*}
\left\|\tilde{g}_{n}\right\|_{E}=Q\left(g_{n}\right)=\operatorname{Inf}_{p \in P}\left\|g_{n}-p\right\|_{E}=\left\|g_{n}-p_{n}\right\|_{E} ; \tag{4.18}
\end{equation*}
$$

we have now

$$
\begin{align*}
& \left\|A_{1} \tilde{g}_{n}\right\|_{E_{1}}=\left\|A_{1} g_{n}\right\|_{E_{1}} \rightarrow 0,  \tag{4.11}\\
& \left\|\tilde{g}_{n}\right\|_{E}=1 . \tag{4.2}
\end{align*}
$$

Therefore, there exists a subsequence $\left\{\tilde{g}_{k}\right\}$ of $\left\{\tilde{g}_{n}\right\}$ and an element $g^{*}$ in $E$ such that :

$$
\begin{equation*}
\tilde{g}_{k} \rightharpoonup g^{*} \text { in } E \text {, } \tag{4.21}
\end{equation*}
$$

and consequently :

$$
\begin{align*}
& A_{0} \tilde{g}_{k} \rightarrow A_{0} g^{*} \text { in } E_{0},  \tag{4.22}\\
& A_{1} g_{k} \rightarrow A_{1} g^{*} \text { in } E_{1} . \tag{4.23}
\end{align*}
$$

Then from (4.23) and (4.19)

$$
\begin{equation*}
A_{1} \tilde{g}_{k} \rightarrow A_{1} g^{*}=0 \text { in } E_{1} \tag{4.24}
\end{equation*}
$$

and from (4.5), (4.22), (4.24) we have

$$
\begin{align*}
& g_{k} \rightarrow g^{*} \text { in } E  \tag{4.25}\\
& \tilde{g}^{*} \in P . \tag{4.26}
\end{align*}
$$

Finally we have from (4.25), (4.26) that

$$
\operatorname{Inf}_{p \in P}\left\|\tilde{g}_{k}-p\right\|_{E} \leqslant\left\|\tilde{g}_{k}-g^{*}\right\| \rightarrow 0
$$

which is in contradiction to $(4.17),(4.18)$; hence $(4.14)$ is proved.
Let us prove now (4.15); we have for all $g$ in $E$ and for all $p$ in $P$ :

$$
\begin{equation*}
\|L g\|_{F}=\|L g-L p\| \leqslant c_{2}\|g-p\|_{E} \tag{4.27}
\end{equation*}
$$

since $L$ is continuous. Taking the «Inf» in (4.26) we get immediately (4.15). In order to apply lemma 4.1 to our case, let us consider for any integer $k$ the space :

$$
\begin{equation*}
\hat{E}^{(k)}=\left\{\underline{v} \mid v_{12}=v_{21}, v_{\imath j} \in H^{k}(\hat{K}) \quad(\imath, j=1,2)\right\} . \tag{4.28}
\end{equation*}
$$

Setting now for any integer $m \geqslant 1$

$$
\left\{\begin{array}{l}
E=\hat{E}^{(m+1)}, \quad E_{0}=\hat{E}^{(0)}, \quad E_{1}=\left(L^{2}(\hat{K})\right)^{7}, \quad F=E  \tag{4.29}\\
A_{0}=I \text { (identity) } \\
A_{1} \underline{v}=\left(\frac{\partial^{m+1}}{\partial x_{1}^{m+1}} \underline{v}, \frac{\partial^{m+1}}{\partial x_{2}^{m+1}} \underline{v}, v_{i / / l}\right) \\
\hat{\Pi}=\text { projection over } S(m, P, \hat{K}) \\
L=I-\hat{\Pi}
\end{array}\right.
$$

we are exactly in the hypotheses of lemma 4.1, since $A_{1} \underline{v}=0$ implies $\underline{v} \in S(m, P, \hat{K})$. Then we get that there exists a constant $\hat{c}$ such that

$$
\begin{equation*}
\|\underline{u}-\hat{\Pi} \underline{u}\| \leqslant \hat{c} \sum_{i, j=1}^{2}\left|u_{i v}\right|_{m+1, \hat{K}} \tag{4.30}
\end{equation*}
$$

for all $\underline{u}$ in $\hat{E}^{(m+1)}$ such that $u_{i j / \imath \jmath}=0$ in $\hat{K}$.
We can now prove theorem 4.1.
Proof of theorem 4.1. - Let $\underline{u} \in V\left(\mathcal{C}_{h}\right)$, and suppose that for all $K$ in $\mathscr{C}_{h}$ we have :

$$
\begin{equation*}
\left(u_{\imath \jmath}\right)_{\mid K} \in H^{m+1}(K) \quad(i, j=1,2) \tag{4.31}
\end{equation*}
$$

Then, for any given $K$ in $\mathcal{G}_{h}$, consider the tensor $\underline{\hat{u}}=\underline{\hat{u}}(K) \in \hat{E}^{(m+1)}$ such that:

$$
\begin{equation*}
\left(u_{\imath \jmath}\right)_{\mid K}=\left(\hat{u}_{r s} \circ F\right) b_{\imath r} b_{\jmath s} \tag{4.32}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3

Of course we have :

$$
\begin{equation*}
\hat{u}_{r s / r s}=0 \tag{4.33}
\end{equation*}
$$

since $\underline{u}$ is self-equilibrating on $K$. Let now :

$$
\begin{equation*}
\underline{\hat{v}}=\underline{\hat{v}}(K)=\hat{\Pi} \underline{\hat{u}} ; \tag{4.34}
\end{equation*}
$$

from (4.29) we obtain :

$$
\begin{equation*}
\|\underline{\hat{u}}-\underline{\hat{v}}\|_{\Lambda} \leqslant \hat{c} \sum_{i, j=1}^{2}\left|\hat{u}_{i j}\right|_{m+1, \hat{\kappa}} \tag{4.35}
\end{equation*}
$$

Let now $\underline{v}(K)$ be defined by :

$$
\begin{equation*}
v_{\imath \jmath}(K)=\left(\hat{v}_{r s} \circ F\right) b_{t r} b_{J s} \tag{4.36}
\end{equation*}
$$

we have :

$$
\begin{align*}
\int_{K}\left(u_{\imath j}-v_{\imath \jmath}\right)\left(u_{i \jmath}-v_{\imath \jmath}\right) \mathrm{d} x & \leqslant\|\underline{u}-\underline{v}\|_{\Lambda}^{2}\|B\|^{4}|\operatorname{det} B| \leqslant  \tag{4.37}\\
& \leqslant \hat{c}^{2} \sum_{i, j=1}^{2}\left|\hat{u}_{\imath \jmath}\right|_{m+1, \hat{K}}^{2}\|B\|^{4}|\operatorname{det} B| \leqslant \\
& \leqslant \hat{c}^{2} \sum_{i, j=1}^{2}\left|u_{\imath \jmath}\right|_{m+1, K}^{2}\|B\|^{4}\left\|B^{-1}\right\|^{4}\|B\|^{2 m+2} \leqslant \\
& \leqslant \hat{c}^{2} \sum_{i, j=1}^{2}\left|u_{\imath \jmath}\right|_{m+1, K}^{2} \sigma_{1}^{2 n+6} \sigma_{2}^{4}|h|^{2 m+2}
\end{align*}
$$

Let now $\underline{v}$ be the element of $V\left(\hat{V}, \mathcal{C}_{h}\right)$ defined by :

$$
\begin{equation*}
\underline{v}_{\mid K}=v(K) \quad \forall K \in \mathscr{C}_{h} . \tag{8}
\end{equation*}
$$

We will have :

$$
\begin{align*}
& \|\underline{u}-\underline{v}\|_{0}^{2}=\sum_{K \in \overparen{C}_{h}} \int_{K}\left(u_{\imath \jmath}-v_{\imath \jmath}\right)\left(u_{\imath \jmath}-v_{\imath \jmath}\right) \mathrm{d} x \leqslant  \tag{4.39}\\
\leqslant & \sum_{K \in \mathscr{Y}_{h}} c^{2} \sum_{i, j=1}^{2}\left|u_{\imath \jmath}\right|_{m+1, K}^{2} \sigma_{1}^{2 m+6} \sigma_{2}^{4}|h|^{2 m+2} \leqslant \hat{c}^{2} \sigma_{1}^{2 m+6} \sigma_{2}^{4}|h|^{2 m+2}\|\underline{u}\|_{m+1, \mathfrak{Y}_{h}}^{2} .
\end{align*}
$$

and (4.3) is proved.
Let us examine, now, the quantity (4.2), that is:

$$
\begin{equation*}
\operatorname{Inf}_{\varphi \in W_{h}}\|\psi-\varphi\|_{W} \tag{4.40}
\end{equation*}
$$

We have the following theorem.
Theorem 4.2. - If $W_{h}=W\left(r, s, \mathfrak{C}_{h}\right)$ and if $\psi_{\mid K} \in H^{9+2}(K)$ for all $K$ in $\mathfrak{G}_{h}$, we have

$$
\begin{equation*}
\operatorname{Inf}_{\varphi \in W_{h}}\|\psi-\varphi\|_{W} \leqslant c|h|^{q}\left(\sum_{K \in \mathcal{F}_{h}}\|\psi\|_{q+2, K}^{2}\right)^{1 / 2} \tag{4.41}
\end{equation*}
$$

where $q=\min (s, r-1)$ and $c$ is a constant independent of $\psi$ and of $\mathscr{C}_{h}$.

Proof. - We shall give only a sketch of the proof, which is based essentially on the Bramble-Hilbert lemma. First of all we remark that :

$$
\begin{equation*}
\|\psi-\varphi\|_{W}^{2}=\sum_{K \in \mathcal{C}_{h}}|\psi-\varphi|_{2, K}^{2}, \tag{4.42}
\end{equation*}
$$

and that

$$
\begin{equation*}
|\psi-\varphi|_{2, K}^{2}=\operatorname{Inf}_{\chi-(\psi-\varphi) \in H_{0}^{2}(K)}|\chi|_{2, K}^{2} \tag{4.43}
\end{equation*}
$$

since $\psi$ and $\varphi$ are biharmonic in $K$.
Therefore if we consider the space

$$
\begin{equation*}
H\left(r, s, \mathcal{G}_{h}\right)=\left\{\chi \mid \chi \in H_{0}^{2}(\Omega), \chi_{\mid \partial K} \in P_{r}(\partial K), \frac{\partial \chi}{\partial n_{\mid \partial K}} \in P_{s}(\partial K), \forall K \in \mathcal{C}_{h}\right\} \tag{4.44}
\end{equation*}
$$

we obtain :

$$
\begin{equation*}
\operatorname{lnf}_{\varphi \in W\left(r, s, \mathscr{F}_{h}\right)}\|\psi-\varphi\|_{W} \leqslant \operatorname{Inf}_{\chi \in H\left(r, s, \mathscr{F}_{h}\right)}|\psi-\chi|_{2, \Omega} \tag{4.45}
\end{equation*}
$$

Now, if $s=r-1$, it is known (see e.g. [5], [10], [15], [36]) that :

$$
\begin{equation*}
\operatorname{Inf}_{\chi \in H\left(r, s \mathscr{C}_{h}\right)}|\psi-\chi|_{2, \Omega} \leqslant c|h|^{r-1}\left(\sum_{K \in \mathscr{C}_{h}}\|\psi\|_{r+1, K}^{2}\right)^{1 / 2} \tag{4.46}
\end{equation*}
$$

with $c$ constant independent of $\psi$ and $\mathscr{G}_{h}$. Then the obvious inclusions:

$$
\begin{array}{ll}
H\left(r, s, \mathcal{G}_{h}\right) \subseteq H\left(r, r-1, \mathcal{C}_{h}\right) & (s \leqslant r-1), \\
H\left(r, s, \mathcal{C}_{h}\right) \subseteq H\left(s+1, s, \mathcal{C}_{h}\right) & (s \geqslant r-1), \tag{4.48}
\end{array}
$$

conclude the proof in the case $s \geqslant 2$.
Let us now consider the case $s=1$; since for all $r \geqslant 3$ we have

$$
\begin{equation*}
W\left(r, 1, \mathfrak{C}_{h}\right) \supseteq W\left(3,1, \mathfrak{G}_{h}\right) \tag{4.49}
\end{equation*}
$$

it will be sufficient to study the case $r=3$. To this aim, let $\varphi \in W\left(3,1, \mathcal{C}_{h}\right)$ be the function defined by

$$
\left\{\begin{array}{l}
\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)=\left(\psi, \psi_{/ 1}, \psi_{2}\right)  \tag{4.50}\\
\text { at each vertex of each } K \in \mathcal{G}_{h}
\end{array}\right.
$$

and let $\chi \in H\left(3,2, \mathfrak{C}_{h}\right)$ be the function of Clough-Tocher type (cfr. e.g. [10]) which verifies (4.50) and also verifies

$$
\begin{equation*}
\frac{\partial \chi}{\partial n}=\frac{\partial \psi}{\partial n} \text { at the middle point of each side of each } K \in \mathcal{C}_{h} . \tag{4.51}
\end{equation*}
$$

It is now well known (cfr. again [10]) that :

$$
\begin{equation*}
|\psi-\chi|_{2, \Omega} \leqslant c|h|\left(\sum_{K \in \mathcal{F}_{h}}\|\psi\|_{3, K}^{2}\right)^{1 / 2} \tag{4.52}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.
with $c$ constant independent of $\psi$ and $\mathscr{G}_{h}$. Considering now the function $\varphi-\chi$, we obvious have that

$$
\begin{equation*}
\varphi-\chi=0 \quad \text { on } \quad \Sigma=\bigcup_{K \in \mathcal{C}_{h}} \partial K \tag{4.53}
\end{equation*}
$$

and for all $K$ in $\mathfrak{G}_{h}$ the function

$$
\begin{equation*}
\frac{\partial}{\partial n}(\varphi-\chi)_{\mid \partial K} \tag{4.54}
\end{equation*}
$$

is a quadratic polynomial on each side which is zero at each corner and which values on the middle point of each side are bounded by some constant (independent of $K$ ) times the norm of $\psi$ in $W^{1, \infty}(K)$. Therefore with classical arguments (cfr. e.g. [36], chap. VI, pp. 24-25) we get

$$
\begin{equation*}
|\varphi-\chi|_{2, \Omega} \leqslant c|h|\left(\sum_{K \in \mathcal{G}_{h}}\|\psi\|_{3, K}^{2}\right)^{1 / 2} \tag{4.55}
\end{equation*}
$$

with $c$ constant independent of $\psi$ and $\mathscr{C}_{h}$. Therefore from (4.52) and (4.55) we get the result for $s=1$ and the proof is completed.

## 5. NUMERICAL SOLUTION

We shall now make some remarks on the effective computation of the solution of the «discretized problem » (2.1). It is easy to verify that, with the indicated choice for $V_{h}$ and $W_{h}$, (2.1) reduces to a system of $N V+N W$ linear equations with $N V+N W$ unknowns, where, of course, $N V$ is the dimension of $V_{h}$ and $N W$ is the dimension of $W_{h}$. In particular, let $\underline{v}^{(1)}, \ldots, \underline{v}^{(N)}$ be a basis in $V_{h}$ and $\varphi^{(1)}, \ldots, \varphi^{(N W)}$ a basis in $W_{h}$; writing .

$$
\begin{align*}
& \underline{u}_{h}=\sum_{i=1}^{N V} U_{i} \underline{v}^{(i)}  \tag{5.1}\\
& \Psi_{h}=\sum_{i=1}^{N W} \Psi_{i} \varphi^{(i)} \tag{5.2}
\end{align*}
$$

problem (2.1) becomes

$$
\left\{\begin{array}{l}
\sum_{i=1}^{N V}\left[\underline{v}^{(i)}, \underline{v}^{(J)}\right] U_{i}-\sum_{i=1}^{N W} b\left(\underline{v}^{(J)}, \varphi^{(i)}\right) \Psi_{i}=-\left[\underline{f}, \underline{v}^{(j)}\right] \quad j=1, \ldots, N V  \tag{5.3}\\
\sum_{i=1}^{N V} b\left(\underline{v}^{(i)}, \varphi^{(\jmath)}\right) U_{i}=-b\left(\underline{f}, \varphi^{(\jmath)}\right) \quad j=1, \ldots, N W
\end{array}\right.
$$

and setting :

| (5.4) | $\mathcal{A}=\left\{A_{\imath \jmath}\right\} ;$ | $A_{i j}=\left[\underline{v}^{(j)}, \underline{v}^{(2)}\right]$ | $(\imath=1, \ldots, N V, j=1, \ldots, N V)$, |
| :---: | :---: | :---: | :---: |
| (5.5) | $\mathfrak{B}=\left\{B_{\imath}\right\}$; | $B_{i j}=b\left(\underline{v}^{(p)}, \varphi^{(\imath)}\right)$ | $(l=1, \ldots, N W ; \jmath=1, \ldots, N V)$, |
| (5.6) | $F_{i}=-\left[\underline{f}, \underline{v}^{(t)}\right]$ | $(l=1, \ldots, N V)$, |  |
| (5.7) | $G_{i}=-b\left(\underline{f}, \varphi^{(2)}\right)$ | $(i=1, \ldots, N W)$, |  |

problem (5.3) can be written as :

$$
\left\{\begin{array}{l}
\mathcal{A} \vec{U}-\mathfrak{B}^{T} \vec{\Psi}=\vec{F}  \tag{5.8}\\
\mathfrak{B} \vec{U}=\vec{G}
\end{array}\right.
$$

where $\mathcal{A}$ is a square $N V \times N V$ matrix, symmetric and positive definite, $\mathfrak{B}$ a rectangle $N W \times N V$ matrix and $\mathcal{B}^{T}$ is the transposed of matrix $\mathfrak{B}$.

The following part of the paragraph will be divided into four steps : $a$ ) computation of the coefficients, $b$ ) computation of the «known vectors» $\vec{F}$ and $\vec{G}, c$ ) solution of the linear system, $d$ ) numerical results.
a) Computation of the coefficients

The practical computation of the coefficients $A_{i j}$ of the matrix $\mathcal{A}$ can be performed without difficulty, since

$$
\begin{equation*}
\left[\underline{v}^{(j)}, \underline{v}^{(i)}\right]=\sum_{K \in \overparen{\mho}_{h}} \int_{K} v_{r s}^{(j)} v_{r s}^{(i)} \mathrm{d} x \tag{5.9}
\end{equation*}
$$

and $v_{r s}^{(l)}$ are smooth known functions (in general polynomials) in each $K$. Moreover, since the elements $\underline{v}$ of $V_{h}$ are independently assumed in each $K$, we have that, if $V_{h}=V\left(\hat{V}, \mathcal{C}_{h}\right)$ and if $n$ and $N(h)$ are respectively the dimension of $\hat{V}$ and the number of elements $K$ in $\mathcal{G}_{h}$, then the dimension $N V$ of $V_{h}$ is given by :

$$
\begin{equation*}
N V=n . N(h) \tag{5.10}
\end{equation*}
$$

Let now $\underline{\hat{v}}^{(1)}, \ldots, \underline{\hat{v}}^{(n)}$ be a basis in $\hat{V}$ and let, for each $K$ in $\mathcal{G}_{h}, \underline{v}^{(1)}(K), \ldots, \underline{v}^{(n)}(K)$ be defined by :

$$
\begin{equation*}
v_{i j}^{(l)}(K)=v_{r s}^{(l)} b_{i r} b_{j s} \quad(l=1, \ldots, n) \tag{5.11}
\end{equation*}
$$

We get, obviously, that $\underline{v}^{(1)}(K), \ldots, \underline{v}^{(n)}(K)$ is a basis in $V(\hat{V}, K)$; ordering now the elements $K$ in $\mathcal{C}_{h}$, and calling them, say, $K_{1}, \ldots, K_{N(h)}$, we can set, for each $j=1, \ldots, n$ and for each $l=1, \ldots, N(h)$ :

$$
\underline{v}^{(n(l-1)+j)}=\left\{\begin{array}{l}
\underline{v}^{(j)}\left(K_{l}\right) \text { on } K_{l}  \tag{5.12}\\
0 \text { on } K_{m} \text { for } m \neq l
\end{array}\right.
$$

we obtain a basis in $V_{h}=V\left(\hat{V}, \mathfrak{C}_{h}\right)$; it is easy to verify that, with a basis of this type, $\mathcal{A}$ becomes a block-diagonal matrix, each block being a square $n \times n$ matrix which is symmetrical and positive definite.

On the other hand, for the computation of the coefficients of $\mathcal{B}$ the use of Green's formula is necessary, since we have

$$
\begin{equation*}
b\left(\underline{v}^{(j)}, \varphi^{(i)}\right)=\sum_{K \in \mathscr{O}_{h}} \int_{K} v_{l m}^{(j)} \varphi_{l l m}^{(i)} \mathrm{d} x \tag{5.13}
\end{equation*}
$$

and if $W_{h}=W\left(r, s, \mathcal{C}_{h}\right)$ (for some value of $r \geqslant 3$ and $s \geqslant 1$ ) the value of the $\varphi^{(i)}$ at the interior of each element $K$ may not, in general, be known.
$n^{\circ}$ décembre 1975, R-3.

Therefore the coefficients $b\left(\underline{v}^{(1)}, \varphi^{(2)}\right)$ must be written in the following form :

$$
\begin{equation*}
b\left(\underline{( }^{(j)}, \varphi^{(\imath)}\right)=\sum_{K \in \mathscr{C}_{h}} \int_{\partial \hat{} K}\left(\varphi_{/ r}^{(l)} v_{r s}^{(j)} n_{s}-\varphi^{(i)} v_{r s / s}^{(j)} n_{r}\right) \mathrm{d} l \tag{5.14}
\end{equation*}
$$

where, for each $K,\left(n_{1}, n_{2}\right)$ is the direction of the outward normal derivative to $\partial K$. The computation of (5.14) is easy, since $v^{(j)}$ and $\varphi^{(i)}$ are known, with their first derivatives, at the interelement boundaries; moreover, if the basis $\underline{v}^{(1)}, \ldots, \underline{v}^{(N V)}$ is chosen of the form (5.12), for each $j$ there is just one term, in the sum which appears in (5.14), which differs from zero.
b) Computation of the «known vectors»

From the computational point of view, the evaluation of terms (5.6) and (5.7) is the greatest difficulty, since in general it is impossible to calculate a priori a particular solution $\underline{f}$ of the equation :

$$
\begin{equation*}
f_{\imath j / \imath \jmath}=p \quad \text { in each } \quad K \in \mathcal{C}_{h} . \tag{5.15}
\end{equation*}
$$

We shall show in the following a procedure which proved efficient in reducing evaluation of (5.6) and (5.7) to the computation of some integrals of the known function $p$, multiplied by suitable polynomials, performed on the elements $K$.

We shall first treat the simplest case, in which $\Omega$ is a square with sides parallel to the axes, and all the elements $K$ are also squares of the same type, with sides of lenght $h$. A more general case, in which all the elements $K$ are triangles of general type, will be treated later on.

Let then $K$ be a square element and let $\left(x_{1}^{\prime}, x_{2}^{\prime}\right)$ and $\left(x_{1}^{\prime \prime}, x_{2}^{\prime \prime}\right)$ be the coordinates of the lower left vertex and, respectively, of the upper right vertex of $K$, as shown in fig. 1.

In each $K$ we choose now $\underline{f}$ of the type :

$$
\begin{align*}
& f_{11}=f_{22}=0 \text { in } K  \tag{5.16}\\
& f_{12}=f_{21}=f=\frac{1}{2} \int_{x^{\prime} 1}^{x_{1}} \int_{x^{\prime}}^{x_{2}} p\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi \tag{5.17}
\end{align*}
$$

and we remark that, in this case, we get :

$$
\begin{array}{ll}
f_{i J / 2 j}=p \text { in } K, & \\
f\left(x_{1}^{\prime}, x_{2}\right)=0 & x_{2}^{\prime} \leqslant x_{2} \leqslant x_{2}^{\prime \prime}, \\
f\left(x_{1}, x_{2}^{\prime}\right)=0 & x_{1}^{\prime} \leqslant x_{1} \leqslant x_{1}^{\prime \prime} . \tag{5.20}
\end{array}
$$



Figure 1

We observe that if $\underline{v}^{(i)}$ is an element of a basis of type (5.12) which is not zero in $K$, then :

$$
\begin{equation*}
\left[\underline{f}, \underline{v}^{(t)}\right]=\int_{K} f_{r s} v_{r s}^{(t)} \mathrm{d} x=2 \int_{K} f v_{12}^{(i)} \mathrm{d} x \tag{5.21}
\end{equation*}
$$

Setting now, for the sake of simplicity,

$$
\begin{equation*}
v=v_{12}^{(i)} \tag{5.22}
\end{equation*}
$$

we define in $K$ the function

$$
g\left(x_{1}, x_{2}\right)=\int_{x_{1}}^{x^{\prime \prime}} \int_{x_{2}}^{x^{\prime \prime}} v\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi
$$

First we observe that, since $v$ is supposed to be a polynomial in $K, g$ is also a polynomial, and its expression can be computed without difficulty. Moreover we have :

$$
\begin{array}{ll}
g_{/ 12}=v \text { in } K, & \\
g\left(x_{1}, x_{2}^{\prime \prime}\right)=0 & x_{1}^{\prime} \leqslant x_{1} \leqslant x_{1}^{\prime \prime} \\
g\left(x_{1}^{\prime \prime}, x_{2}\right)=0 & x_{2}^{\prime} \leqslant x_{2} \leqslant x_{2}^{\prime \prime} \tag{5.26}
\end{array}
$$

Therefore the quantity (5.21), integrating by parts and using (5.18)-(5.20) and (5.24)-(5.26), becomes :

$$
\begin{align*}
\int_{K} f v \mathrm{~d} x=\int_{K} f g_{/ 12} \mathrm{~d} x= & \int_{K} f_{l 12} g \mathrm{~d} x-\int_{\partial K} f_{/ 2} g n_{1} \mathrm{~d} l+  \tag{5.27}\\
& +\int_{\partial K} f g_{/ 1} n_{2} \mathrm{~d} l=\int_{K} f_{/ 12} g \mathrm{~d} x=\frac{1}{2} \int_{K} p g \mathrm{~d} x
\end{align*}
$$

and finally

$$
\begin{equation*}
\left[\underline{f}, \underline{v}^{(i)}\right]=\int_{K} p g \mathrm{~d} x \tag{5.28}
\end{equation*}
$$

Since $p$ is the known function and $g$ is polynomial which can be easily deduced from $\underline{v}^{(i)}$, terms (5.6) can be easily computed using expression (5.28).

Let now deal with terms (5.7); by means of Green's formula and of (5.16), (5.17) we get

$$
\begin{equation*}
b\left(\underline{f}, \varphi^{(i)}\right)=2 \sum_{K \in \mathcal{G}_{h}} \int_{\partial K}\left(f \varphi_{/ 2}^{(i)} n_{1}-f_{/ 1} \varphi^{(i)} n_{2}\right) \mathrm{d} l . \tag{5.29}
\end{equation*}
$$

It can be easily verified that if, for instance, $W_{h}=W\left(3,1, \mathcal{C}_{h}\right)$ (this procedure can nevertheless be followed also for the general case), then we can construct a basis for $W_{h}$ in the following way: of all vertices of the elements $K$, we define as «nodes » those which are internal to $\Omega$, and we associate, with each node $P$, three functions $\varphi_{0}, \varphi_{1}, \varphi_{2}$ (of $W_{h}$ ) such that

$$
\begin{equation*}
\varphi_{r}=\frac{\partial \varphi_{r}}{\partial x_{1}}=\frac{\partial \varphi_{r}}{\partial x_{2}}=0 \quad(r=0,1,2) \tag{5.30}
\end{equation*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.
at each node different from $P$, and moreover :

$$
\begin{array}{ll}
\varphi_{0}(P)=1, & \frac{\partial \varphi_{0}}{\partial x_{1}}(P)=\frac{\partial \varphi_{0}}{\partial x_{2}}(P)=0 \\
\frac{\partial \varphi_{1}}{\partial x_{1}}(P)=1, & \varphi_{1}(P)=\frac{\partial \varphi_{1}}{\partial x_{2}}(P)=0 \\
\frac{\partial \varphi_{2}}{x_{2}}(P)=1, & \varphi_{2}(P)=\frac{\partial \varphi_{2}}{x_{1}}(P)=0 \tag{5.33}
\end{array}
$$

As $P$ describes all nodes, the set of the functions $\varphi_{0}, \varphi_{1}, \varphi_{2}$ verifying (5.30)-(5.33) describes a basis of $W\left(3,1, \mathcal{G}_{h}\right)$; we shall also say, in other words, that « we have chosen as degrees of freedom in $W\left(3,1, \mathfrak{C}_{h}\right)$ the value of $\varphi$ and of its first derivatives at each node ».

Let now $\varphi=\varphi^{(i)}$ be an element of such basis, and let $P \equiv\left(P_{1}, P_{2}\right)$ be the node of $\mathscr{G}_{h}$ corresponding to $\varphi$, in the sense that $P$ is the unique node in which $\varphi=\varphi_{/ 1}=\varphi_{/ 2}=0$ do not hold. The sum which appears in (5.29) is now reduced to the sum of the four integrals over the four sides which have $P$ as a vertex (see fig. 2 , where the indicated sides are called $L_{12}, L_{24}, L_{34}, L_{13}$ ).

If we denote by $f^{(l)}$ the restriction to $K_{l}$ of $f$, we get


Figure 2

$$
\begin{align*}
b(\underline{f}, \varphi)= & \sum_{K \in \mathcal{O}_{h}} 2 \int_{\partial K}\left(f \varphi_{/ 2} n_{1}-f_{/ 1} \varphi n_{2}\right) \mathrm{d} l=  \tag{5.34}\\
= & 2 \int_{L_{12}}\left(f^{(1)}-f^{(2)}\right) \varphi_{/ 2} \mathrm{~d} l+2 \int_{L_{24}}\left(f_{11}^{(4)}-f_{11}^{(2)}\right) \varphi \mathrm{d} l \\
& \quad+2 \int_{L_{34}}\left(f^{(3)}-f^{(4)}\right) \varphi_{/ 2} \mathrm{~d} l+2 \int_{L_{13}}\left(f_{/ 1}^{(3)}-f_{/ 1}^{(1)}\right) \varphi \mathrm{d} l
\end{align*}
$$

and from (5.19), (5.20) :

$$
\begin{align*}
& b(\underline{f}, \varphi)=2 \int_{L_{12}} f^{(1)} \varphi_{/ 2} \mathrm{~d} l-2 \int_{L_{24}} f_{11}^{(2)} \varphi \mathrm{d} l  \tag{5.35}\\
&+2^{\circ} \int_{L_{34}} f^{(3)} \varphi_{/ 2} \mathrm{~d} l-2 \int_{L_{13}} f_{11}^{(1)} \varphi \mathrm{d} l .
\end{align*}
$$

We shall only show how to calculate the integral :

$$
\begin{equation*}
\int_{L_{12}} f^{(1)} \varphi_{/ 2} \mathrm{~d} l \tag{5.36}
\end{equation*}
$$

the integrals over the other sides can be calculated in a similar manner. Integrating (5.36) by parts we get :

$$
\begin{align*}
\int_{L_{12}} f^{(1)} \varphi_{/ 2} \mathrm{~d} l & =f^{(1)}(P) \varphi(P)-\int_{L_{12}} f_{12}^{(1)} \varphi\left(P_{1}, x_{2}\right) \mathrm{d} x  \tag{5.37}\\
& =\varphi(P) \frac{1}{2} \int_{K_{1}} p\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi-\frac{1}{2} \int_{K_{1}} p\left(\xi_{1}, \xi_{2}\right) \varphi\left(P_{1}, \xi_{2}\right) \mathrm{d} \xi \\
& =\frac{1}{2} \int_{K_{1}} p\left(\xi_{1}, \xi_{2}\right)\left[\varphi(P)-\varphi\left(P_{1}, \xi_{2}\right)\right] \mathrm{d} \xi .
\end{align*}
$$

If we consider now the contribution of the other three integrals appearing in (5.35) we obtain :

$$
\begin{align*}
b(\underline{f}, \varphi) & =\int_{K_{1}} p\left(x_{1}, x_{2}\right)\left[\varphi(P)-\varphi\left(P_{1}, x_{2}\right)-\varphi\left(x_{1}, P_{2}\right)\right] \mathrm{d} x_{1} \mathrm{~d} x_{2}  \tag{5.38}\\
& -\int_{K_{2}} p\left(x_{1}, x_{2}\right) \varphi\left(x_{1}, P_{2}\right) \mathrm{d} x_{1} \mathrm{~d} x_{2}-\int_{K_{3}} p\left(x_{1}, x_{2}\right) \varphi\left(P_{1}, x_{2}\right) \mathrm{d} x_{1} \mathrm{~d} x_{2}
\end{align*}
$$

We remark that, since $\varphi$ is a polynomial on each $L_{i j}$, functions $\varphi\left(P_{1}, x_{2}\right)$ and $\varphi\left(x_{1}, P_{2}\right)$ will also be polynomials in each $K_{l}$, and therefore (5.7) can be easily calculated by (5.38).

We shall examine now the more general case in which the elements $K$ are triangles; we shall show that, even in this case, quantities (5.6) and (5.7) can be calculated by evaluating integrals over some elements $K$ of the function $p\left(x_{1}, x_{2}\right)$ multiplied by suitable polynomials. Let $K$ a triangle in $\mathfrak{C}_{h}$; we first remark that, by translation and eventually exchanging the $x_{1}$-axis with the $x_{2}$-axis, we can always reduce to the situation of fig. 3 . We denote by $P \equiv(O, O), P^{\prime} \equiv\left(x_{1}^{\prime}, x_{2}^{\prime}\right), P^{\prime \prime} \equiv\left(x_{1}^{\prime \prime}, x_{2}^{\prime \prime}\right)$ the vertices of $K$ and by $L_{1}=P P^{\prime}$ $L_{2}=P P^{\prime \prime}, L_{3}=P^{\prime} P^{\prime \prime}$ the sides of $K$, whose equations are supposed to be, respectively:

$$
\begin{gather*}
x_{2}=r_{1}\left(x_{1}\right)=\alpha_{1} x_{1}, \quad x_{2}=r_{2}\left(x_{1}\right)=\alpha_{2}\left(x_{1}\right)  \tag{5.39}\\
x_{3}=r_{3}\left(x_{1}\right)=\alpha_{3} x_{1}+\beta
\end{gather*}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.


Figure 3
Let us study, first, the quantity (5.6); we suppose that $\underline{v}=\underline{v}^{(i)}$ is an element of the basis (5.12) which is different from zero on $K$, and we want to calculate the quantity

$$
\begin{equation*}
\int_{K} f_{r s} v_{r s} \mathrm{~d} x \tag{5.40}
\end{equation*}
$$

If we choose again $f$ of the type

$$
\begin{equation*}
f_{11}=f_{22}=0, \quad f_{12}=f_{21}=f \tag{5.41}
\end{equation*}
$$

and we set, for the sake of simplicity,

$$
v=v_{12}=v_{21},
$$

then (5.40) can be written as

$$
\begin{equation*}
2 \int_{K} f v \mathrm{~d} x \tag{5.42}
\end{equation*}
$$

In order to evaluate (5.42), we calculate first a function $g\left(x_{1}, x_{2}\right)$ which satisfies the following conditions:

$$
\begin{align*}
g_{/ 12} & =\frac{\partial^{2} g}{\partial x_{1} \partial x_{2}}=v \text { in } K  \tag{5.43}\\
g & =g_{/ 2}=0 \text { on } L_{3} \tag{5.44}
\end{align*}
$$

Such a function can be calculated explicitly by setting

$$
\begin{equation*}
g\left(x_{1}, x_{2}\right)=g^{1}\left(x_{1}, x_{2}\right)+g^{2}\left(x_{2}\right)+g^{3}\left(x_{1}\right) \tag{5.45}
\end{equation*}
$$

where the functions $g^{i}$ are such that :

$$
\begin{equation*}
g_{/ 12}^{1}=\frac{\partial^{2} g^{1}}{\partial x_{1} \partial x_{2}}=v \text { in } K \tag{5.46}
\end{equation*}
$$

$$
\begin{align*}
g_{12}^{2}\left(x_{2}\right) & =-g_{12}^{1}\left(r_{3}^{-1}\left(x_{2}\right), x_{2}\right)  \tag{5.47}\\
g^{3}\left(x_{1}\right) & =-g^{1}\left(x_{1}, r_{3}\left(x_{1}\right)\right)-g^{2}\left(r_{3}\left(x_{1}\right)\right) \tag{5.48}
\end{align*}
$$

It can be easily verified that the function defined by (5.45)-(5.48) verifies (5.43), (5.44); we also remark that, since $v$ is a polynomial, the functions $g^{i}$ can be easily calculated and are also polynomials.

By (5.43) and by Green's formula we get

$$
\begin{equation*}
\int_{K} f v \mathrm{~d} x=\int_{K} f g_{/ 12} \mathrm{~d} x=\int_{K} f_{l 12} g \mathrm{~d} x+\int_{\partial K} f g_{/ 2} \mathrm{~d} l-\int_{\partial K} f_{l 1} g \mathrm{~d} l \tag{5.49}
\end{equation*}
$$

We choose $f$ in such a way that :

$$
\begin{align*}
f_{12} & =\frac{1}{2} p \text { in } K  \tag{5.50}\\
f & =0 \text { on } L_{1}  \tag{5.51}\\
f & =0 \text { on } L_{2} \tag{5.52}
\end{align*}
$$

An explicit solution of (5.50)-(5.52) can be obtained in the following way. Let $\tilde{p}\left(x_{1}, x_{2}\right)$ be defined by

$$
\tilde{p}\left(x_{1}, x_{2}\right)= \begin{cases}p\left(x_{1}, x_{2}\right) & \text { if }\left(x_{1}, x_{2}\right) \in K  \tag{5.53}\\ 0 & \text { if }\left(x_{1}, x_{2}\right) \notin K\end{cases}
$$

and let $z^{1}$ be defined by,

$$
\begin{equation*}
z^{1}\left(x_{1}, x_{2}\right)=\frac{1}{2} \int_{0}^{x_{1}} \int_{x_{2}^{\prime \prime}}^{x_{2}} \tilde{p}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi \tag{5.54}
\end{equation*}
$$

of course $z^{1}\left(x_{1}, x_{2}\right)$ satisfies the following conditions :

$$
\begin{gather*}
z_{/ 12}^{1}=\frac{1}{2} \tilde{p}\left(x_{1}, x_{2}\right)  \tag{5.55}\\
z^{1}\left(0, x_{2}\right)=0 ; \quad z^{1}\left(x_{1}, x_{2}^{\prime \prime}\right)=0 \tag{5.56}
\end{gather*}
$$

Let now $\chi\left(x_{2}\right)$ be defined by

$$
\begin{equation*}
\chi\left(x_{2}\right)=z^{1}\left(r_{1}^{-1}\left(x_{2}\right), x_{2}\right) \tag{5.57}
\end{equation*}
$$

we set

$$
z^{2}\left(x_{1}, x_{2}\right)= \begin{cases}\chi\left(x_{2}\right) & \text { for } x_{2}>0  \tag{5.58}\\ 0 & \text { for } x_{2} \leq 0\end{cases}
$$

$\mathrm{n}^{\circ}$ décembre 1975, R-3.
and we have that function $f$, given by

$$
\begin{equation*}
f\left(x_{1}, x_{2}\right)=z^{1}\left(x_{1}, x_{2}\right)-z^{2}\left(x_{1}, x_{2}\right) \tag{5.59}
\end{equation*}
$$

i.e.
(5.60) $f\left(x_{1}, x_{2}\right)= \begin{cases}\frac{1}{2} \int_{0}^{x_{1}} \int_{x_{2}^{\prime \prime}}^{x_{2}} \tilde{p}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi-\frac{1}{2} \int_{x_{2}^{\prime \prime}}^{x_{2}} \int_{0}^{r_{1}^{\prime 1}\left(x_{2}\right)} \tilde{p}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi \\ \frac{1}{2} \int_{0}^{x_{1}} \int_{x_{2}^{\prime \prime}}^{x_{2}} \tilde{p}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi \text { for } x_{2} \leqslant 0 & \text { for } x_{2}>0\end{cases}$
satisfies conditions (5.50)-(5.52). We have, in particular, the following expressions for $f_{/ 1}$ and $f_{12}$ :

$$
\begin{gather*}
f_{/ 1}\left(x_{1}, x_{2}\right)=\frac{1}{2} \int_{x_{2}^{\prime \prime}}^{x_{2}} \tilde{p}\left(x_{1}, \xi_{2}\right) \mathrm{d} \xi_{2}  \tag{5.61}\\
f_{l 2}\left(x_{1}, x_{2}\right)=\left\{\begin{array}{l}
\frac{1}{2}\left[\int_{0}^{x_{1}} \tilde{p}\left(\xi_{1}, x_{2}\right) \mathrm{d} \xi_{1}-\int_{0}^{r_{1}-1}\left(x_{2}\right)\right. \\
p \\
\left(\xi_{1}, x_{2}\right) \mathrm{d} \xi_{1} \\
\left.-\frac{1}{\alpha_{1}} \int_{x^{\prime \prime} 2}^{x_{2}} \tilde{p}\left(r_{1}^{-1}\left(x_{2}\right), \xi_{2}\right) \mathrm{d} \xi_{2}\right] \text { for } x_{2}>0 \\
\frac{1}{2} \int_{0}^{x_{1}} \tilde{p}\left(\xi_{1}, x_{2}\right) \mathrm{d} \xi_{1} \text { for } x_{2} \leqslant 0
\end{array}\right.
\end{gather*}
$$

By formulas (5.44), (5.51), (5.52), we have first that (5.49) becomes

$$
\begin{equation*}
\int_{K} f v \mathrm{~d} x=\int_{K} f_{/ 12} g \mathrm{~d} x-\int_{L_{1}} f_{1} g n_{2}^{(1)} \mathrm{d} l-\int_{L_{2}} f_{/ 1} g n_{2}^{(2)} \mathrm{d} l \tag{5.63}
\end{equation*}
$$

where $\left(n_{1}^{(2)}, n_{2}^{(2)}\right)$ is the normal outward direction to the side $L_{1}$. Substituting (5.50) and (5.61) in (5.63) we get :

$$
\begin{equation*}
\int_{K} f v \mathrm{~d} x=\frac{1}{2} \int_{K} p g \mathrm{~d} x-\frac{c_{1}}{2} \int_{0}^{x_{1}^{\prime}} \int_{x_{2}^{\prime \prime}}^{r_{1}\left(x_{1}\right)} \tilde{p} G^{(1)} \mathrm{d} x-\frac{c_{2}}{2} \int_{0}^{x_{1}^{\prime \prime}} \int_{x_{2}^{\prime \prime}}^{r_{2}\left(x_{1}\right)} \tilde{p} G^{(2)} \mathrm{d} x \tag{5.64}
\end{equation*}
$$

where $G^{(1)}\left(x_{1}, x_{2}\right)=g\left(x_{1}, r_{1}\left(x_{1}\right)\right), G^{(2)}\left(x_{1}, x_{2}\right)=g\left(x_{1}, r_{2}\left(x_{1}\right)\right)$ and $c_{1}, c_{2}$ are constants depending on $\alpha_{1}, \alpha_{2}$. We remark now that $\tilde{p}\left(x_{1}, x_{2}\right)$ is zero for $x_{2}^{\prime \prime} \leqslant x_{2} \leqslant r_{2}\left(x_{1}\right)$; Setting now

$$
G\left(x_{1}, x_{2}\right)= \begin{cases}g\left(x_{1}, x_{2}\right)-c_{1} G^{(1)}\left(x_{1}, x_{2}\right) & \text { if } 0 \leqslant x_{1} \leqslant x_{1}^{\prime}  \tag{5.65}\\ g\left(x_{1}, x_{2}\right) & \text { if } x_{1}^{\prime} \leqslant x_{1}\end{cases}
$$

we finally get

$$
\begin{equation*}
\int_{K} f v \mathrm{~d} x=\frac{1}{2} \int_{K} p G \mathrm{~d} x \tag{5.66}
\end{equation*}
$$

where $G$ is a polynomial in each of the two regions $x_{1} \leqslant x_{1}^{\prime}$ and $x_{1}^{\prime} \leqslant x_{1}$.
Let us now consider the terms of the type :

$$
\begin{equation*}
b(f, \varphi)=\sum_{K \in \mathcal{C}_{h}} \int_{K}\left(f_{\imath j} \varphi_{/ t \jmath}-p \varphi\right) \mathrm{d} x \tag{5.67}
\end{equation*}
$$

where $\varphi$ is an element of the basis for $W_{h}$. It is easy to verify that if, for instance, $W_{h}=W\left(3,1, \mathfrak{G}_{h}\right)$ and we choose as degrees of freedom the values of the functions and of their first derivatives, then the sum which appears in (5.67) reduces to a few terms. Let us consider just one of them, say a triangle $K$ as in figure 3. By Green's formula and due to the given choice for $f$, we have :

$$
\begin{equation*}
\int_{K}\left(f_{i j} \varphi_{l i j}-p \varphi\right) \mathrm{d} x=\int_{\partial K}\left(f \varphi_{/ 2} n_{1}-f_{l 1} \varphi n_{2}\right) \mathrm{d} l ; \tag{5.68}
\end{equation*}
$$

let us first consider the term :

$$
\begin{equation*}
\int_{\partial K} f_{/ 1} \varphi n_{2} \mathrm{~d} l \tag{5.69}
\end{equation*}
$$

We have

$$
\begin{align*}
\int_{\partial K} f_{l 1} \varphi n_{2} \mathrm{~d} l= & \sum_{i=1}^{3} \int_{L_{i}} f_{/ 1} \varphi n_{2}^{(1)} \mathrm{d} l  \tag{5.70}\\
= & n_{2}^{(1)} c_{1} \int_{0}^{x_{1}^{\prime}} f_{/ 1}\left(\xi_{1}, r_{1}\left(\xi_{1}\right)\right) \varphi\left(\xi_{1}, r_{1}\left(\xi_{1}\right)\right) \mathrm{d} \xi_{1} \\
& +n_{2}^{(2)} c_{2} \int_{0}^{x_{1}^{\prime \prime}} f_{l 1}\left(\xi_{1}, r_{2}\left(\xi_{1}\right)\right) \varphi\left(\xi_{1}, r_{2}\left(\xi_{1}\right)\right) \mathrm{d} \xi_{1} \\
& +n_{2}^{(3)} c_{3} \int_{x_{1}^{\prime}}^{x_{1}^{\prime \prime}} f_{l 1}\left(\xi_{1}, r_{3}\left(\xi_{1}\right)\right) \varphi\left(\xi_{1}, r_{3}\left(\xi_{1}\right)\right) \mathrm{d} \xi_{1}
\end{align*}
$$

and from (5.61)

$$
\begin{align*}
\int_{\partial K} f_{11} \varphi n_{2}= & \frac{n_{2}^{(1)} c_{1}}{2} \int_{0}^{x_{1}^{\prime}} \int_{r_{2}\left(x_{1}\right)}^{r_{1}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right) \Phi_{1}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi  \tag{5.71}\\
& +\frac{n_{2}^{(3)} c_{3}}{2} \int_{x_{1}^{\prime}}^{x_{1}^{\prime \prime}} \int_{r_{2}\left(x_{1}\right)}^{r_{3}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right) \Phi_{3}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi .
\end{align*}
$$

$n^{\circ}$ décembre 1975, R-3.

Since $\Phi_{i}\left(x_{1}, x_{2}\right)=\varphi\left(x_{1}, r_{i}\left(x_{1}\right)\right)(i=1,3)$ are polynomials, formula (5.71) is of the expected type. We shall deal now with the term

$$
\begin{equation*}
\int_{\partial \mathrm{K}} f \varphi_{/ 2} n_{1} \mathrm{~d} l \tag{5.72}
\end{equation*}
$$

which reduces, by (5.51), (5.52), to

$$
\begin{equation*}
\int_{L_{3}} f \varphi_{/ 2} n_{1}^{(3)} \mathrm{d} l . \tag{5.73}
\end{equation*}
$$

Let now $Q\left(x_{1}\right)$ be a primitive of the polynomial $\varphi_{12}\left(x_{1}, r_{3}\left(x_{1}\right)\right)$, which is itself a polynomial and easily computable; changing the variable and integrating by parts we get :

$$
\begin{align*}
n_{1}^{(3)} \int_{L_{3}} f \varphi_{/ 2} \mathrm{~d} l= & n_{1}^{(3)} c_{3} \int_{x_{1}^{\prime}}^{x_{1}^{\prime \prime}} f\left(\xi_{1}, r_{3}\left(\xi_{1}\right)\right) \varphi_{/ 2}\left(\xi_{1}, r_{3}\left(\xi_{1}\right)\right) \mathrm{d} \xi_{1}  \tag{5.74}\\
= & n_{1}^{(3)} c_{3} \int_{x_{1}^{\prime}}^{x_{1}^{\prime \prime}} f_{/ 1}\left(\xi_{1}, r_{3}\left(\xi_{1}\right)\right) Q\left(\xi_{1}\right) \mathrm{d} \xi_{1} \\
& -\alpha_{3} n_{1}^{(3)} c_{3} \int_{x_{1}^{\prime}}^{x_{1}^{\prime \prime}} f_{/ 2}\left(\xi_{1}, r_{3}\left(\xi_{1}\right)\right) Q\left(\xi_{1}\right) \mathrm{d} \xi_{1}
\end{align*}
$$

From (5.61), (5.62) we obtain :

$$
\begin{align*}
\int_{\partial K} f \varphi_{/ 2} n_{1} \mathrm{~d} l= & -\frac{n_{1}^{(3)} c_{3}}{2} \int_{x_{1}^{\prime}}^{x_{1}^{\prime}} \int_{r_{2}\left(x_{1}\right)}^{r_{3}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right) Q\left(\xi_{1}\right) \mathrm{d} \xi  \tag{5.75}\\
& +\frac{n_{1}^{(3)} c_{3}}{2} \int_{K} p\left(\xi_{1}, \xi_{2}\right) Q\left(\xi_{1}\right) \mathrm{d} \xi \\
& +\frac{n_{1}^{(3)} c_{3} \alpha_{3}}{2 \alpha_{1}} \int_{x_{1}^{\prime}}^{\bar{x}_{1}} \int_{r_{2}\left(x_{1}\right)}^{r_{3}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right) Q\left(\xi_{1}\right) \mathrm{d} \xi
\end{align*}
$$

and finally :

$$
\begin{align*}
& \int_{K}\left(f_{i j} \varphi_{i i j}-p \varphi\right) \mathrm{d} x  \tag{5.76}\\
& =\frac{1}{2}\left\{\left(n_{1}^{(3)} c_{3}-n_{2}^{(1)} c_{1}\right) \int_{0}^{x_{1}^{\prime}} \int_{r_{2}\left(x_{1}\right)}^{r_{1}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right)\left(Q\left(\xi_{1}\right)+\Phi_{1}\left(\xi_{1}, \xi_{2}\right)\right) \mathrm{d} \xi\right. \\
& +c_{3}\left(\frac{n_{1}^{(3)} \alpha_{3}}{\alpha_{1}}-n_{2}^{(3)}\right) \int_{x_{1}^{\prime}}^{\bar{x}_{1}} \int_{r_{2}\left(x_{1}\right)}^{r_{3}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right)\left(Q\left(\xi_{1}\right)+\Phi_{3}\left(\xi_{1}, \xi_{2}\right)\right) \mathrm{d} \xi \\
& \left.\quad-n_{2}^{(3)} c_{3} \int_{\bar{x}_{1}}^{x_{1}^{\prime \prime}} \int_{r_{2}\left(x_{1}\right)}^{r_{3}\left(x_{1}\right)} p\left(\xi_{1}, \xi_{2}\right) \Phi_{3}\left(\xi_{1}, \xi_{2}\right) \mathrm{d} \xi\right\}
\end{align*}
$$

which is again of the desired type.

## c) Solution of the linear system

We have seen in a) that, if the chosen basis for $V_{h}$ is of type (5.12), then $\mathcal{A}$ is a block diagonal matrix, each block being an $n \times n$ square matrix, symmetrical and positive definite. Therefore the best manner for solving the system (5.8) is an a priori inversion of . $\mathcal{t}$ (cfr. e.g. [32]), which is easily performed and leads, by the substitution

$$
\begin{equation*}
\vec{U}=\mathcal{A}^{-1} \mathfrak{B}^{T} \vec{\Psi}+\mathcal{A}^{-1} \vec{F} \tag{5.77}
\end{equation*}
$$

to the linear system

$$
\begin{equation*}
\mathfrak{B A}^{-1} \mathfrak{B}^{T} \vec{\Psi}=\mathfrak{B A}^{-1} \vec{F}+\vec{G} \tag{5.78}
\end{equation*}
$$

in the $N W$ unknowns $\Psi_{1}, \ldots, \Psi_{N W}$. Let us set:

$$
\begin{align*}
& \mathfrak{H}=\mathscr{B} \mathcal{A}^{-1} \mathcal{B}^{T}  \tag{5.79}\\
& \vec{P}=\mathscr{B} \mathcal{A}^{-1} \vec{F}+\vec{G} . \tag{5.80}
\end{align*}
$$

Proposition 5.1. - Hypothesis H 1 implies that $\mathfrak{H}$ is symmetrical and positive definite.

Proof. - From (5.79) we have obviously that $\mathscr{H}$ is symmetrical, since $\mathcal{A}^{-1}$ is symmetrical. Moreover, since $\mathcal{A}^{-1}$ is positive definite, we have for every vector $\vec{\Phi}=\Phi_{1}, \ldots, \Phi_{N W}$ that :

$$
\begin{equation*}
\left(\mathcal{B A}^{-1} \mathfrak{B}^{T} \vec{\Phi}, \vec{\Phi}\right)=\left(\mathcal{A}^{-1} \mathcal{B}^{T} \vec{\Phi}, \mathfrak{B}^{T} \vec{\Phi}\right) \geqslant 0 \tag{5.81}
\end{equation*}
$$

and moreover,

$$
\begin{equation*}
\left(\mathcal{B} \mathcal{A}^{-1} \mathcal{B}^{T} \vec{\Phi}, \vec{\Phi}\right)=0 \Leftrightarrow \mathcal{B}^{T} \vec{\Phi}=0 . \tag{5.82}
\end{equation*}
$$

Let now $\varphi$ be the element of $W_{h}$ defined by

$$
\begin{equation*}
\varphi=\sum_{i=1}^{N W} \varphi^{(i)} \Phi_{i} ; \tag{5.83}
\end{equation*}
$$

if $\mathcal{B}^{T} \vec{\Phi}=0$ then, for all $\underline{v}$ in $V_{h}$, it follows

$$
\begin{equation*}
b(\underline{v}, \varphi)=\sum_{j=1}^{N V} V_{j} b\left(\underline{\dot{v}}^{(j)}, \varphi\right)=\sum_{j=1}^{N V} V_{j}\left(\mathcal{B}^{T} \vec{\Phi}\right)_{j}=0 \tag{5.84}
\end{equation*}
$$

and, from hypothesis H1, (5.84) implies $\varphi=0$ and therefore $\vec{\Phi}=0$. So $\Re^{T} \vec{\Phi}=0$ iff $\vec{\Phi}=0$ and from (5.81), (5.82) we obtain the result.

The numerical solution of

$$
\begin{equation*}
\mathscr{H} \vec{\Psi}=\vec{P} \tag{5.85}
\end{equation*}
$$

can therefore be performed in one of the classical ways for solving linear systems with a symmetrical positive definite matrix.
$\mathrm{n}^{\circ}$ décembre 1975, R-3.

We shall now make some comments on the «topological structure» of $\mathscr{H}$ (i.e. the position of the non-zero elements). It is easy to verify that, if we choose as degrees of fredom in $W_{h}=W\left(r, s, \mathscr{C}_{h}\right)$ the values of $\varphi$ and of some of its partial derivatives at certain points on the sides of the elements $K$, then the coefficient $H_{i j}$ of $\mathcal{H}$ is different from zero only if $\varphi^{(i)}$ and $\varphi^{(j)}$ have its «corresponding points » on the boundary of the same element.


Figure 4

It follows that the topological structure of $\mathscr{H}$ is depending only on the choice of the degrees of freedom in $W_{h}$. We remark that, however, in order to have a $C^{1}$-continuity for the functions of $W\left(r, s, \mathfrak{C}_{h}\right)$, we are compelled to choose, as degrees of freedom, at least the values of $\left(\varphi, \varphi_{/ 1}, \varphi_{/ 2}\right)$ at the nodes,


Figure 5
and this justifies conditions $r \geqslant 3, s \geqslant 1$. We also remark that in many cases the degrees of freedom of $W\left(r, s, \mathfrak{C}_{h}\right)$ will coincide with those used in the classical, conforming and non conforming, finite element «displacement» methods for biharmonic equations. Nevertheless the coincidence of the degrees of freedom does not imply that the trial functions will coincide at the interior of the elements and, in many cases, even on the interelement boundaries, as in the non-conforming methods. Similarly, the structure of the matrices for

[^1]the displacement methods and for the hybrid methods will coincide, but not the values of non-zero coefficients; again the «known vectors» are computed in a similar fashion but they do not coincide. We also remark that, even if the topological structure of $\mathscr{H}$ remains unchanged, changing $V_{h}$, the values of the non-zero coefficients (and, of course, of the known vector) depend on the choice of $V_{h}$; we have therefore many «hybrid analogues » of the same displacement method.

We give in the following some examples of choices of degrees of freedom in $W\left(r, s, \mathcal{G}_{h}\right)$ for some value of $r$ and $s$ in the case $\hat{K}=$ triangle (fig. 4) and in the case $\hat{K}=$ unit square (fig. 5). The reader will recognize many of the classical «structures» used in the displacement methods (as e.g. Zienkiewicz, Clough-Tocher, Bramble-Zlamal, Adini, Fraeijs de Veubeke-Sander, Bogner-Fox-Schmit) and also other different «structures» that can be easily used by the «hybrid analogue» approach.

## d) Numerical results

We shall report here some of the results obtained in the numerical experiments performed on the Honeywell 6030 of the «Centro di Calcoli Numerici dell'Università di Pavia », in the simplest case of a square plate; the «reference element $\hat{K} »$ was also a square and the choices $W\left(3,1, \mathcal{G}_{h}\right)$ and $W\left(3,3, \mathcal{G}_{h}\right)$ has been tested for $W_{h}$ (that is, the hybrid analogous of the "Adini element» and of the «Bogner-Fox-Schmit element»). For the case $W_{h}=W\left(3,1, \mathcal{C}_{h}\right)$, different choices for $\hat{V}$ have also been tested, that is :

$$
\begin{aligned}
& \hat{V} 1 \equiv\left\{\begin{array}{l}
v_{11}=a_{0}+a_{1} x+a_{2} y \\
v_{12}=v_{21}=b_{0}+b_{1} x+b_{2} y \\
v_{22}=c_{0}+c_{1} x+c_{2} y
\end{array}\right. \\
& \hat{V} 2 \equiv\left\{\begin{array}{l}
v_{11}=a_{0}+a_{1} x+a_{2} x y \\
v_{12}=v_{21}=b_{0}+b_{1} x+b_{2} y \\
v_{22}=c_{0}+c_{1} y+c_{2} x y
\end{array}\right. \\
& \hat{V} 3 \equiv\left\{\begin{array}{l}
v_{11}=a_{0}+a_{1} x+a_{2} y+a_{3} x y \\
v_{12}=v_{21}=b_{0}+b_{1} x+b_{2} y \\
v_{22}=c_{0}+c_{1} x+c_{2} y+c_{3} x y
\end{array}\right. \\
& \hat{V} 4 \equiv\left\{\begin{array}{l}
v_{11}=a_{0}+a_{1} x+a_{2} y+a_{3} x y \\
v_{12}=v_{21}=b_{0}+b_{1} x+b_{2} y+b_{3} x^{2}+b_{4} y^{2} \\
v_{22}=c_{0}+c_{1} x+c_{2} y+c_{3} x y
\end{array}\right.
\end{aligned}
$$

(We note that, in the case $\hat{V} 1$, hypothesis $C(\hat{V}, \hat{W})$ is not satisfied; the results, however, are «good», at least for the tested values of $h$ ). In the case $W_{h}=W\left(3,3, \mathfrak{G}_{h}\right)$ the only choice

$$
\hat{V} 5 \equiv\left\{\begin{array}{l}
v_{11}=a_{0}+a_{1} x+a_{2} y+a_{3} x y+a_{4} y^{2}, \\
v_{12}=v_{21}=b_{0}+b_{1} x+b_{2} y+b_{3} x^{2}+b_{4} y^{2} \\
v_{22}=c_{0}+c_{1} x+c_{2} y+c_{3} x y+c_{4} x^{2},
\end{array}\right.
$$

has been tested. Figures $6,7,8$ report the plots of the errors

$$
\begin{aligned}
& \operatorname{ERR}^{(1)} \simeq\left\|w-\psi_{h}\right\|_{H^{1}}\|w\|_{H^{1}}^{-1} \\
& \operatorname{ERR}^{(0)} \simeq\left\|w-\psi_{h}\right\|_{L^{2}}\|w\|_{L^{2}}^{-1} \\
& \operatorname{ERR}^{(2)} \simeq \frac{\left|u-u_{h}\right|}{|u|} \text { in the middle point of the plate, }
\end{aligned}
$$

for the different choices of $W_{h}$ and $V_{h}$, with different values of $h$ (or, more precisely, of $N=\frac{1}{h}=$ number of elements on each side). Further results, with an accurate comparison between some of the classical displacement methods and their hybrid analogues will be published in a forthcoming paper by the same authors.

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In figures 6-7-8 the following symbols are used :


Figure 6
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Figure 7
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Figure 8

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[^0]:    (1) For the defintions of the spaces $H^{m}(\Omega)$ and $H_{0}^{m}(\Omega)$, see eg [28]

[^1]:    $\mathrm{n}^{\circ}$ décembre 1975, R-3.

