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MIXED AND HYBRID FINITE ELEMENT METHODS

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MIXED AND HYBRID METHODS*

METHODES MIXTES ET HYBRIDES

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Abstract

This work, essentially expository in nature, is a survey of mixed and hybrid finite element methods and is meant to serve as a basic reference on the subject. Detailed descriptions of these methods including convergence results and error estimates are given for a second order elliptic problem in \mathbb{R}^n , for n=2,3. References for extensions of these methods to more complicated problems are given.

Résumé

Ce travail est un exposé des méthodes d'éléments finis mixtes et hybrides qui devrait servir de référence sur ce sujet. Des descriptions détaillées de ces méthodes, incluant des résultats de convergence et des estimations des erreurs, sont données pour un problème elliptique du deuxième ordre dans \mathbb{R}^n , pour n=2,3. On donne des références pour des extensions de ces méthodes à des problèmes plus compliqués.

Key words: mixed finite element methods, hybrid finite element methods

Mots-clés: méthodes d'éléments finis mixtes, méthodes d'éléments finis hybrides



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CHAPTER I

INTRODUCTION

1 - PRELIMINARY REMARKS

The terms "mixed finite element method", "hybrid finite element method", and even "mixed - hybrid finite element method" come essentially from the vocabulary of structural mechanics as do the terms "conforming finite element method" and "equilibrium finite element method" ; cf. in particular the historical account of hybrid and mixed elements given by PIAN (1983). Consider for simplicity the linear elasticity problem describing the displacement u of an elastic body $\Omega.$ A finite element method based on the variational formulation associated with the principle of minimization of the potential energy form is called a displacement method as the procedure yields an approximation of the displacement field u. However, in structural mechanics one is often more interested in the stress tensor $rac{\sigma}{\widetilde{\omega}}$ than in the displacement. Methods have been devised for computing an approximation of $ilde{\sigma}$ from the approximation of u. Alternatively the stress may be characterized as the minimum of the complementary energy form on the space of symmetric tensors satisfying a relation expressing the equilibrium between the internal and boundary forces acting on Ω , and finite element methods based on the associated variational form are termed equilibrium methods. The drawback of these methods is that it is not usually easy to construct an approximation space of tensors satisfying the equilibrium relation and having the required amount of regularity, c^0 in this instance. Mixed methods and hybrid methods were devised to avoid this difficulty.

With the aid of Lagrangian multipliers the above contrained minimization problem may be turned into a saddle point problem. The Hellinger-Reissner principle characterizes the pair (g,u) as the saddle point of a Lagrangian where g is allowed to vary over the space of symmetric tensors subject only to regularity contraints, the equilibrium relation now being incorporated in the Lagrangian. Finite element methods based on the corresponding variational principle are called mixed methods. Mixed methods were first introduced in the literature in HERMANN (1967) for the plate problem. We shall say, more generally, that a finite element method is a mixed method if it involves the simultaneous approximation of two or more vector fields defined on the physical domain. These will in general be the principal unknown and an expression involving one or more of its derivatives.

An alternative approach is taken in the development of hybrid methods. The idea is to retain the equilibrium requirement for functions in the space in which the stress is approximated but to relax the regularity requirement. Given a finite element mesh on the domain Ω , a mesh dependent Lagrangian is constructed for which the unique saddle point is $(\underline{\sigma},\phi)$ where $\underline{\sigma}$ varies over a space of tensors satisfying the equilibrium condition but defined only on the union of the interiors of the elements. The Lagrangian multiplier ϕ is a vector field defined on the union of the boundaries of the elements which turns out to be none other than the trace of the displacement u. Finite element methods based on the associated variational formulation are hybrid methods. It is equally possible, starting from the original displacement formulation, to define a Lagrangian, again mesh dependent, by introducing a multiplier λ to relax the regularity requirement on the space in which the displacement u is sought. The space of multipliers is again a space of vector fields on the union of the boundaries of the elements, and u is looked for in a space of vector fields defined on the interiors of the elements. There is a unique saddle point (u,λ) for this Lagrangian and, morever, the multiplier λ is the normal constraint $\sigma \cdot
u$ on the boundaries of the elements, u denoting a field of unit normal vectors on the union of the boundaries. Methods obtained from such a Lagrangian are also termed hybrid methods. Again we shall adopt a general terminology saying that the method is a hybrid method if it involves the simultaneous approximation of a vector field defined on the union of the elements of the discretization and another defined on the union of the boundaries of the elements. The first reference for hybrid methods is JONES (1964).

An aperçu of these methods is given in CIARLET (1978, Chapter 7).

Besides structural mechanics mixed and hybrid methods have been found useful in many fields. In fluid mechanics for the incompressible Stokes problem, for

instance, the role of the displacement is played by the fluid pressure, that of the stress by the fluid velocity, that of the equilibrium condition by the incompressibility condition.

The above ideas will be developed in detail in Sections 2 and 3 in the context of a simpler model problem. First, however, we need to describe the functional framework in which we shall work.

Let Ω be a bounded, open subset of R^n , with n=2 or n=3, having a boundary Γ , $\Gamma=\partial\Omega$, which is piecewise C^1 . For each natural number m denote by $H^m(\Omega)$ the Sobolev space of order m of scalar valued functions on Ω , defined recursively by

(1.1)
$$\begin{cases} H^{0}(\Omega) = L^{2}(\Omega) ; \\ H^{m}(\Omega) = \left\{ v \in H^{m-1}(\Omega) ; \forall \alpha = (\alpha_{1}, \dots, \alpha_{n}) \in \mathbb{N}^{n} \text{ with } \sum_{i=1}^{n} \alpha_{i} = m, \\ \partial^{\alpha} v = \frac{\partial^{m} v}{\partial x_{1} \cdot \dots \partial x_{n}^{n}} \in L^{2}(\Omega) \right\} \quad \text{for all } m \geq 1. \end{cases}$$

We put

$$(1.2) \begin{cases} \left\| \mathbf{v} \right\|_{0,\Omega} = \left(\int_{\Omega} \mathbf{v}^{2}(\mathbf{x}) d\mathbf{x} \right)^{\frac{1}{2}}, \\ \left\| \mathbf{v} \right\|_{m,\Omega} = \left(\int_{\Omega} \left\| \partial^{a} \mathbf{v} \right\|_{0,\Omega}^{2} \right)^{\frac{1}{2}} & \text{for all } m \geq 1, \\ \left\| \mathbf{v} \right\|_{m,\Omega} = \left(\left\| \mathbf{v} \right\|_{m-1,\Omega}^{2} + \left\| \mathbf{v} \right\|_{m,\Omega}^{2} \right)^{\frac{1}{2}} & \text{for all } m \geq 1. \end{cases}$$

The product space

$$(\operatorname{H}^{m}(\Omega))^{n} = \{\underline{q} = (q_{\underline{i}})_{1 \leq \underline{i} \leq \underline{n}} ; q_{\underline{i}} \in \operatorname{H}^{m}(\Omega) \text{ for all } \underline{i} = 1, \dots, n\}$$

will be equiped with the following seminorm and norm :

(1.3)
$$\begin{cases} \left\| \mathbf{q} \right\|_{m,\Omega} = \left(\sum_{i=1}^{n} \left\| \mathbf{q}_{i} \right\|_{m,\Omega}^{2} \right)^{\frac{1}{2}}, \\ \left\| \mathbf{q} \right\|_{m,\Omega} = \left(\sum_{i=1}^{n} \left\| \mathbf{q}_{i} \right\|_{m,\Omega}^{2} \right)^{\frac{1}{2}}. \end{cases}$$

We shall make use of the following result from functional analysis :

THEOREM 1.1. The mapping $v \longrightarrow v_{|\Gamma}$ defined a priori for functions v regular on $\bar{\Omega}$, can be extended to a continuous linear mapping called the trace map of $H^1(\Omega)$ into $L^2(\Gamma)$. In other words there exists a constant C, depending only on Ω , such that we have

(1.4)
$$\|\mathbf{v}\|_{0,\Gamma} \le \mathbf{c} \|\mathbf{v}\|_{1,\Omega}$$
 for all $\mathbf{v} \in \mathbf{H}^1(\Omega)$.

It can be shown that on the kernel $H_0^1(\Omega)$ of the trace mapping the correspondence $v \longrightarrow |v|_{1,\Omega}$ defines a norm equivalent to the norm $v \longrightarrow ||v||_{1,\Omega}$; i.e. the Friedrichs-Poincaré inequality holds on $H_0^1(\Omega)$. The image of the trace mapping, denoted $H^{\frac{1}{2}}(\Gamma)$, is a Hilbert space with norm

(1.5)
$$\|\psi\|_{L_{2},\Gamma} = \inf_{\{\mathbf{v} \in H^{1}(\Omega); \mathbf{v} \mid \Gamma = \psi\}} \|\mathbf{v}\|_{1,\Omega}.$$

Thus for each vector $\mathbf{q} = (\mathbf{q_i})_{1 \leq i \leq n} \in (\mathrm{H}^1(\Omega))^n$ the n traces $\mathbf{q_i}$ are defined and belong to $L^2(\Gamma)$, and in particular, the linear combination $\mathbf{q} \cdot \mathbf{\nu} = \sum_{i=1}^n \mathbf{q_i} \mathbf{p_i}$ where $\mathbf{p_i} = (\mathbf{p_i})_{1 \leq i \leq n}$ denotes the unit exterior normal vector to \mathbf{Q} , is in $L^2(\Gamma)$. Imposing less regularity on each component $\mathbf{q_i}$, we can still define the normal trace as an element of a space of distributions on Γ , cf. Theorem 1.2 below. First, though, we need to introduce two more spaces. The dual space of $\mathbf{H}^{\frac{1}{2}}(\Gamma)$ is denoted $\mathbf{H}^{-\frac{1}{2}}(\Gamma)$ and is a Hilbert space with norm

(1.6)
$$\|\mu\|_{-\frac{1}{2},\Gamma} = \sup_{\{\psi \in H^{\frac{1}{2}}(\Gamma); \|\psi\|_{\frac{1}{2},\Gamma} = 1\}} \langle \mu,\psi \rangle_{\Gamma},$$

where $\langle \cdot, \cdot \rangle_{\Gamma}$ is the duality pairing between $\operatorname{H}^{\frac{1}{2}}(\Gamma)$ and $\operatorname{H}^{-\frac{1}{2}}(\Gamma)$. The space $\operatorname{H}(\operatorname{div};\Omega)$ is defined by

(1.7)
$$H(\operatorname{div};\Omega) = \{\underline{q} = (q_i)_{1 \le i \le n} \in (L^2(\Omega))^n ; \operatorname{div} \underline{q} = \sum_{i=1}^n \frac{\partial q_i}{\partial x_i} \in L^2(\Omega)\}$$

and is a Hilbert space with norm

THEOREM 1.2 . The mapping $q \longrightarrow q \cdot \nu$ defined a priori from $(H^1(\Omega))^n$ into $L^2(\Gamma)$ can be extended to a continuous, linear mapping from $H(\operatorname{div};\Omega)$ onto $H^{-\frac{1}{2}}(\Gamma)$. Further we have the following characterization of the norm on $H^{-\frac{1}{2}}(\Gamma)$:

(1.9)
$$\|\mu\|_{-\frac{1}{2},\Gamma} = \inf_{\{\underline{q} \in H(\operatorname{div};\Omega); \underline{q} \cdot \nu = \mu\}} \|\underline{q}\|_{H(\operatorname{div};\Omega)}.$$

A demonstration of the first part of this theorem can be found in TEMAM (1977, Theorem 1.2, p.9). The characterization of the norm was given by THOMAS (1977, Chapter I). To avoid, whenever possible, working in this space $\operatorname{H}^{-\frac{1}{2}}(\Gamma)$ which contains all the functions of $L^2(\Gamma)$, we define, with the aid of Theorem 1.2, the space

(1.10)
$$H(\operatorname{div};\Omega) = \{q \in H(\operatorname{div};\Omega) : q \cdot \nu \in L^2(\Gamma)\}$$

which is a Hilbert space with norm

(1.11)
$$\| g \|_{H \text{ (div; }\Omega)} = \{ \| g \|_{H(\text{div;}\Omega)}^p + \| g \cdot \nu \|_{O,\Gamma}^p \}^{\frac{1}{2}}.$$

We shall make use of the following version of Green's formula :

(1.12)
$$\int_{\Omega} (\mathbf{v} \operatorname{div} \mathbf{q} + \operatorname{grad} \mathbf{v} \cdot \mathbf{q}) \, d\mathbf{x} = \int_{\Gamma} \mathbf{v} \, \mathbf{q} \cdot \mathbf{v} \, d\sigma \qquad \text{for all } \mathbf{v} \in \mathrm{H}^{1}(\Omega) \text{ and }$$

$$\mathbf{q} \in H(\operatorname{div};\Omega). \quad \Box$$

Remark 1.1 . The reader not familiar with the notation $H^{\frac{1}{2}}(\Gamma)$ or $H^{-\frac{1}{2}}(\Gamma)$ may find it somewhat surprizing. It is natural in that these spaces can also be defined as Sobolev spaces of order s, for s not necessarily an integer. The study of such spaces is not within the scope of this work but may be found in LIONS & MAGENES (1968, Chapter I) or in ADAMS (1975, Chapter 7).

Other examples of functional spaces will be used in the following, and corresponding trace theorems will be stated. However, before ending this section we would like to give necessary and sufficient conditions for a locally regular scalar function to belong to $\mathrm{H}^1(\Omega)$ and for a locally regular vector function to belong to $\mathrm{H}(\mathrm{div};\Omega)$. More precisely, let T_{h} be a decomposition of Ω into compact sets T , $\Omega = \bigcup_{\mathrm{T}} \mathrm{T}$, such that the interiors of the sets T of T_{h} are pairwise disjoint and such that the boundary $\partial \mathrm{T}$ of an element T of T_{h} is piecewise C^1 . For the sake of simplicity we shall make the following abuses of notation ; we shall write

$$H^{m}(T)$$
 for $H^{m}(\mathring{T})$, $\|\cdot\|_{m,T}$ for $\|\cdot\|_{m,\mathring{T}}$; $\|\cdot\|_{H(\operatorname{div};T)}$ for $\|\cdot\|_{H(\operatorname{div};\mathring{T})}$;

etc.

(1.13)
$$|v_{T_1}|_{T'} - |v_{T_2}|_{T'}$$
 for all $T' = T_1 \cap T_2$ with T_1 , $T_2 \in T_h$.

Similarly a function $\underline{q} \in (L^2(\Omega))^n$, whose restriction $\underline{q}_{|T}$ may be identified with a function $\underline{q}_{T} \in \mathcal{H}(\text{div};T)$ for each $T \in T_h$, belongs to $\mathcal{H}(\text{div};\Omega)$ if and only if for each interface $T' = T_1 \cap T_2$ with T_1 , $T_2 \in T_h$, the normal trace of \underline{q}_{T_1} coincides with the negative of that of \underline{q}_{T_2} :

$$(1.14) \qquad (q_{T_1} \cdot \nu_{T_1}) \mid_{T'} + (q_{T_2} \cdot \nu_{T_2}) \mid_{T'} = 0 \qquad \text{for all } T' = T_1 \cap T_2$$
 with T_1 , $T_2 \in T_b$,

where $\nu_{\rm T}$ is the unit exterior normal vector of T.

Clearly we have $\nu_{T_1}=-\nu_{T_2}$ on $T'=T_1\cap T_2$ so that if ν is any nontrivial vector orthogonal to T' we have

$$q_{T_1} \cdot \nu = q_{T_2} \cdot \nu$$
 on $T' = T_1 \cap T_2$.

2 - CONFORMING AND EQUILIBRIUM VARIATIONAL FORMULATIONS FOR A MODEL BOUNDARY VALUE PROBLEM.

Many physical phenomena can be modelled by a system of first order partial differential equations on a domain $\Omega \subset R^n$ of the form

(2.1)
$$p_{i} = \sum_{j=1}^{n} a_{ij} \frac{\partial u}{\partial x_{j}} \qquad \text{in } \Omega, \ 1 \leq i \leq n,$$

where the vector function $\mathbf{p} = (\mathbf{p_i})_{1 \le i \le n}$ satisfies the equilibrium relation

(2.2)
$$\operatorname{div} p + f = 0 \qquad \text{in } \Omega,$$

and the boundary conditions are of the form

$$(2.3) (\underline{p} \cdot \underline{\nu} - \underline{g}) + t (\underline{u} - \underline{u}) = 0 on \Gamma,$$

where Γ is the boundary of Ω ; cf. DUVAUT & LIONS (1972, pp. 16-17) for example. For a review of such boundary conditions, cf. DAUTRAY & LIONS (1985, pp.230-240). The coefficients a_{ij} of the system, defined on Ω and the function t defined on Γ are known functions dependant on the physical problem. The function f defined on Ω and the functions g and \bar{u} on Γ are data functions. The function t is allowed to be infinite, $t=+\infty$, on a part Γ_0 of Γ thus implying that $u=\bar{u}$ on Γ_0 . The steady state heat equation with convection terms neglected is an example of such a model; u denotes the temperature, p is the heat flow, and $p \cdot \nu$, ν a unit vector normal to Γ , is the flux across Γ . Linear models for elasticity problems involve an analogous system relating the displacement field $u=(u_i)_{1\leq i\leq n}$ and the stress field $\sigma=(\sigma_{ij})_{1\leq i,j\leq n}$. Hooke's law gives the stress field in terms of the linearized displacement field; $\varepsilon_{ij}(u)=\frac{\lambda_i}{2}(\frac{\partial u_i}{\partial x_i}+\frac{\partial u_j}{\partial x_i})$ is the analogue of (2.1).

For these examples corresponding to fundamental problems of engineering, we see that it is at least as important to calculate an approximation of p, or even the moments of $p \cdot \nu_{\Gamma}$ on Γ or of $p \cdot \nu$ on a surface contained in Ω with normal vector ν , as to calculate an approximation of u or of the trace u or of the trace of u on a surface lying in Ω . On the other hand, from a strictly theoretical point of view, an apparently obvious approach would seem to be to begin by simplifying the study of system (2.1), (2.2) by eliminating the unknown p, thus considering u as the solution of the second order partial differential equation

(2.4)
$$-\sum_{i,j=1}^{n} \frac{\partial}{\partial x_{i}} \left(a_{ij} \frac{\partial u}{\partial x_{j}}\right) - f \qquad \text{in } \Omega.$$

The classical methods for the numerical approximation of problem (2.1) - (2.3) are based on the following approach: one determines a function \mathbf{u}_h satisfying (2.4) in a more or less weak sense--using for instance standard finite difference or finite element methods. Once the function \mathbf{u}_h has been calculated, if it is necessary to obtain an approximation of \mathbf{p} , the simplest procedure is to consider the function \mathbf{p}_h whose \mathbf{i}^{th} component, $1 \le \mathbf{i} \le \mathbf{n}$, is given by $(\mathbf{p}_h)_{\mathbf{i}} = \sum_{j=1}^n \mathbf{a}_{\mathbf{i}\mathbf{j}} = \frac{\partial \mathbf{u}_h}{\partial \mathbf{x}_j}$. This naive treatment has the inconvenience of producing a function \mathbf{p}_h which satisfies only in a very weak sense the equilibrium condition and for which the two normal traces $\mathbf{p}_h \cdot \mathbf{v}_T$, across an interface $\mathbf{T}' \subset \Omega$ do not agree. Experience shows that it is better, a priori, to choose a method which solves directly for the desired quantities rather than to try to obtain them in an indirect fashion as in the example above. Our objective here is to propose several such methods.

We shall consider in this chapter, by way of introduction, the case in which the coefficient functions a_{ij} are constant in Ω , and the matrix (a_{ij}) is symmetric and positive definite. Thus the relations (2.1) are invertible :

$$\frac{\partial u}{\partial x_i} = \sum_{j=1}^n A_{ij} p_j, \qquad 1 \le i \le n,$$

and the inverse matrix $(A_{\mbox{ij}})$ is also symmetric and positive definite. For the moment we shall restrict our attention to the two following model problems :

The homogeneous Dirichlet problem. Given f in $L^2(\Omega)$, find $p \in H(\text{div};\Omega)$ and $u \in H^1_0(\Omega)$ satisfying (2.1) and (2.2).

The homogeneous Neumann problem. Given f in $L^2(\Omega)$ with \int_{Ω} f dx = 0, find $p \in H_0(\text{div};\Omega)$ and $u \in H^1(\Omega)/R$ satisfying (2.1) and (2.2), where for the sake of simplicity we suppose that Ω is connected.

Just as $H^1_0(\Omega)$ denotes the kernel of the trace mapping of $H^1(\Omega)$ into $L^2(\Gamma)$, $H_0(\operatorname{div};\Omega)$ denotes the kernel of the normal trace mapping of $H(\operatorname{div};\Omega)$ into $L^2(\Gamma)$. A specific representative of the equivalence class of u in the quotient space $H^1(\Omega)/R$, for u a solution of the homogeneous Neumann problem, may be chosen by requiring further that a relation such as $\int_{\Omega} u \ dx = 0$ or such as $\int_{\Gamma} u \ dx = 0$ be satisfied. (Using the standard abuse of notation we write u for an element of $H^1(\Omega)$ as well as for its equivalence class in $H^1(\Omega)/R$).

The solution (p,u) in the case of the homogeneous Dirichlet or the homogeneous Neumann problem may be characterized as the unique solution of a minimization problem. For the homogeneous Dirichlet problem put

(2.6D)
$$W^{f} = \{q \in H(div;\Omega) ; div q + f = 0\}, \qquad V = H_0^1(\Omega),$$

and for the homogeneous Neumann problem

$$(2.6N) Wf = {q \in H0(div; \Omega) ; div q + f = 0}, V = H1(\Omega)/R$$

Then in either case (p,u) is the unique element of $W^f \times V$ satisfying (2.1) or equivalently (2.5). On $W^f \times V$ define the functional

(2.7)
$$I(q,v) = \frac{1}{2} \int_{\Omega} \sum_{i=1}^{n} ((q_i - \sum_{j=1}^{n} a_{ij} \frac{\partial v}{\partial x_j}) (\sum_{k=1}^{n} A_{ik} q_k - \frac{\partial v}{\partial x_i})) dx.$$

For each pair (\underline{q},v) in $W^{f}\times V$ we have

$$I(q,v) = \frac{1}{2} \int_{\Omega} \sum_{i,k=1}^{n} ((q_i - \sum_{j=1}^{n} a_{ij} - \frac{\partial v}{\partial x_j}) (A_{ik} (q_k - \sum_{\ell=1}^{n} a_{k\ell} - \frac{\partial v}{\partial x_\ell}))) dx \ge 0,$$

with the equality being realized, as [A] is positive definite, only in case (q,v) satisfies the relation $q_i = \sum\limits_{j=1}^n a_{ij} \frac{\partial v}{\partial x_j}$ for all $i, i=1,\ldots,n$; i.e. if (q,v) satisfies (2.1). Thus the solution (p,u) of (2.1) - (2.2) with the homogeneous Dirichlet or the homogeneous Neumann condition is the unique element of $w^f \times v$ satisfying

$$(2.8) I(p,u) = 0$$

or in other words the unique element minimizing I on W^fxV . This presentation was inspired by that in BARLOW (1986).

With this characterization of the solution (p,u) of (2.1) - (2.3), it is natural to try to construct an approximation in the following manner. Let V_h be a finite dimensional subspace of V and W_h^f a finite dimensional, affine subspace of $H(\operatorname{div};\Omega)$ or $H_0(\operatorname{div};\Omega)$ contained in W_h^f . We consider the following problem : find (p,u_h) satisfying

$$\begin{cases} (\underline{p}_h, \underline{u}_h) \in W_h^f \times V_h , \\ [2pt] I(\underline{p}_h, \underline{u}_h) - \inf_{(\underline{q}_h, v_h) \in W_h^f \times V_h} I(\underline{q}_h, v_h) . \end{cases}$$

Problem (2.9) has a unique solution $(\underline{p}_h, \underline{u}_h)$; the uniqueness is an immediate consequence of the characterizations given by (2.15) and (2.16) below. Furthermore, the calculation of the value $I(\underline{p}_h, \underline{u}_h) > 0$ furnishes, a posteriori, a control of the error committed in the approximation.

Formulation (2.9) thus seems to provide a natural coupling between the approximations p_h of p and u_h of u. In fact it does not at all. To see this one has only to note that, with the choice (2.6D) as well as with the choice (2.6N), we have

(2.10)
$$I(q,v) = J_{*}(q) + J^{f}(v),$$

with

(2.11)
$$J^{f}(v) = \frac{1}{2} \int_{\Omega} \sum_{i,j=1}^{n} a_{ij} \frac{\partial v}{\partial x_{j}} \frac{\partial v}{\partial x_{i}} dx - \int_{\Omega} fv dx,$$

(2.12)
$$J_{*}(q) = \frac{1}{2} \int_{\Omega} \sum_{i,j=1}^{n} A_{ij} q_{j} q_{i} dx.$$

Hence if (p,u) is the solution of (2.8), then u is characterized as being the unique solution of the minimization problem : find u satisfying

(2.13)
$$\begin{cases} u \in V ; \\ J^{f}(u) = \inf_{v \in V} J^{f}(v) ; \end{cases}$$

and \underline{p} as being that of the minimization problem : find p satisfying

(2.14)
$$\begin{cases} \underline{p} \in W^{f}; \\ J_{*}(\underline{p}) = \inf_{\underline{q} \in W^{f}} J_{*}(\underline{q}). \end{cases}$$

Since $I(\underline{p}, u) = 0$, we have the relation

(2.15)
$$J^{f}(u) + J_{\downarrow}(p) = 0$$

known in mechanics as the complementary energy principle.

The solution (p_h, u_h) of (2.9) is characterized with the aid of the following two minimization problems : find u_h satisfying

$$\begin{cases} u_h \in V_h, \\ J^f(u_h) = \inf_{v_h \in V_h} J^f(v_h); \end{cases}$$

find p_h satisfying

(2.17)
$$\begin{cases} p_h \in W_h^f, \\ J_*(p_h) = \inf_{q_h \in W_h^f} J_*(q_h). \end{cases}$$

Since the approximate solution $u_h \in V_h$ is an element of V, we say that u_h is obtained by a conforming approximation method. Similarly, since the approximate solution $p_h \in W_h^f$ is an element of W^f we say that p_h is obtained by an equilibrium approximation method. If \tilde{J} is the energy defined by $\tilde{J} = J^f(u) = -J_*(p)$, cf. (2.15), we remark that the conforming approximation permits us to calculate an approximation $J^f(u_h)$ of \tilde{J} which is of necessity excessive, whereas the equilibrium approximation permits us to calculate an approximation $-J_*(p_h)$ of \tilde{J} which necessarily too small:

$$(2.18) - J_*(\underline{p}_h) \le - J_*(\underline{p}) = \bar{J} = J^f(u) \le J^f(u_h).$$

Remark 2.1 . The possibility to thus obtain an a posteriori bound on the energy \bar{J} is what motivated the development of equilibrium approximation methods, cf. FRAEIJS DE VEUBEKE (1965), (1973), (1975), FRAEIJS DE VEUBEKE & HOGGE (1972).

Remark 2.2 . Let g be a function in $L^2(\Gamma)$ satisfying \int_{Γ} g ds = 0. If (\underline{p},u) is a solution of the system

$$\begin{cases}
p = \text{grad } u & \text{in } \Omega, \\
\text{div } p = 0 & \text{in } \Omega, \\
p \cdot \nu = g & \text{on } \Gamma,
\end{cases}$$

for Ω a connected open subset of R^n , then as above one can verify that the function u, defined to within a constant, minimizes on $H^1(\Omega)/R$ the functional

(2.19)
$$J_{g}(v) = \frac{1}{2} \int_{\Omega} |\operatorname{grad} v|^{2} dx - \int_{\Gamma} gv d\sigma$$

whereas the vector function p minimizes the functional

(2.20)
$$J_{*}(q) = \frac{1}{2} \int_{\Omega} |q|^{2} dx$$

on the affine manifold of vector functions \underline{q} with div $\underline{q}=0$ on Ω and $\underline{q} \cdot \nu = \underline{q}$ on Γ . If \underline{u} minimizes $J_{\underline{q}}(\underline{v})$ then we have div (grad $\underline{u})=0$; and if \underline{p} minimizes $J_{\underline{q}}(\underline{q})$, we have rot $\underline{p}=0$ since then \underline{p} has a scalar potential. These principles are known in fluid mechanics as <u>Dirichlet's principle</u> and <u>Kelvin's principle</u>, respectively; cf SERRIN J. (1959, Section 24). In the standard terminology of fluid mechanics, the vector variables representing the flow rate are denoted \underline{u} or \underline{v} , while the scalar variable is the scalar potential of the velocity field or the pressure and is generally denoted ϕ , ψ , or \underline{p} .

Remark 2.3 . We shall see later on, cf. Remark 3.1; that the two minimization problems (2.13) and (2.14) (or similarly (2.19) and (2.20)) are dual problems in the sense of the mathematical theory of duality in convex analysis. Usually (2.13) is considered to be the primal problem and (2.14) its dual. In keeping with standard practice we shall say that a formulation is primal, respectively dual, when the elimination of certain variables yields the conforming formulation (2.13), respectively the equilibrium formulation (2.14).

To conclude this presentation of conforming and equilibrium approximations we give the variational formulations associated with these minimization problems :

Conforming variational formulation. Find a solution u of

(2.21)
$$\begin{cases} u \in V, \\ \sum_{i,j=1}^{n} \int_{\Omega} a_{ij} \frac{\partial u}{\partial x_{i}} \frac{\partial v}{\partial x_{i}} dx = \int_{\Omega} fv dx & \text{for all } v \in V, \end{cases}$$

which leads one to seek an approximation \mathbf{u}_{h} of \mathbf{u} satisfying

$$\begin{cases} u_h \in V_h, \\ \sum_{i,j=1}^n \int_{\Omega} a_{ij} \frac{\partial u_h}{\partial x_j} \frac{\partial v_h}{\partial x_i} dx = \int_{\Omega} f v_h dx & \text{for all } v_h \in V_h. \end{cases}$$

Equilibrium variational formulation. Find a solution p of

which leads one to seek an approximation $\ensuremath{\underline{p}}_h$ of $\ensuremath{\underline{p}}$ satisfying

$$\begin{cases} p_h \in W_h^f, \\ \sum_{i,j=1}^n \int_{\Omega} A_{ij} p_{hj} q_{hi} dx = 0 \end{cases}$$
 for all $q_h \in W_h^0$;

where W^0 , respectively W_h^0 , is W^f , respectively W_h^f , for f = 0 and hence is a subspace of $H(\text{div};\Omega)$ or $H_0(\text{div};\Omega)$ depending on whether the problem considered is the homogeneous Dirichlet problem or the homogeneous Neumann problem.

3 - FIRST EXAMPLES OF MIXED AND HYBRID FORMULATIONS

The actual construction of a solution p_h of problem (2.24) (equilibrium formulation) poses several practical difficulties . If T_h is a triangulation of $\bar{\Omega}$, $\bar{\Omega} = \bigcup_{T \in T_h} T$, and W_h^0 is a finite dimensional subspace of W^0 associated with T_h , then for a test function q_h to belong to W_h^0 it must satisfy the following two constraints:

- i) The normal traces of \underline{q}_h must be "continuous across the interface" between any two abutting finite elements T_1 and T_2 of T_h , i.e. $\underline{q}_h \cdot \nu_{T_1}$ and $\underline{q}_h \cdot \nu_{T_2}$ must coincide on $T' = T_1 \cap T_2$ (cf. the second part of Theorem 1.3).
- ii) Within each finite element T of T_h the divergence of q_h must vanish; i.e. $\operatorname{div}(q_h|_T) = 0$ for all $T \in T_h$; thus, in particular $\int_{\partial T} q_h \cdot \nu_T d\sigma = 0$.

We do not intend to attack head on the practical problem of constructing a basis for the subspace \mathbb{W}_h^0 . To circumvent this difficulty we shall free ourselves of one of the constraints that should be satisfied by the test functions by using the technique of Langrangian multipliers.

The dual mixed formulation. The affine subspace defined by (2.6D) in the case of the homogeneous Dirichlet problem and by (2.6N) in the case of the homogeneous Neumann problem may be characterized as being the set of functions $q \in W$ satisfying the constraint div q + f = 0 in Ω , where W is the space defined by

$$(3.1D) W = H(div;\Omega),$$

$$(3.1N) W = H_0(\operatorname{div};\Omega).$$

Thus we are led to substitute for the minimization problem (2.14) the saddle point problem of the Lagrangian $\mathcal{L}_{\star}(\cdot,\cdot)$ defined on $\mathbb{W}\times L^2(\Omega)$ by

(3.2)
$$L_*(q, v) = J_*(q) + \int_{\Omega} (\operatorname{div} q + f) v \, dx.$$

Thus we seek $(p,u) \in W \times L^2(\Omega)$ satisfying

$$\mathcal{I}_{\star}(\underline{p},v) \leq \mathcal{I}_{\star}(\underline{p},u) \leq \mathcal{I}_{\star}(\underline{q},u) \qquad \text{for all } (q,v) \in \mathbb{W} \times L^{2}(\Omega).$$

The solution is characterized as being the solution of the system

$$(3.4) \begin{cases} (\underline{p}, u) \in W \times L^{2}(\Omega) , \\ \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} p_{j} q_{i} dx + \int_{\Omega} u div \underline{q} dx = 0 & \text{for all } \underline{q} \in W; \\ \int_{\Omega} (\text{div } \underline{p}) v dx = - \int_{\Omega} f v dx & \text{for all } v \in L^{2}(\Omega). \end{cases}$$

If $(\underline{p}, \mathbf{u})$ satisfies (3.4) it is clear that \underline{p} satisfies (2.23).

Denote, for the moment, by $\hat{\mathbf{u}}$ the solution of (2.21). Using Green's formula together with relation (2.5), we see that

$$\begin{split} \sum_{i,j=1}^n \, \int_{\Omega} \, \mathbf{A}_{ij} \, \, \mathbf{p}_j \, \, \mathbf{q}_i \, \, \mathrm{d}\mathbf{x} &= \int_{\Omega} (\mathbf{grad} \, \, \hat{\mathbf{u}}) \cdot \mathbf{q} \, \, \mathrm{d}\mathbf{x} \\ &= - \, \int_{\Omega} \, \hat{\mathbf{u}} \, \, \mathrm{div} \, \, \mathbf{q} \, \, \mathrm{d}\mathbf{x} \quad \qquad \text{for all } \, \mathbf{q} \in \mathbb{W}. \end{split}$$

Hence the Lagrange multiplier associated with the constraint div q=f is none other than

$$(3.5)$$
 $u = \hat{u}$.

The primal mixed formulation . In virtue of (3.5), the solution (p,u) actually belongs to W \times V. Using Green's formula we establish that the solution (p,u) satisfies

$$\begin{cases} \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} p_{j} q_{i} dx - \int_{\Omega} (\operatorname{grad} u) \cdot q dx = 0 & \text{for all } q \in W; \\ - \int_{\Omega} p \cdot \operatorname{grad} v dx = - \int_{\Omega} f v dx & \text{for all } v \in V. \end{cases}$$

As W is dense in $\left(L^2(\Omega)\right)^n$, the pair (\underline{p},u) is the solution of the system

$$(3.6) \begin{cases} (\underline{p}, u) \in (L^{2}(\Omega))^{n} \times V, \\ \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} p_{j} q_{i} dx - \int_{\Omega} (\underline{grad} u) \cdot \underline{q} dx = 0 & \text{for all } \underline{q} \in (L^{2}(\Omega))^{n}, \\ - \int_{\Omega} \underline{p} \cdot \underline{grad} v dx = - \int_{\Omega} fv dx & \text{for all } v \in V. \end{cases}$$

This formulation, (3.6), can be interpreted as being a characterization of the saddle point of the functional $\mathcal{L}^f(\cdot,\cdot)$ defined on $(L^2(\Omega))^n \times V$ by

(3.7)
$$\mathcal{I}^{f}(q,v) = \bar{J}_{x}(q) - \int_{\Omega} q \cdot \operatorname{grad} v \, dx + \int_{\Omega} f \, v \, dx,$$

where \tilde{J}_{\star} is the continuous extension of J_{\star} to the space $(L^2(\Omega))^n$. This functional can be found in the literature written in the form

$$(3.7 \text{ bis}) \qquad x^{f}(q,v) = -\frac{1}{2} \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} q_{j} q_{i} dx + \sum_{i=1}^{n} \int_{\Omega} (\sum_{j=1}^{n} A_{ij} q_{j} - \frac{\partial v}{\partial x_{i}}) q_{i} dx + \int_{\Omega} fv dx.$$

The elimination of p from formulation (3.6) gives us again the conforming variational formulation (2.21) of the problem. This type of method is analyzed in BABUSKA, ODEN, & LEE (1977). Such methods have proved useful in theoretical studies such as the development of new models for nonlinear elasticity, cf. CIARLET & DESTUYNDER (1979a), (1979b), DESTUYNDER (1986).

Remark 3.1 . We have on the one hand

$$\sup_{\mathbf{v}\in\mathbf{V}} \mathcal{L}^{\mathbf{f}}(\mathbf{q},\mathbf{v}) = \begin{cases} J_{\mathbf{x}}(\mathbf{q}) & \text{if } \mathbf{q}\in\mathbf{W}^{\mathbf{f}}, \\ +\infty & \text{otherwise,} \end{cases}$$

and on the other

$$\inf_{q\in W} \, \mathcal{I}^f(q,v) \, = \, -\, \frac{1}{2} \, \sum_{i\,,\,j=1}^n \, \int_\Omega \, a_{i\,j} \, \frac{\partial v}{\partial x_j} \, \frac{\partial v}{\partial x_i} \, \, dx \, + \, \int_\Omega \, fv \, \, dx \, .$$

We deduce that

$$\inf_{\mathbf{q} \in \mathbb{W}} \sup_{\mathbf{v} \in \mathbb{V}} \mathcal{I}^{\mathbf{f}}(\mathbf{q}, \mathbf{v}) = \inf_{\mathbf{q} \in \mathbb{W}^{\mathbf{f}}} J_{\star}(\mathbf{q}) = J_{\star}(\mathbf{p})$$

and that

$$\sup_{\mathbf{v} \in \mathbb{V}} \inf_{\mathbf{q} \in \mathbb{W}} \mathbf{f}(\mathbf{q}, \mathbf{v}) = -\inf_{\mathbf{v} \in \mathbb{V}} \mathbf{J}^{\mathbf{f}}(\mathbf{v}) = -\mathbf{J}^{\mathbf{f}}(\mathbf{u}).$$

Then using the complementary energy principle (2.15), we obtain

(3.8)
$$\inf_{\mathbf{q} \in \mathbb{W}} J_{\star}(\mathbf{q}) = \inf_{\mathbf{q} \in \mathbb{W}} \sup_{\mathbf{v} \in \mathbb{V}} \mathcal{I}^{\mathbf{f}}(\mathbf{q}, \mathbf{v}) = \sup_{\mathbf{v} \in \mathbb{V}} \inf_{\mathbf{q} \in \mathbb{W}} \mathcal{I}^{\mathbf{f}}(\mathbf{q}, \mathbf{v}) = -\inf_{\mathbf{v} \in \mathbb{V}} J^{\mathbf{f}}(\mathbf{v}).$$

In this form we recognize the minimization problems (2.13) and (2.14) as dual problems in the sense of optimization theory, cf. for example CIARLET (1982, p.222).

The dual hybrid formulation. To a decomposition T_h of $\bar{\Omega}$, $\bar{\Omega}=\cup$ T, is associated the affine subspace $Y^{\bar{f}}$ of $(L^2(\Omega))^n$ defined by

(3.9)
$$Y^{f} = \{q \in (L^{2}(\Omega))^{n} ; \forall T \in T_{h}, \operatorname{div}(q_{|T}) + f_{|T} = 0,$$
 in the distributional sense \}.

In other words, Y f is the set of vector functions $q \in (L^2(\Omega))^n$ satisfying for each $T \in T_h$

$$\int_{T} \underline{q} \cdot \underline{\text{grad}} \ v_{\underline{T}} \ dx = \int_{T} f v_{\underline{T}} \ dx \qquad \qquad \text{for all } v_{\underline{T}} \in H^{1}_{0}(T).$$

For each $T \in T_h$, the restriction q of the function $q \in Y^f$ to T belongs to the space $H(\operatorname{div};T)$. Further, the normal trace $q \mid_{T} \cdot \nu_{T}$ is defined, cf. Theorem 1.2, to be an element of $H^{-\frac{1}{2}}(\partial T)$, the dual space of $H^{\frac{1}{2}}(\partial T)$; we denote by $\langle \cdot, \cdot \rangle_{\partial T}$ this duality. We ascertain that the affine subspace W^f may be characterized as being the set of functions $q \in Y^f$ satisfying the constraint

$$\sum_{T \in T_h} \langle q|_{T} \cdot \nu_{T}, v \rangle_{\partial T} = 0$$
 for all $v \in V$.

This leads us to associate to the decomposition $T_{\rm h}$ of $\bar{\Omega}$ the Lagrangian $D_*(\cdot,\cdot)$ defined on Y^fxV by

$$(3.10) D_{\star}(q,v) = J_{\star}(q) - \sum_{T \in T_h} \langle q_{|T} \cdot \nu_{T}, v \rangle_{\partial T}.$$

(The choice of $D_{\star}(q,v)$ as given in (3.10) instead of as $J_{\star}(q) + \sum_{T \in T} \langle q |_{T} \cdot \nu_{T}, v \rangle_{\partial T}$ is made simply to obtain (3.13) below). Thus we seek an element $h(p,u) \in Y^f \times V$ satisfying

$$(3.11) D_{\star}(\underline{p},v) \leq D_{\star}(\underline{p},u) \leq D_{\star}(\underline{q},u) \text{for all } (\underline{q},v) \in Y^{f} \times V.$$

A pair (p,u) is such an element if and only if it is a solution of the system

$$(3.12) \begin{cases} (\underline{p}, \mathbf{u}) \in \mathbf{Y}^{f} \times \mathbf{V}, \\ \sum_{i,j=1}^{n} \int_{\Omega} \mathbf{A}_{ij} \mathbf{p}_{j} \mathbf{q}_{i} d\mathbf{x} + \sum_{T \in T_{h}} \langle \underline{q}|_{T} \cdot \underline{\nu}_{T}, \mathbf{u} \rangle_{\partial T} = 0 & \text{for all } \underline{q} \in \mathbf{Y}^{0}, \\ \sum_{T \in T_{h}} \langle \underline{p}|_{T} \cdot \underline{\nu}_{T}, \mathbf{v} \rangle_{\partial T} = 0 & \text{for all } \mathbf{v} \in \mathbf{V}; \end{cases}$$

where Y^0 is Y^f for $f \equiv 0$. If (p,u) is a solution of problem (3.12), then p is the solution of (2.22). It is clear that u is not uniquely defined as an element of V; only its traces on the interfaces are determined by the first equation of (3.12). However, if \hat{u} is the solution of (2.20), we can choose

$$(3.13)$$
 $u = \hat{u}$.

Remark 3.2 . In fact with the dual hybrid formulation the Lagrangian multiplier should not be considered as a function in V defined on all Ω but as an element of the product space Π $H^{\frac{1}{2}}(\partial T)$ satisfying certain matching conditions at the $T \in T_h$ interfaces. This will be the point of view taken for the analysis of the method given in Chapter 5. A variation of considerable practical interest consists of taking the multiplier to be an element of Π $L^2(\partial T)$ satisfying again certain $T \in T_h$

The primal hybrid formulation. To a decomposition T_h of Ω we associate in a manner similar to the above a subspace X of $L^2(\Omega)$ defined by

(3.14)
$$X = \{v \in L^2(\Omega), \forall T \in T_h, v_T \in H^1(T)\}.$$

The space V may be characterized as being the space of functions $v \in X$ satisfying the constraint

$$\sum_{T \in T_h} \langle \underline{q} \cdot \underline{v}|_T, v_T \rangle_{\partial T} = 0$$
 for all $\underline{q} \in W$.

Thus we replace the minimization problem (2.13) by the saddle point problem for the Lagrangian $D^{f}(\cdot, \cdot)$ defined on $X \times W$ by

$$D^{f}(v,q) = \bar{J}^{f}(v) - \sum_{T \in T_{h}} \langle q \cdot \nu_{|T}, v_{T} \rangle_{\partial T},$$

where $ilde{J}^{ extbf{f}}$ is the extension of $ilde{J}^{ extbf{f}}$ to X defined as follows :

(3.16)
$$\bar{J}^{f}(v) = \frac{1}{2} \sum_{T \in T_{h}} \left(\int_{T} \sum_{i,j=1}^{n} a_{ij} \frac{\partial v}{\partial x_{j}} \frac{\partial v}{\partial x_{i}} dx \right) - \int_{\Omega} fv dx.$$

Thus we seek a solution (p,u) of the system

$$(3.17) \begin{cases} (\underline{p}, \mathbf{u}) \in \mathbb{W} \times \mathbb{X} , \\ \sum_{T \in T_{h}} \left\{ \int_{T} \sum_{i,j=1}^{n} a_{ij} \frac{\partial \mathbf{u}}{\partial \mathbf{x}_{j}} \frac{\partial \mathbf{v}}{\partial \mathbf{x}_{i}} d\mathbf{x} - \langle \underline{p} \cdot \underline{\nu}_{T}, \mathbf{v}_{T} \rangle_{\partial T} \right\} = \int_{\Omega} \text{ for all } \mathbf{v} \in \mathbb{X}, \\ \sum_{T \in T_{h}} \langle \mathbf{u}_{T}, \underline{q} \cdot \underline{\nu}_{T} \rangle_{\partial T} = 0 \\ \text{ for all } \underline{q} \in \mathbb{W}. \end{cases}$$

The solution (p,u) of problem (3.17) is such that u is a solution of problem (2.21) and if \hat{p} is the solution of (2.23) we may choose

$$p = \hat{p}$$
,

as indeed only the normal traces of \underline{p} on the interfaces are determined by (3.17).

Remark 3.3 . The hybrid formulations have been presented here as subordinate to a decomposition $T_{\hat{\mathbf{h}}}$ of $\hat{\Omega}$ which will, of course, coincide with the triangulation of $\hat{\Omega}$ by elements T on which will be constructed approximations of functions. One

can, however, generalize this type of hybrid formulation to one subordinate to an arbitrary decomposition of $\bar{\Omega}$ into subdomains. A particular case is that for which the decomposition is the simplest possible, a decomposition into a single "subdomain". In this case the primal hybrid formulation of the homogeneous Dirichlet problem becomes: find a saddle point $(u,\lambda) \in H^1(\Omega) \times H^{-\frac{1}{2}}(\Gamma)$ of the Lagrangian $\mathcal{I}^f(\cdot,\cdot)$ defined by $\mathcal{I}^f(v,\mu) = J^{-\frac{1}{2}}(v) - \langle \mu,v \rangle_{\Gamma}$ for all $(v,\mu) \in H^1(\Omega) \times H^{-\frac{1}{2}}(\Gamma)$. It was for this type of formulation, adapted to a nonhomogeneous Dirichlet problem, that the first mathematical analysis of a finite element method with Lagrangian multipliers was given, cf. BABUSKA (1973).

4 - GENERAL ORIENTATION

The elementary examples of mixed formulations and of hybrid formulations presented in Section 3 may all be written in the following general form : find a pair (ϕ,λ) satisfying

where W and M are Hilbert spaces, a (\cdot, \cdot) is a symmetric bilinear form on W × W, b (\cdot, \cdot) is a bilinear form on W × M, f is a linear form on W, and g is a linear form on M. Then given a finite dimensional subspace W_h of W and a finite dimensional subspace M_h of M we consider the problem : find a pair (ϕ_h, λ_h) satisfying

$$\begin{cases} \phi_h \in W_h, & \lambda_h \in M_h, \\ a(\phi_h, \psi_h) + b(\psi_h, \lambda_h) = f(\psi_h) & \text{for all } \psi_h \in W_h, \\ b(\phi_h, \mu_h) & = g(\mu_h) & \text{for all } \mu_h \in M_h. \end{cases}$$

This finite dimensional problem may then be posed in the obvious manner as a problem of solving a linear system of $N_{\hat{h}}$ equations in $N_{\hat{h}}$ unknowns, where

$$(4.3) N_h = \dim W_h + \dim M_h.$$

The matrix of this system is symmetric but can not be positive definite. Thus the first problem is to find sufficient conditions, bearing on the pair (W_h, M_h) of subspaces, for the matrix of the associated linear system to be invertible. In a

more general way Chapter III treats the numerical analysis of this type of problem in an abstract setting. The a priori estimates of the error committed in approximating the solution (ϕ,λ) of (4.1) by the solution (ϕ_h,λ_h) of (4.2) are obtained as a function of the error of approximation of (ϕ,λ) in $\mathbb{W}_h \times \mathbb{M}_h$.

In the examples given in Section 3, the spaces W and M appearing in the mixed or hybrid formulations are either $L^2(\Omega)$, $H^1(\Omega)$, or closed subspaces of these spaces for the space of scalar functions and either $(L^2(\Omega))^n$, $H(\text{div};\Omega)$, or closed subspaces of these spaces for the space of vector functions. Chapter II contains the major results concerning interpolation by finite elements of scalar functions in $H^1(\Omega)$ and of vector functions in $H(\text{div};\Omega)$.

Chapters II and III are independent. Both are fundamental for the analysis of the examples developed in the subsequent chapters.

Chapters IV and V treat essentially the approximation of the model problem governed by the first order system of partial differential equations ${\sf S}_{\sf S}$

$$(4.4) p_i = \sum_{j=1}^n a_{ij} \frac{\partial u}{\partial x_j}, 1 \le i \le n,$$

(4.5)
$$\operatorname{div} p + f = 0$$
,

posed on a bounded, open domain $\Omega\subset\mathbb{R}^n$ with n=2 or n=3 with boundary Γ piecewise C^1 . The coefficients a_{ij} are supposed to be bounded and measurable in Ω and to satisfy the uniform ellipticity hypothesis

(4.6)
$$\sum_{i,j=1}^{n} a_{ij}(x) \xi_{i}\xi_{j} \geq \alpha_{p} |\xi|^{2} \text{ a.e. on } \Omega \qquad \text{for all } \xi \in \mathbb{R}^{n},$$

for some $\alpha_p > 0$. Unlike in Sections 2 and 3, we shall not suppose the symmetry of [a], i.e. that $a_{ij}=a_{ji}$. Also we shall not maintain throughout the following the notation particular to Sections 2 and 3.

With the above hypotheses concerning the coefficients a the matrix [a] may be inverted to obtain

$$\frac{\partial \mathbf{u}}{\partial \mathbf{x}_{i}} = \sum_{j=1}^{n} \mathbf{A}_{ij} \mathbf{p}_{j}, \qquad 1 \leq i \leq n ;$$

the coefficients $\textbf{A}_{\mbox{ij}}$ are then bounded and measurable in Ω and satisfy the uniform ellipticity hypothesis

(4.8)
$$\sum_{i,j=1}^{n} A_{ij}(x) \xi_{i}\xi_{j} \ge \alpha_{d} |\xi|^{2} \text{ a.e. on } \Omega \qquad \text{for all } \xi \in \mathbb{R}^{n},$$

for some $\alpha_d > 0$.

Chapter IV is concerned with mixed methods. Primal mixed methods are treated rather rapidly, the dual mixed methods in some more detail. Several algorithms for the solution of the linear system are also given here. In Chapter V hybrid methods are studied. Primal hybrid and dual hybrid methods as well as the method of hybridization of mixed methods are included. In Chapter 6 some further examples of mixed methods and of hybrid methods are mentioned, and some extensions and variants of the theory in the literature are cited.

CHAPTER II

FINITE ELEMENT TYPE INTERPOLATION OF SCALAR AND VECTORIAL FUNCTIONS

5 - LAGRANGIAN INTERPOLATION; AN H¹ -APPROXIMATION OF SCALAR FUNCTIONS

In the examples of mixed and hybrid formulations presented in Section 3, an essential role is played by two Hilbert spaces : the space $\mathrm{H}^1(\Omega)$ of scalar functions $v \in L^2(\Omega)$ such that grad $v \in \left(L^2(\Omega)\right)^n$ and the space $\mathrm{H}(\mathrm{div};\Omega)$ of vector functions $q \in \left(L^2(\Omega)\right)^n$ such that $\mathrm{div}\ q \in L^2(\Omega)$. In this chapter we shall describe several choices of "finite-element-type" subspaces of $\mathrm{H}^1(\Omega)$ and of $\mathrm{H}(\mathrm{div};\Omega)$.

For this general presentation we shall use the most elementary setting for finite element methods (some generalizations will be presented in Section 7):

- i) Ω will denote a bounded, open, polyhedral subset of \mathbb{R}^n , n=2 or 3;
- ii) T_h will be a triangulation of $\overline{\Omega}$ by n simplexes T of diameter no greater than h, (T is a triangle for n = 2, a tetrahedron for n = 3), with

$$(5.1) \qquad \overline{\Omega} = \bigcup_{\mathbf{T} \in T_{\mathbf{h}}} \mathbf{T};$$

for the definition of a triangulation. cf., CIARLET (1978, p.38) for example. The boundary of an n simplex T, i.e. the union of its n+l (n-1)-dimensional faces, will be denoted by ∂T , and in a general manner we shall use the symbol T' to represent an arbitrary (n-1)-dimensional face of T and we shall write T'<T; T' will be an edge of T when n = 2, a triangular face when n = 3. We shall also use the following notation:

- mes T = the Euclidian measure of T in \mathbb{R}^n (geometric area if n = 2, geometric volume if n = 3);
- (5.2) h_{T} = the diameter of T, which here, with a triangulation by simplexes, is just the length of the longest edge;
 - $ho_{\rm T}$ = the radius of the circle inscribed in T if n = 2, of the sphere inscribed in T if n = 3;

and

$$\begin{array}{ccc} \text{(5.3)} & \text{h} = \max & \text{h}_{\text{T}}. \\ & \text{T} \in T_{\text{h}} \end{array}$$

Let H be a set of positive numbers. A family of triangulations $\{T_h; h \in H\}$ will be said to be $\underline{regular}$ if

(5.4)
$$\begin{cases} & \inf & h = 0 \\ & h \in H \end{cases}$$
 and
$$\begin{pmatrix} & \inf & \min & \frac{\rho_T}{h} > 0. \\ & h \in H \text{ T } \in T_h \end{cases}$$

This definition is equivalent to that given by CIARLET (1978, p.124). We shall say that the family $\{T_h : h \in H\}$ is uniformly regular if

(5.5)
$$\begin{cases} &\inf & h = 0 \\ &h \in H \end{cases}$$
 and
$$&\inf & \{(\min & \rho_T) / (\max & h_T)\} > 0.$$

$$h \in H \quad T \in T_h \quad T \in T_h$$

Thus a family of triangulations $\{T_h : h \in H\}$ is uniformly regular if and only if it is regular and it satisfies the <u>inverse hypothesis</u>: there exists a constant C, independent of h, such that for each $h \in H$

(5.6)
$$\min_{T \in T_{h}} h_{T} \ge C h,$$

cf. CIARLET (1978, p.140).

<u>Remark 5.1</u>. The regularity hypothesis (5.4) is quite natural for the theory of interpolation by finite elements. However, in the chapters that follow we shall be rather abusive of the inverse hypothesis (5.6), using it often just to simplify

the exposition of certain results.

For each natural number k, we denote by $P_{k}(T)$ the space of restrictions to T of polynomial functions of degree k in n variables ; we have

(5.7)
$$\dim P_{k}(T) = \frac{(n+k)!}{n!k!} = \begin{cases} 1/2 & (k+1) & (k+2) & \text{if } n=2 \\ \\ 1/6 & (k+1) & (k+2) & (k+3) & \text{if } n=3. \end{cases}$$

To an n simplex T and a positive integer k is associated $\Sigma_T^{(k)}$, the principal lattice of order k of the triangle T, i.e. the set of points of T each of whose barycentric coordinates, relative to the vertices of T, is a multiple of 1/k. We have

(5.8) card
$$\Sigma_{T}^{(k)} = \frac{(n+k)!}{n!k!}$$
.

The Lagrangian interpolant $\mathcal{L}_{T}^{(k)}v$ on T of order k, $k\geq 1$, of a function $v\in C^{0}(T)$ is by definition the unique function in $P_{t}(T)$ such that

(5.9)
$$\mathcal{I}_{T}^{(k)}v(a) = v(a) \qquad \text{for all } a \in \Sigma_{T}^{(k)}.$$

Given a triangulation T_h of $\overline{\Omega}$ and a function $v \in C^0(\overline{\Omega})$, for a positive integer k, we shall denote by $\mathcal{L}_h^{(k)}$ v the function on $\overline{\Omega}$ which coincides with the polynomial function $\mathcal{L}_T^{(k)}(v_{|T|})$ on T for each simplex $T \in T_h$. One says that $\mathcal{L}_h^{(k)}v$ is the Lagrangian interpolant of order k on T_h of the function v. It is essential to remark that $\mathcal{L}_h^{(k)}v$ is, for each positive integer k, a continuous function on $\overline{\Omega}$. Thus one can easily verify that the mapping $\mathcal{L}_h^{(k)}$ is surjective from $C^0(\overline{\Omega})$ onto the space $L_h^{(k)}$ of simplicial Lagrangian interpolants defined by

(5.10)
$$L_h^{(k)} = \{ v \in C^0(\overline{\Omega}) ; \forall T \in T_h, v_{|_T} \in P_k(T) \}.$$

For each k,k \geq 1, the space $L_h^{(k)}$ is a finite dimensional subspace of $\text{H}^1(\Omega).$ In dimensions n=2 and n=3 one may define $\mathcal{I}_T^{(k)}v$ for $v\in \text{H}^2(T)\subset C^0(T)$ but not for $v\in \text{H}^1(T)^{(1)}$. With the aid of the following notation, given for each positive integer m ,

⁽¹⁾ We have adopted the notational convention $H^{m}(A) = H^{m}(A)$, A the interior of A.

(5.11)
$$H^{m}(T_{h}) = \{v \in L^{2}(\Omega) ; \forall T \in T_{h}, v_{|_{T}} \in H^{m}(T)\},$$

one may characterize $L_h^{(k)}$ as being the fixed point space of the operator $\mathcal{L}_h^{(k)}$ on $\mathrm{H}^1(\Omega)\cap\,\mathrm{H}^2(T_h)$:

(5.12)
$$L_{h}^{(k)} = \{ v \in H^{1}(\Omega) \cap H^{2}(T_{h}) ; v = \mathcal{L}_{h}^{(k)} v \}.$$

Thus we have

THEOREM 5.1. For each positive integer k and each triangulation T_h of $\overline{\Omega}$, the mapping $\mathcal{L}_h^{(k)}$ defines a projection of $H^1(\Omega) \cap H^2(T_h)$ onto $L_h^{(k)}$.

We also have

THEOREM 5.2. Suppose we are given a regular family $\{T_h : h \in H\}$ of triangulations of $\overline{\Omega}$. For each positive integer k, there exists a constant C, independent of h, such that

(5.13)
$$|v-\mathcal{L}_{h}^{(k)}v|_{1,\Omega} + h^{-1} ||v-\mathcal{L}_{h}^{(k)}v||_{0,\Omega} \le C h^{m} |v|_{m+1,\Omega}$$

for each $v \in H^{m+1}(\Omega)$, $1 \le m \le k$.

We shall not give here the demonstrations of these results fundamental for the study of the convergence of finite element methods but refer the reader to RAVIART & THOMAS (1983, pp.79-103), for example, for the details of the proofs.

We would like now to make precise the idea of interpolation of the trace on the set $\partial T_h = \bigcup \partial T$ of a scalar function v continuous on $\overline{\Omega}$. For a triangulation $T \in T_h$ of $\overline{\Omega}$ and a positive integer k, the space $L_{\partial T_k}^{(k)}$ is defined by

(5.14)
$$L_{\partial T_{h}}^{(k)} = \{ \psi = (\psi_{T})_{T \in T_{h}} \in \Pi \quad L^{2}(\partial T) ; \exists v \in L_{h}^{(k)} \text{ s.t.}$$

$$\forall T \in T_{h}, \forall_{\partial T} = \psi_{T} \}.$$

The space $L_{\partial T_h}^{(k)}$ may also be characterized as being the space to which $\psi=(\psi_T)_{T\in T_h}$ belongs if and only if it satisfies the following properties :

$$\begin{cases} \text{ i)} \quad \psi_{\text{T}} \in P_{k}(\partial \mathbf{T}) & \text{for each } \mathbf{T} \in T_{h} \\ \\ \text{ii)} \quad \psi_{\mathbf{T}_{1}} = \psi_{\mathbf{T}_{2}} \text{ on } \mathbf{T}' & \text{whenever T' is a face common to the simplexes } \mathbf{T}_{1} \text{ and } \mathbf{T}_{2} \text{ of } T_{h}, \end{cases}$$

where $P_k(\partial T)$ is defined to be the space of restrictions to the boundary ∂T of T of polynomials of degree not greater than k in n variables or equivalently the space of functions in n variables, polynomial of degree not greater than k on each face T' of T and continuous on ∂T . One may easily verify that

$$(5.16) \quad \dim P_k(\partial T) = \begin{cases} \dim P_k(T) & \text{if } k \leq n \\ \\ \dim P_k(T) - \dim P_{k-n-1}(T) & \text{if } k > n. \end{cases}$$

Let $\mathcal{I}_{\partial T_h}^{(k)}$ be the mapping which associates to an element $\psi = (\psi_T)_T \in T_h \in \mathbb{I}$ $\mathbb{C}^0(\partial T)$ the element $(\mathcal{I}_{\partial T}^{(k)} \psi_T)_{T \in T_h}$, where $\mathcal{I}_{\partial T}^{(k)} \psi_T$, the k^{th} order Lagrangian interpolant of ψ_T , is the function coinciding with $\mathcal{I}_{T'}^{(k)}(\psi_{T|T'})$ on each face T' of T. As the elements of $L_h^{(k)}$ are of class C^0 (cf. CIARLET (1978, p.95) for example), for each function $v \in H^1(\Omega) \cap H^2(\partial T_h)$, we have

$$\mathcal{L}_{\partial T}^{(k)}(v_{\mid \partial T}) = (\mathcal{L}_{T}^{(k)}(v_{\mid T}))_{\mid \partial T} \qquad \text{for all } T \in T_{h},$$

or in other words, if we put

$$\psi(\mathbf{v}) = (\psi(\mathbf{v})\mathbf{T})\mathbf{T} \in T_{\mathbf{h}} = (\mathbf{v} \mid \partial \mathbf{T})\mathbf{T} \in T_{\mathbf{h}},$$

then we have

$$(5.18) \mathcal{L}_{\partial T_{h}}^{(k)} \psi_{(v)} = ((\mathcal{L}_{h}^{(k)} v) |_{\partial T})_{T \in T_{h}}.$$

 $\begin{array}{ll} \underline{\textit{THEOREM 5.3}} &. & \textit{For each positive integer k and each triangulation T_h of $\overline{\Omega}$, the mapping $v \to t$_{\partial T_h}^{(k)} \psi(v)$, with $\psi_{(v)}$ given by (5.17)), is a surjection of $H^1(\Omega) \cap H^2(T_h)$ onto $L_{\partial T_h}^{(k)}$. It is a bijection of $L_h^{(k)}$ onto $L_{\partial T_h}^{(k)}$ only in case $k \le n$.} \end{array}$

<u>Proof</u>. In light of Theorem 5.1, it suffices to show that the mapping $v \to \mathcal{L}_h^{(k)} \psi(v)$ is a surjection of $L_h^{(k)}$ onto $L_{\partial T_h}^{(k)}$ and is also a bijection only if $k \le n$. For $k \le n$, the claim follows immediately from the facts that according to (5.16) the spaces $P_k^{(T)}$ and $P_k^{(\partial T)}$ have the same dimension and that for each ψ , $\psi = (\psi_T)_{T \in T_h}$,

belonging to $L_{\partial T_h}^{(k)}$, there exists a unique function v in $L_h^{(k)}$ whose restriction to

belonging to $L_{\partial T_h}^{\text{NT}}$, there exists a unique like the boundary of T is equal to ψ_T for each T in T_h .

For k > n, the dimension of $P_k(T)$ is greater than that of $P_k(\partial T)$, so a bijection from $L_h^{(k)}$ onto $L_{\partial T_h}^{(k)}$ can not exist. To show surjectivity in this case, we have only to remark that given $\psi = (\psi_T)_{T \in T_h} \in L_{\partial T_h}^{(k)}$, the function $v \in L_h^{(k)}$, for example, defined by its values at the lattice points $a \in U$ $\sum_{T=T_h}^{(k)} as$ follows:

$$\mathbf{v(a)} \,=\, \left\{ \begin{array}{ll} \psi_{\mathbf{T}}(\mathbf{a}) & \text{if } \mathbf{a} \in \partial \mathbf{T} \text{ for some } \mathbf{T} \in T_{\mathbf{h}}, \\ \\ 0 & \text{otherwise }, \end{array} \right.$$

defines an antecedent of ψ .

Before announcing an error bound for interpolation in $L_{\partial T_h}^{(k)}$, we would like to define a norm on the space Π $H^{\frac{1}{2}}(\partial T)$ which will be seen to be well adapted to the analysis of dual bybrid methods. the analysis of dual hybrid methods. Thus, given a triangulation $T_{
m h}$ of $\overline{\Omega}$, we put

$$\left\langle \begin{array}{c} \left\| \| \mathbf{v} \| \right\|_{1,T} = \left\{ \left\| \mathbf{v} \right\|_{1,T}^{2} + \mathbf{h}^{-2} \| \mathbf{v} \|_{0,T}^{2} \right\}^{\frac{1}{2}} & \text{for all } \mathbf{v} \in \mathbf{H}^{1}(\mathbf{T}), \ \mathbf{T} \in T_{\mathbf{h}}, \\ \left\| \| \psi \| \right\|_{\frac{1}{2},\partial T} = \inf_{\left\{ \mathbf{v} \in \mathbf{H}^{1}(\mathbf{T}); \mathbf{v} \right\|_{\partial T} = \psi \right\}} \| \| \mathbf{v} \| \|_{1,T} & \text{for all } \psi \in \mathbf{H}^{\frac{1}{2}}(\partial \mathbf{T}), \ \mathbf{T} \in T_{\mathbf{h}}, \\ \left\| \| \psi \| \right\|_{\frac{1}{2},\partial T_{\mathbf{h}}} = \left\{ \sum_{\mathbf{T} \in T_{\mathbf{h}}} \| \psi_{\mathbf{T}} \| \right\|_{\frac{1}{2},\partial T}^{2} & \text{for all } \psi = (\psi_{\mathbf{T}})_{\mathbf{T} \in T_{\mathbf{h}}} \in \Pi \quad \mathbf{H}^{\frac{1}{2}}(\partial \mathbf{T}). \end{aligned}$$

We may now state

Suppose $\{T_h : h \in H\}$ is a regular family of triangulations of Ω . For each positive integer k, there exists a constant C, independent of h, such that if $v \in H^{m+1}(\Omega)$, $1 \le m \le k$, and $\psi = \psi_{(v)}$ is determined by (5.17), then

$$\| | \psi - x_{\partial T_{\mathbf{h}}}^{(\mathbf{k})} \psi | \|_{L_{2}, \partial T_{\mathbf{h}}} \le C h^{\mathbf{m}} | \mathbf{v} |_{\mathbf{m}+1, \Omega}.$$

Proof . With the norms defined in (5.19) we have

$$\left\| \left\| \psi - z_{\partial T_{\mathbf{h}}}^{(\mathbf{k})} \psi \right\| \right\|_{L_{\mathbf{h}}, \partial T_{\mathbf{h}}}^{2} \leq \left\| \left\| \mathbf{v} - \mathcal{I}_{\mathbf{h}}^{(\mathbf{k})} \mathbf{v} \right\|_{1, \Omega}^{2} + \mathbf{h}^{-2} \left\| \left\| \mathbf{v} - \mathcal{I}_{\mathbf{h}}^{(\mathbf{k})} \mathbf{v} \right\|_{0, \Omega}^{2} \right.,$$

so the demonstration is completed on applying Theorem 5.2.

6 - INTERPOLATION OF THE MOMENTS OF A VECTORIAL FUNCTION; AN H(div) APPROXIMATION OF VECTOR VALUED FUNCTIONS

We retain here the assumptions made in Section 5 concerning the domain $\Omega \subset \mathbb{R}^n$, n=2 or 3, and the triangulation T_h , n=0 T. Given a positive integer k, we shall associate to each n-simplex T $\in T_h$ a subspace $D_k(T)$ of the product space $(P_k(T))^n$ and an interpolation operator $E_T^{(k)}$ defined on a subspace $H(\operatorname{div};T)$ of $H(\operatorname{div};T)$ and having values in the space $D_k(T)$. Denoting by $H(\operatorname{div};T_h)$ the space defined by

(6.1)
$$H(\operatorname{div};T_{h}) = (\underline{q} \in (L^{2}(\Omega))^{n}; \forall T \in T_{h}, \underline{q}_{T} \in H(\operatorname{div};T)),$$

we can associate to each vector function $\underline{q} \in \mathit{H}(\mathit{div};T_h)$ the function $E_h^{(k)}\underline{q} \in (L^2(\Omega))^n$ determined by

(6.2)
$$(E_{\mathbf{h}}^{(\mathbf{k})} \underline{\mathbf{q}}) |_{\mathbf{T}} = E_{\mathbf{T}}^{(\mathbf{k})} (\underline{\mathbf{q}}|_{\mathbf{T}})$$
 for all $\mathbf{T} \in T_{\mathbf{h}}$.

The spaces $\mathrm{D_k^{}(T)}$ and the operators $\mathrm{E_T^{\,(k)}}$ will be constructed in such a way that the following two hypotheses are satisfied

(6.3)
$$\operatorname{div} \ \underline{q} = 0 \Longrightarrow \operatorname{div} (E_{T}^{(k)} \underline{q}) = 0 \qquad \text{for all } T \in T_{h},$$
 and for all $\underline{q} \in H(\operatorname{div}; T)$,

and

(6.4)
$$E_{h}^{(k)} \underline{q} \in H(\operatorname{div};\Omega)$$
 for all $\underline{q} \in H(\operatorname{div};\Omega) \cap H(\operatorname{div};T_{h})$.

Property (6.3) is a local property of the interpolation of vector functions in H(div;T). Its analogue for the approximation of scalar functions in $H^1(T)$ is that

the interpolant of a function whose gradient vanishes also has a trivial gradient; i.e., the interpolant of a constant function is a constant. However we did not think it necessary to state such a natural requirement.

Property (6.4) is a continuity condition for the interpolant of a vector function having a certain amount of regularity on $\overline{\Omega}$. Here of course the scalar analogue, that the interpolant $\mathcal{L}_h^{(k)}$ v be in $\mathrm{H}^1(\Omega)$ whenever $\mathbf{v} \in \mathrm{H}^1(\Omega) \cap \mathrm{H}^2(T_h)$, is satisfied as the space $\mathrm{L}_h^{(k)}$ is a subspace of $\mathrm{C}^0(\overline{\Omega})$.

 $\underbrace{\textit{Remark 6.1}}$. Almost reflexively, a mathematician's first idea is to try to reduce the study of interpolation of vector functions to that of scalar functions as given in the preceeding section by choosing

$$D_{k}(T) = (P_{k}(T))^{n},$$

$$H(\text{div},T) = (H^{2}(T))^{n},$$

$$E_{T}^{(k)} = (\mathcal{I}_{T}^{(k)})^{\otimes n}$$

(kth order Lagrangian interpolation on T carried out component by component).

One may check that condition (6.4) is satisfied; however, (6.3) is not. To see this consider in 2 dimensions, on the reference triangle \hat{T} with vertices (0,0), (1,0) and (0,1), the function $q=(q_1,q_2)$ given by

$$\begin{cases} q_1(x_1, x_2) = \frac{1}{2} x_1^2 - x_1 x_2, \\ \\ q_2(x_1, x_2) = \frac{1}{2} x_2^2 - x_1 x_2. \end{cases}$$

This function q is divergence free. Yet its Lagragian interpolant in $(P_1(T))^2$ is none other than the function x ————> ($\frac{1}{2}$ x_1 , $\frac{1}{2}$ x_2) in \hat{T} whose divergence is the constant function x ———> 1.

We shall give two families of examples of H(div) approximation; the first the original examples due to THOMAS (1977, Chapter III), cf. also RAVIART & THOMAS (1977), for the 2 dimensional case and to NEDELEC (1980) for the 3 dimensional case, we present in some detail; the second examples due to BREZZI, DOUGLAS, & MARINI (1985) in the 2-dimensional case, and to BREZZI, DOUGLAS, DURAN, & FORTIN (1987) in the 3-dimensional case will be described more briefly. Yet another family of examples of H(div) approximation may be found in NEDELEC (1986).

For each positive integer k we put

(6.5)
$$D_k = (P_{k-1})^n \oplus x P_{k-1}$$

and define $D_k(T)$ to be the space of restrictions to T of the functions of D_k . In other words, the vector function $\underline{q}=(q_1,\ldots,q_n)$ belongs to $D_k(T)$ if and only if there exists n+1 scalar polynomial functions $q_i^*\in P_{k-1}(T)$, $i=0,1,\ldots,n$, such that

(6.6)
$$q_{i}(x) = q_{i}^{*}(x) + x_{i} q_{0}^{*}(x)$$
 for all $x \in T$, $i=1,...,n$.

The expression (6.6) for q may be made unique by further requiring that q^* be homogeneous of degree k-1. Thus we have

(6.7)
$$\dim P_{k}(T) = (n+1) \dim P_{k-1}(T) - \dim P_{k-2}(T),$$

which together with (5.7) gives us

dim
$$D_k(T) = (n+k) \frac{(n+k-2)!}{(n-1)!(k-1)!}$$
,

or more specifically

$$\dim D_{k}(T) = \begin{cases} k(k+2), & \text{when } n=2, \\ k(k+1), & \text{when } n=3. \end{cases}$$

For each nonnegative integer ℓ we denote by $M_{\ell}(T)$ the space of linear forms on $(C^0(T))^n$ defined by

(6.8)
$$M_{\ell}(T) = \{(q \to \int_{T} q.r. dx) \in \mathcal{L}((C^{0}(T))^{n}, R) ; r \in (P_{\ell})^{n}\},$$

i.e. an element of $M_{\ell}(T)$ assigns to the vector function \mathbf{q} on T one of its moments on T of order no greater than ℓ . In keeping with classical notation we shall write $M_{\ell}(T) = \{0\}$ for ℓ a negative integer, in particular for ℓ --1. Similarly for each nonnegative integer ℓ and each face T' of T we denote by $M_{\ell}(T')$ the space of linear forms on $(C^0(T))^n$ defined by

(6.9)
$$M_{\ell}(T') = \{(q + \int_{T}, q \cdot \nu_{T} w \, d\sigma) \in \mathcal{L}((C^{0}(T))^{n}, R) ; w \in P_{\ell}\},$$

where $\nu_{\rm T}$ is the unit exterior normal to ∂T , constant on each face T' of T. Thus an element of $M_{\ell}({\rm T}')$ assigns to the vector function q on T one of the moments on T' of order no greater than ℓ of its exterieur normal trace $q.\nu_{\rm T}$. Now we define, for k a positive integer,

(6.10)
$$M_{T}^{(k)} - (\bigcup_{T' < T} M_{k-1}(T')) \cup M_{k-2}(T),$$

where T'<T mean T' is an (n-1)-dimensional face of T. Without change of notation we shall consider henceforth that the forms in $M_{\rm T}^{(k)}$ are defined on the space $H({\rm div};T)$,

(6.11)
$$H(\operatorname{div};T) = \{ \underline{q} \in H(\operatorname{div};T) : \underline{q} \cdot \underline{\nu}_T \in L^2(\partial T) \}.$$

The dimension of the space $M_{T}^{(k)}$ is given by

$$\dim M_T^{(k)} = (n+1) \dim P_{k-1}(\mathbb{R}^{n-1}) + n \dim P_{k-2}(T),$$

or more explicitly, by

(6.12)
$$\dim M_{T}^{(k)} = (n+k) \frac{(n+k-2)!}{(n-1)!(k-1)!}.$$

Thus we see that

(6.13)
$$\dim D_k(T) = \dim M_T^{(k)}, \qquad k=1, 2,$$

THEOREM 6.1. Let T be an n simplex and k a positive integer. For each $q \in \mathit{H}(\text{div};T)$ there exists one and only one function, which we shall denote $E^{(k)}_{T}q$, in $D_{k}(T)$ for which

(6.14)
$$m(E_T^{(k)}\underline{q}) = m(\underline{q}) \qquad \text{for all } m \in M_T^{(k)}.$$

The mapping $E_T^{(k)}: H(\text{div},T) \to D_k(T)$ thus defined further satisfies the commuting diagram property

$$H(\operatorname{div};T) \xrightarrow{\operatorname{div}} L^{2}(T)$$

$$\downarrow E_{T}^{(k)} \qquad \downarrow \Pi_{T}^{(k-1)}$$

$$D_{k}(T) \xrightarrow{\operatorname{div}} P_{k-1}(T)$$

i.e.

(6.15)
$$\operatorname{div}(E_{T}^{(k)}q) = \Pi_{T}^{(k-1)} (\operatorname{div} q) \qquad \text{for all } q \in H(\operatorname{div};T),$$

where ${\rm I\!I}_{\,\,T}^{\,\,(k-1)}$ denotes the orthogonal projection in $L^2({\rm T})$ onto $P_{k-1}({\rm T})$.

<u>Proof</u>. Suppose for the moment that the operator $E_T^{(k)}$ from H(div;T) to $D_k^{(T)}$ satisfying (6.14) is well defined and let $q \in H(\text{div};T)$. Then it is clear that the divergence of any element of $D_k^{(T)}$, and in particular of $E_T^{(k)}q$, is a polynomial of degree not exceeding k-1. Thus using Green's formula with (6.14) we have

$$\int_{T} \operatorname{div} \left(E_{T}^{(k)} q \right) w \, dx = \int_{T} \operatorname{div} q w \, dx \qquad \text{for all } w \in P_{k-1}(T)$$

or equivalently (6.15).

To show that given $q \in \mathit{H}(\operatorname{div};T)$, (6.14) does—indeed determine a unique element of $D_k(T)$, in view of (6.13) it suffices to show that if $p \in D_k(T)$ and m(p)=0 for each $m \in \mathit{M}_T^{(k)}$, then p must be trivial. It is easy to check that the normal trace on any face T' of T of any element of $D_k(T)$ is a polynomial of degree at most k-1. Thus, the normal trace of p on T' must be trivial. The divergence of p, a priori a polynomial of degree at most k-1 on T, can then be shown to vanish by using Green's formula. It follows that p belongs to the subspace $(P_{k-1}(T))^n$ of $D_k(T)$; to see this simply apply Euler's identity for homogeneous functions: grad $f(x) \cdot x = \alpha f(x)$ whenever f is homogeneous of degree α . Next we show that if p is a vector normal to a face T' of T that $p \cdot p'$ vanishes on T. As $p \cdot p'$ vanishes on T', we have that $p \cdot p' = bw$, where p is the first degree polynomial whose zeroes form the (n-1)-dimensional hyperplane containing T', and p is p in p and we have

$$\int_{T} \tilde{p} \cdot (w \nu') dx = \int_{T} bww dx = 0.$$

As b is either always positive or always negative on T, w and hence $p \cdot \nu$ vanish on T. The theorem now follows on noting that for any n of the n+1 faces of T, the corresponding normal vectors form a basis for \mathbb{R}^n .

For each $T \in T_h$ we have now constructed an interpolation operator $E_T^{(k)}$ from $H(\operatorname{div};T)$ into $D_k(T) \subset (P_k(T))^n$ satisfying (6.15) and a fortiori (6.3). To a triangulation T_h of $\overline{\Omega}$ we associate the space $H(\operatorname{div};T_h)$ defined by (6.1) together with (6.11):

(6.16)
$$H(\operatorname{div};T_{h}) = \{ \underline{q} \in (L^{2}(\Omega))^{n} ; \forall T \in T_{h}, \\ \underline{q}_{T} \in H(\operatorname{div};T) \text{ and } \underline{q}_{T} \cdot \underline{\nu}_{T} \in L^{2}(\partial T) \},$$

and the equilibrium interpolation operator $E_h^{(k)}$ of order k on $H(\text{div}; T_h)$ constructed from the operators $E_T^{(k)}$ with the aid of (6.2).

To see that $E_h^{(k)}$ satisfies (6.4), we note that if q is sufficiently regular in the sense that $q \in H(\operatorname{div};\Omega) \cap H(\operatorname{div};T_h)$ then Theorem 1.3 implies

(6.17)
$$q_{|T_1} \cdot \nu_{T_1} + q_{|T_2} \cdot \nu_{T_2} = 0$$
 on T', for all $T' = T_1 \cap T_2$ with $T_1, T_2 \in T_h$.

Thus, for each $T' = T_1 \cap T_2$ with $T_1, T_2 \in T_h, r = E_h^{(k)}$ g satisfies

$$\int_{T_{1}} (r_{1} \cdot \nu_{T_{1}} + r_{1} \cdot \nu_{T_{2}} \cdot \nu_{T_{2}}) w ds = 0 for all w \in P_{k-1}.$$

The restrictions $r_{\mid T} \in D_k(T)$, $T \in T_h$, are polynomial of degree at most k and their normal traces on ∂T are piecewise polynomial of degree at most k-1; indeed, on each face T' of T, $T \in T_h$, $x \cdot \nu_T$ is constant. With this observation of central importance we see that r satisfies the continuity condition (6.17) and hence $r = E_h(x) = E_h(x)$ i.e. (6.4) holds.

We denote by $\mathbf{E}_{h}^{(k)}$ the finite dimensional space of equilibrium interpolants

(6.18)
$$\mathbb{E}_{h}^{(k)} = \{ \underline{q} \in \mathbb{H}(\operatorname{div};\Omega) ; \forall T \in T_{h}, \underline{q}_{T} \in \underline{D}_{k}(T) \},$$

which, as we have just seen, is the space of functions $\underline{q} \in (L^2(\Omega))^n$ for which the restriction \underline{q} may be identified with a function in $D_k(T)$ for each $T \in T_h$ and for which (6.17) holds. In analogy with (5.12), $E_h^{(k)}$ may also be characterized as a space of fixed points as follows:

(6.18 bis)
$$E_h^{(k)} = \{q \in H(div;\Omega) \cap H(div;T_h) ; E_h^{(k)}(q) = q\}.$$

Thus we have

THEOREM 6.2 . For each positive integer k, and each triangulation T_h of $\overline{\Omega}$, the mapping $E_h^{(k)}$ is a projection from $H(\text{div};\Omega) \cap H(\text{div};T_h)$ onto $E_h^{(k)}$.

The next theorem gives a bound for the error of interpolation by $E_h^{(k)}$:

THEOREM 6.3. Suppose $\{T_h : h \in H\}$ is a regular family of triangulations of $\overline{\Omega}$ and k is a positive integer. There exists a constant C, independent of h, such that

and

(6.20)
$$\|\operatorname{div}\left(\underline{q}-E_{h}^{(k)}\underline{q}\right)\|_{0,\Omega} \leq \operatorname{Ch}^{\ell}\operatorname{div}\underline{q}|_{\ell,\Omega} \quad \text{for all } \underline{q} \in \left(\operatorname{H}^{1}(\Omega)\right)^{n} \text{ with } \\ \operatorname{div}\underline{q} \in \operatorname{H}^{\ell}(\Omega), \ 0 \leq \ell \leq k.$$

<u>Proof</u>. Estimate (6.20) is a classical corollary of property (6.15) identifying div $(E_T^{(k)} \neq 1)$ as the image under the orthogonal L²-projection onto $L_h^{(k)}$ of div q, and we shall not give the details of its demonstration. The argument justifying (6.19) is more technical: for each $T \in T_h$ we must find a bound for the quantity $\|q - E_T^{(k)} \neq 1\|_{0,T} = \sum_{i=1}^n \|q_i - (E_T^{(k)} \neq 1)\|_{0,T} \|_{0,T}^2 \|_{1,T}^2$; however, this can not be done component by component as the interpolation space is not a simple tensor product of n copies of the same space of polynomials. Thus it is necessary to work directly with the subspaces of vector functions of n variables with values in \mathbb{R}^n . We employ the technique of reducing the question of the error bound on T to that on the reference triangle \hat{T} as in CIARLET & RAVIART (1972). For each $T \in T_h$, there exists an affine transformation $F: \hat{T} \to T = F(\hat{T})$, bijective from \hat{T} onto T. We denote by DF the linear tangent mapping and by J the Jacobian of the transformation. We have here $|T| = (\text{mes } T)/(\text{mes } \hat{T})$. To each scalar function $\hat{v} \in H^1(\hat{T})$, classically is associated the function $v \in H^1(T)$, defined by

$$(6.21) v(x) = \hat{v}(\hat{x}) for each x = F(\hat{x}) \in T.$$

Here, to each vector function $\hat{q} \in H(\text{div}; \hat{T})$ we associate the function q, that we shall show belongs to H(div; T), defined by

(6.22)
$$q(x) = \frac{1}{|J|} DF \hat{q}(\hat{x}) \qquad \text{for each } x = F(\hat{x}) \in T.$$

The transformation (6.22) was constructed in such a way that

(6.23)
$$\int_{T} q \cdot \operatorname{grad} v \, dx = \int_{\hat{T}} \hat{q} \cdot \operatorname{grad} \hat{v} \, d\hat{x}.$$

Using Green's formula we deduce, cf. THOMAS (1976, Lemma 5.1), that

(6.24)
$$\int_{T} v \operatorname{div} q \, dx = \int_{\hat{T}} \hat{v} \operatorname{div} \hat{q} \, d\hat{x},$$

and we have, with $\hat{\nu}$ denoting of course the exterieur unit normal of $\partial \hat{T}$,

(6.25)
$$\int_{\partial T} \mathbf{v} \, \mathbf{q} \cdot \mathbf{\nu} \, d\sigma = \int_{\partial \hat{T}} \mathbf{v} \, \hat{\mathbf{q}} \cdot \hat{\mathbf{\nu}} \, d\hat{\sigma},$$

for all sufficiently regular functions v and q. Thus $q \longrightarrow q$ is a bijective affine transformation from $H(div; \hat{T})$ onto H(div; T) such that

(6.26) (mes T) div
$$q = (\text{mes } \hat{T})$$
 div \hat{q}

and on each face $T' = F(\hat{T}')$ of T

6.27) (mes T')
$$\mathbf{q} \cdot \mathbf{v} = (\text{mes } \hat{\mathbf{T}}') \hat{\mathbf{q}} \cdot \hat{\mathbf{v}}.$$

Furthermore, it is immediately verified that this transformation takes $\mathbf{D}_k(\hat{\mathbf{T}})$ onto $\mathbf{D}_k(\mathbf{T})$:

$$(6.28) q \in D_{k}(T) \Leftrightarrow \tilde{q} \in D_{k}(\hat{T}).$$

We conclude from (6.27) and (6.28) that the interpolant in $D_k(T)$ of a function q, image of the function \hat{q} defined on \hat{T} , is none other than the image of the interpolant in $D_k(\hat{T})$ of \hat{q} :

(6.29)
$$(E_{\mathbf{T}}^{(\mathbf{k})} \mathbf{q}) \cdot \mathbf{F} = \frac{1}{|\mathbf{J}|} \mathrm{DF}(E^{(\mathbf{k})} \mathbf{\hat{q}}).$$

This property allows us to reduce the problem of obtaining an error bound for interpolation on T to that of obtaining an error bound for interpolation on \hat{T} , the reference n-simplex. As the subspace $(P_{k-1})^n$ of D_k is invariant under the interpolation operator, we know that there exists a constant \hat{C} depending only on \hat{T} and k such that.

$$\|\hat{\mathbf{q}} - E_{\hat{\mathbf{T}}}^{(k)}\hat{\mathbf{q}}\|_{O,\hat{\mathbf{T}}} \le \hat{\mathbf{C}} \|\hat{\mathbf{q}}\|_{\ell,\hat{\mathbf{T}}} \qquad \text{for all } \hat{\mathbf{q}} \in (\mathbf{H}^{\ell}(\hat{\mathbf{T}}))^n \text{ with } 1 \le \ell \le k.$$

With this bound and change of variable formulas we obtain for each T in T_h a local bound for $\|q - E_T^{(k)}q\|_{0,T}$ and then the global bound (6.19).

Remark 6.2 . The introduction of the spaces D_k is fundamental for our construction of an $H(\text{div};\Omega)$ approximation. We have the strict inclusions $(P_{k-1})^n \subset D_k \subset (P_k)^n$, and we have chosen to index by k the space D_k strictly between $(P_{k-1})^n$ and $(P_k)^n$ though in the literature one also finds this same space indexed by k-1. Our choice was made with consideration for the error bounds. As Lagrangian

interpolation, $x_h^{(k)}$, by functions locally in P_k leads to error bounds in $H^1(\Omega)$ of order h^k , so the interpolation $E_h^{(k)}$, with E as in equilibrium, by functions locally in D_k leads to an error in $H(\operatorname{div};\Omega)$ of order h^k .

The transformation $\hat{q} \longrightarrow q$ introduced during the demonstration of Theorem 6.3 in (6.22) is none other than the transformation known in mechanics as the Piola transformation of a vector field defined on a manifold of R^n , cf. for example MARSDEN & HUGHES (1983, pp 116-118).

Next we present a few results concerning the approximation of the normal traces on the set $\partial T_{\rm h} = \bigcup$ ∂T of sufficiently regular vector functions q, for $T \in T_{\rm h}$ example $q \in (H^1(\Omega))^n$. Given a triangulation T_h of $\overline{\Omega}$ and a positive integer k, we denote by $E_{\partial T_{h}}^{(k)}$ the space defined by

$$(6.30) E_{\partial T_{h}}^{(k)} = \{\mu = (\mu_{T})_{T \in T_{h}} \in \prod_{T \in T_{h}} L^{2}(\partial T) ; \exists q \in E_{h}^{(k)} \text{ s.t. } \forall T \in T_{h},$$

$$q_{|\partial T} \cdot \nu_{T} = \mu_{T}\},$$

where E $_{
m h}^{(k)}$ is the space defined by (6.18). The space E $_{
m \partial T}^{(k)}$ may also be characterized as the space of functions μ = $(\mu_{\rm T})_{{\rm T}\in T_{\rm h}}$ satisfying the following two properties :

$$(6.31) \begin{cases} \text{i) } \mu_{\mathrm{T}} \in \mathrm{D}_{\mathrm{k}}(\partial \mathtt{T}), & \text{for each } \mathtt{T} \in T_{\mathrm{h}} \\ \\ \text{and} & \\ \\ \text{ii) } \mu_{\mathrm{T}_{1}} + \mu_{\mathrm{T}_{2}} = 0 & \text{on each } \mathtt{T}', \mathtt{T}' = \mathtt{T}_{1} \, \cap \, \mathtt{T}_{2} \, \text{with} \, \mathtt{T}_{1}, \mathtt{T}_{2} \in T_{\mathrm{h}}, \end{cases}$$

where by $D_{1}(\partial T)$ we mean the space

(6.32)
$$D_{k}(\partial T) = \{ \mu \in L^{2}(\partial T) ; \exists q \in D_{k}(T) , q \cdot \nu_{T} = \mu \},$$

or equivalently the space of functions polynomial of degree at most (k-1) on each face T' of T with no continuity constraint at the intersection of two faces. Thus we have

(6.33)
$$\dim D_{k}(\partial T) = (n+1) \frac{(n+k-2)!}{(n-1)!(k-1)!}.$$

Let $E_{\partial T_h}^{(k)}$ be the operator which associates to each element $\mu = (\mu_T)_{T \in T_h} \in \Pi$ $L^2(\partial T)$ the element $(E_{\partial T}^{(k)})_{T \in T_h}$ where $E_{\partial T}^{(k)}$ is the orthogonal projection from $L^2(\partial T)$ onto $D_k(\partial T)$, $T \in T_h$. Associating to each vector function $q \in H(\operatorname{div}; T_h) \cap H(\operatorname{div}; \Omega)$ the element

(6.34)
$$\mu_{(q)} = (q \cdot \nu_T)_{T \in T_h}$$

we clearly have

$$(6.35) E_{\partial T_h}^{(k)} \mu_{(\underline{q})} = ((E_h^{(k)} \underline{q}) \cdot \nu_{\underline{T}})_{\underline{T} \in T_h} \in E_{\partial T_h}^{(k)}.$$

Hence we deduce :

Here we shall give only one result concerning error bounds for interpolation in $E_{\partial T_h}^{(k)}$. However, beforehand, in a manner analogous to (5.19), we define several norms that we shall find useful. Given a triangulation T_h of $\overline{\Omega}$, we put

$$\left\{ \begin{array}{l} \left\| \left\| \mathbf{q} \right\| \right\|_{\mathbf{H}(\operatorname{div}; \mathbf{T})} = \left(\left\| \mathbf{q} \right\|_{0, \mathbf{T}}^{2} + \mathbf{h}^{2} \right\| \operatorname{div} \mathbf{q} \left\|_{0, \mathbf{T}}^{2} \right)^{\frac{1}{2}} \text{ for all } \mathbf{q} \in \mathbf{H}(\operatorname{div}; \mathbf{T}), \mathbf{T} \in T_{\mathbf{h}}, \\ \left\| \left\| \boldsymbol{\mu} \right\| \right\|_{-\frac{1}{2}, \partial \mathbf{T}} = \inf_{\left\{ \mathbf{q} \in \mathbf{H}(\operatorname{div}; \mathbf{T}); \mathbf{q} \cdot \boldsymbol{\nu}_{\mathbf{T}} = \boldsymbol{\mu} \right\}} \left\| \left\| \mathbf{q} \right\| \right\|_{\mathbf{H}(\operatorname{div}; \mathbf{T})} \text{ for all } \boldsymbol{\mu} \in \mathbf{H}^{-\frac{1}{2}}(\partial \mathbf{T}), \mathbf{T} \in T_{\mathbf{h}}, \\ \left\| \left\| \boldsymbol{\mu} \right\| \right\|_{-\frac{1}{2}, \partial T_{\mathbf{h}}} = \left\{ \sum_{\mathbf{T} \in T_{\mathbf{h}}} \left\| \boldsymbol{\mu}_{\mathbf{T}} \right\|_{-\frac{1}{2}, \partial \mathbf{T}}^{2} \right\}^{\frac{1}{2}} \text{ for all } \boldsymbol{\mu} = (\boldsymbol{\mu}_{\mathbf{T}}) \in \prod_{\mathbf{T} \in T_{\mathbf{h}}} \mathbf{H}^{-\frac{1}{2}}(\partial \mathbf{T}). \end{aligned}$$

The following result will be needed for our analysis of primal hybrid methods :

THEOREM 6.5 Suppose $\{T_h : h \in H\}$ is a regular family of triangulations of $\overline{\Omega}$. For each positive integrer k, there exists a constant C, independent of h, such that if $q \in (H^l(\Omega))^n$, $1 \le l \le k$, and $\mu = \mu(q)$ is determined by (6.34), then

(6.37)
$$\left| \left| \left| q - E_{\partial T_{\mathbf{h}}}^{(\mathbf{k})} \mu \right| \right| \right|_{-\frac{1}{2}, \partial T_{\mathbf{h}}} \leq \operatorname{Ch}^{\ell} \left| q \right|_{\ell, \Omega} .$$

To obtain this result we employ the bound (6.19) of order ℓ and the bound (6.20) of order (ℓ -1).

The second family of examples, is defined for integers k, $k \ge 2$. One chooses for interpolation space on the simplex T the space

(6.38)
$$D_{k}^{*}(T) = (P_{k-1}(T))^{n},$$

where we have kept the indexing convention of Remark 6.2 instead of that of BREZZI, DOUGLAS, & MARINI (1985). For $q \in \mathit{H}(\operatorname{div};T)$ the interpolant $\mathit{E}_{T}^{*(k)}q$ is defined to be the unique element of $\mathit{D}_{k}^{*}(T)$ such that

$$m(E_T^{*(k)}\underline{q}) = m(\underline{q})$$
 for all $m \in M_T^{*(k)}$,

where the set of moments $M_{\mathrm{T}}^{*(k)}$ used to define the interpolant is the union

$$\textit{M}_{\texttt{T}}^{*(k)} = (\bigcup_{\texttt{T}' < \texttt{T}} \textit{M}_{(k-1)}(\texttt{T}')) \cup \textit{M} \underset{k-2}{\texttt{grad}} (\texttt{T}) \cup \textit{M}_{k-1}^{\perp}(\texttt{T}),$$

with $M_{\ell}(T')$ defined by (6.9),

$$M_{\ell}^{grad}$$
 (T) = {($q \longrightarrow \int_{T} q \cdot grad \ w \ dx$) $\in \mathcal{L}((C^{0}(T))^{n}, \mathbb{R})$; $w \in P_{\ell}$ }

and

$$M_{\ell}^{\perp}$$
 (T) = { (q -----> $\int_{T} q \cdot w \, dx \in \ell((C^{0}(T))^{n}, \mathbb{R}) ; w \in (P_{\ell})^{n},$
 $w \cdot \nu_{T} = 0 \text{ on } \partial T, \text{ and } div w = 0$ }.

Define the operator $E_h^{\star(k)}$ by (6.2) and the space $E_h^{\star(k)}$ by analogy with (6.18) or equivalently (6.18 bis). It can be shown that (6.15) still holds with $E_T^{(k)}$ replaced by $E_T^{\star(k)}$ and $\Pi^{(k-1)}_T$ replaced by $\Pi^{(k-2)}_T$, and thus so do (6.3) and (6.4) if $E_h^{(k)}$ is also replaced by $E_h^{\star(k)}$. In fact the analogues of Theorem 6.1 and Theorem 6.2 obtained by replacing $E_T^{(k)}$, $E_h^{(k)}$, $D_k^{(T)}$, and $E_h^{(k)}$ by $E_T^{\star(k)}$, $E_h^{\star(k)}$, $D_k^{\star(T)}$ and $E_h^{\star(k)}$ respectively remain valid. One can also demonstrate the following analogue of Theorem 6.3: if $\{T_h$; $h \in H\}$ is a regular family of triangulations of Π , then

and

Remark 6.4 . The space $E_h^{*(k)}$,

(6.41)
$$E_h^{\star(k)} = \{\underline{q} \in H(\operatorname{div};\Omega) \cap H(\operatorname{div}; T_h) ; \underline{q} = E_h^{\star(k)}\underline{q}\} ,$$

is clearly a proper subspace of $E_h^{(k)}$. Hence fewer degrees of freedom are needed to determine the interpolants in $E_h^{*(k)}$ than in $E_h^{(k)}$. The interpolation errors in these two approximation spaces are of the same order when measured in the L^2 -norm $\| \|_{0,\Omega}$; however, interpolation in $E_h^{*(k)}$ is less precise by one order than that in $E_h^{(k)}$ when measured in the $H(\operatorname{div};\Omega)$ norm unless $|\operatorname{div} q|_{k-1,\Omega} = 0$, i.e. unless $\operatorname{div} q \in P_{k-2}$. Moreover one may verify that

$$D_k^*(T) = \{ q \in D_k(T) ; \text{ div } q \in P_{k-2}(T) \}.$$

Remark 6.5 . In fact one of the most important distinctions between the approximations in $E_h^{(k)}$ and in $E_h^{*(k)}$ lies not in the interpolations themselves but in the fact that for the dual mixed method for the elliptic problem the $E_h^{(k)}$ interpolation for the vector variable is used in conjuntion with a locally P_{k-1} interpolation for the scalar variable whereas when interpolation in $E_h^{*(k)}$ is used for the vector variable the natural domain for approximating the scalar variable is the space of locally P_{k-2} functions. Thus, for example, when a piecewise constant approximation is sufficient for the scalar variable but more precision than that provided by the $E_h^{(1)}$ interpolant is desired for the vector variable, one has the option of approximating the vector function in $E_h^{*(2)}$ while retaining the locally P_0 approximation of the scalar function instead of using the full P_1 , $E_h^{(2)}$ approximation. The gain in accuracy thus obtained, however, is, a priori, only in the $(L^2(\Omega))^n$ norm.

7 - EXAMPLES OF FINITE ELEMENT TYPE APPROXIMATION FOR NON TRIANGULAR GEOMETRIC FORMS

Using the families of examples of finite element type approximation developed in Sections 5 and 6, we can construct approximations of order k, for each positive integer k, of functions in $H^1(\Omega)$ and of functions in $H(\operatorname{div};\Omega)$ in the academic setting Ω a polyhedral domain triangulated by n-simplexes. To simplify vocabulary we fix n=2. Thus for the preceeding theory Ω was supposed to be an open, polygonal domain triangulated by triangles. Here we shall indicate a few generalizations related to the geometry of the triangulation : rectangular finite elements, curved finite elements such as curvilinear triangles, convex quadrilaterals...

We begin with the case of rectangular finite elements. We shall describe families of examples of approximation of functions in $\mathrm{H}^1(\Omega)$ and of functions in $\mathrm{H}(\mathrm{div};\Omega)$ when Ω is an open set of \mathbb{R}^2 "triangulated" in the form $\overline{\Omega}=\cup$ T by rectangles T having sides parallel to the axes. These families will be the rectangular analogues of the triangular families defined in the two preceding sections.

We first consider the approximation of H^1 functions. For any two nonnegative integers k and ℓ , denote by $\mathrm{P}_{k,\ell}$ the space of polynomials in two variables, of degree at most k in the first variable and of degree at most ℓ in the second variable, and by $\mathrm{P}_{k,\ell}(\mathrm{T})$ the space of restrictions to the rectangle T of the polynomials in $\mathrm{P}_{k,\ell}$. To a rectangle T and a positive integer k is associated the set $\sum_T^{(k)}$ of points of the grid on T obtained by subdividing each edge of T into k equal parts. The Lagrangian interpolant $\mathcal{L}_T^{(k)}v$ on T of order k of a function $v\in C^0(\mathrm{T})$ is the unique function in $\mathrm{P}_{k,k}(\mathrm{T})$ which coincides with v at each point of $\sum_T^{(k)}$. For T_h a triangulation of $\overline{\Omega}$ by rectangles, we define the interpolation operator $\mathcal{L}_h^{(k)}$ from $\mathrm{H}^1(\Omega)\cap\mathrm{H}^2(T_h)$ onto $\mathrm{L}_h^{(k)}$, where $\mathrm{L}_h^{(k)}$ is defined by

(7.1)
$$L_{h}^{(k)} = \{ v \in H^{1}(\Omega) ; \forall T \in T_{h}, v_{T} \in P_{k,k}(T) \},$$

by requiring that $\mathcal{L}_h^{(k)}$ v agree on T with $\mathcal{L}_T^{(k)}(v_T)$ for each $T \in T_h$. The error estimate for Lagrangian interpolation given by Theorem 5.2 remain valid when the family of triangulations is a regular family of triangulations by rectangles; i.e. if $\{T_h; h \in \mathcal{H}\}$ is a family of triangulations T_h of $\overline{\Omega}$ by rectangles T of diameter $h_T \leq h$ which is regular in the sense that the ratio of the width to the length of each rectangle T, $T \in T_h$, $h \in \mathcal{H}$, is bounded below by a positive constant independent of h, then there exists a constant C such that

(7.2)
$$\|\mathbf{v} - \mathbf{z}_{h}^{(k)}\mathbf{v}\|_{1,\Omega} + h^{-1} \|\mathbf{v} - \mathbf{z}_{h}^{(k)}\mathbf{v}\|_{0,\Omega} \le Ch^{m} \|\mathbf{v}\|_{m+1,\Omega}$$

for all $v \in H^{m+1}(\Omega)$ with $1 \le m \le k$.

The demonstration for this case may be found in RAVIART & THOMAS (1983, pp. 79-103), for example.

For the approximation of vector functions on a rectangle T we define the space $D_k(T)$ to be the product space $P_{k,k-1}(T) \times P_{k-1,k}(T)$. The equilibrium interpolant $E_T^{(k)}q$ of order k, k \geq 1, on T of the vector function $q \in \mathit{H}(\mathrm{div};T)$ is the unique function in $D_k(T)$ such that

$$(7.3) \qquad \int_{\mathbb{T}'} (E_{\mathbb{T}}^{(k)} q) \cdot \nu_{\mathbb{T}} w \, d\sigma = \int_{\mathbb{T}'} q \cdot \nu_{\mathbb{T}} w \, d\sigma \qquad \text{for all } w \in P_{k-1}(\mathbb{T}') \text{ with } \mathbb{T}' < \mathbb{T},$$

and

For T_h a triangulation of $\overline{\Omega}$ by rectangles we define the interpolation operator $E_h^{(k)}$ from $H(\operatorname{div};\Omega) \cap H(\operatorname{div};T_h)$ onto $E_h^{(k)}$, where $E_h^{(k)}$ is defined by

(7.5)
$$E_{h}^{(k)} = \{ \underline{q} \in H(\operatorname{div};\Omega) ; \forall T \in T_{h}, \underline{q}|_{T} \in D_{k}(T) \},$$

by requiring that $E_h^{(k)}$ q agree on T with $E_T^{(k)}$ (q) for each $T \in T_h$. Using Green's formula one can show that on each rectangle T

(7.6)
$$\operatorname{div}(E_{T}^{(k)}q) = \Pi_{T}^{(k-1)}(\operatorname{div}q) \qquad \text{for all } q \in H(\operatorname{div};T),$$

where $\Pi_T^{(k-1)}$ is the orthogonal L^2 projection of $L^2(T)$ onto $P_{k-1,k-1}(T)$. The error estimate given by Theorem 6 remains valid in the rectangular case : if $\{T_h^-; h\in H\}$ is a regular family of triangulations of $\overline{\Omega}$ by rectangles then there is a constant C such that

and

For the demonstration in the case n=2 treated here, see THOMAS (1977, Chapter III); for the generalization to the case n=3, see NEDELEC (1980).

The analogue for rectangular elements of the second family described in Section 6 is obtained by choosing for $k\geq 2$ for interpolation space on the rectangle T the space

(7.9)
$$D_{k}^{*}(T) = (P_{k-1}(T))^{2} + (rot x_{1}^{k} x_{2}) + (rot x_{1}^{k} x_{2}^{k})$$

and for interpolation operator $E_T^{*(k)}$ on H(div;T) the mapping associating to $q \in H(\text{div};T)$ the unique function in $D_k^*(T)$ such that

(7.10)
$$\int_{T'} (E_T^{*(k)} q) \cdot \nu_T w d\sigma = \int_{T'} q \cdot \nu_T w d\sigma \qquad \text{for all } w \in P_{k-1}(T')$$
 with $T' < T$.

and

(7.11)
$$\int_{\mathbb{T}} (E_{\mathbb{T}}^{*(k)} q) \cdot \mathbf{r} \, d\mathbf{x} = \int_{\mathbb{T}} q \cdot \mathbf{r} \, d\mathbf{x} \qquad \text{for all } \mathbf{r} \in (P_{k-2}(\mathbb{T}))^{2}.$$

Then to a triangulation T_h of $\overline{\Omega}$ by rectangles we associate the interpolation space $E_h^{*(k)}$ and interpolation operator $E_h^{*(k)}$ from $H(\operatorname{div};\Omega) \cap H(\operatorname{div};T_h)$ onto $E_h^{*(k)}$ defined by analogy with the preceding example. The analogue of (7.6) holds in this case as do the error bounds (6.39) and (6.40). For the details of the construction and the demonstration of the error bounds see BREZZI, DOUGLAS, & MARINI (1985, paragraph 5). The extension to the 3 dimensional case is given in BREZZI, DOUGLAS, DURAN, & FORTIN (1987, paragraph 3).

<u>Remark 7.1</u>. In the rectangular case the difference in size of the spaces $D_k(T)$ and $D_k^*(T)$ is more striking than in the triangular case. For T a rectangle,

$$\begin{cases} \dim D_k(T) = 2k^2 + 2k, \\ \dim D_k^*(T) = k^2 + k + 2, \end{cases}$$

and indeed, if NE is the number of elements in $T_{\rm h}$,

dim
$$E_h^{(k)}$$
 - dim $E_h^{*(k)}$ - $(k^2 + k-2)$ NE -
10 NE if k=3.

However, at the level of construction of the spaces $E_h^{(k)}$ and $E_h^{\star(k)}$, we note that the degrees of freedom which assure that the continuity condition (6.17) holds, i.e. those which assure the connections between the finite elements and thus contribute to the size of the algebraic system to be solved, are exactly the same for the two methods.

The situation is not the same in 3 dimensions. If $D_k(T)$ and $D_k^*(T)$ denote the local approximation spaces on a rectangular solid T as defined in NEDELEC (1980) and in BREZZI, DOUGLAS, DURAN & FORTIN (1987) respectively, we have

$$\begin{cases} \dim D_{k}(T) = 3k^{2}(k+1), \\ \dim D_{k}^{*}(T) = \frac{1}{2} k(k^{2}+3k+8). \end{cases}$$

Given a triangulation T_h of $\overline{\Omega}$ by rectangular solids, the degrees of freedom needed to guarantee the required continuity in the corresponding finite dimensional subspaces of $H(\operatorname{div};\Omega)$, $E_h^{(k)}$ and $E_h^{*(k)}$, are the moments of the normal component with respect to

$$\begin{cases} P_{k-1,k-1}(T') & \text{(i.e. } k^2 \text{ moments) for elements of } E_h^{(k)}, \\ P_{k-1}(T') & \text{(i.e. } \frac{1}{2} k(k+1) \text{ moments) for elements of } E_h^{\star(k)}. \end{cases}$$

Here the gain in calculation time with interpolation by $E_h^{*(k)}$, the operator associated with $E_h^{*(k)}$, instead of $E_h^{(k)}$, the operator associated with $E_h^{(k)}$, is clear. If NF is the number of 2-dimensional faces T' of elements T of T_h with T' $\subset \Omega$, then the difference in calculation time can be estimated in terms of the number

$$(3k^2-3k)NF = \begin{cases} 6NF, & \text{if } k=2, \\ \\ 18NF, & \text{if } k=3. \end{cases}$$

Nevertheless, concerning the precision of the interpolation in 3 dimensions as well as in 2 dimensions, we note as in the simplicial case (cf. Remark 6.5) that while the order of approximation is the same for $E_h^{(k)}$ and $E_h^{*(k)}$ in $(L^2(\Omega))^n$, the accuracy in $H(\text{div};\Omega)$ of $E_h^{*(k)}$ is of one order less than that of $E_h^{(k)}$.

Remark 7.2. Another family of examples of H(div) approximation on rectangles in 2 dimensions and in 3 dimensions may be found in BREZZI, DOUGLAS, FORTIN, & MARINI (1987). For rectangular elements in 3 dimensions as well as for prisms, other families are given in NEDELEC (1986).

We cannot develop here a theory of approximation as general as that constructed above for finite elements of such an elementary form as a triangle or a rectangle for the case of finite elements of other geometric forms. However, we shall describe an example where the triangulation T_h of $\overline{\Omega}$ consists of convex quadrilaterals. Thus for $T \in T_h$ there exists a bijection $F \in (P_{1,1})^2$ such that $T = F(\hat{T})$, with \hat{T} the reference square having vertices (0,0), (1,0), (1,1) and (0,1). For each point $\hat{x} \in \hat{T}$, the linear tangent mapping $DF(\hat{x})$ and its Jacobien $J(\hat{x}) = \det(DF(\hat{x}))$ are affine functions of the coordinates (\hat{x}_1,\hat{x}_2) .

For a scalar function $v \in H^2(T)$ we may define a function $\mathcal{L}_T^{(1)}v$ in the space

$$(7.17) \{w \in H^{1}(T) ; w \cdot F \in P_{1,1}(\hat{T})\}$$

by interpolating the values of v at the four vertices of the quadrilateral T.

Similarly for a vector function $\underline{q}\in H(\mathrm{div};T)$ we may interpolate, on each of the four edges, the average value of the normal trace to define a function $E_{\mathrm{T}}^{(1)}\underline{q}$ in the space

(7.18)
$$\{ \underline{r} \in H(div; T) ; |J|DF^{-1}(\underline{r} \circ F) \in P_{2,1}(\hat{T}) \times P_{1,2}(\hat{T}) \}.$$

Even though the spaces defined in (7.17) and (7.18) are spaces of polynomial functions only if T is a parallelogram, we note that in any case the trace of $\mathcal{L}_{T}^{(1)}$ on each edge T' of T is an affine function and the normal trace of $E_{T}^{(1)}$ q on T' is constant. Furthermore, $|J| \operatorname{div}(E_{T}^{(1)}q)$ is constant on T and we have

(7.19)
$$|J| \operatorname{div} (E_{T}^{(1)} \underline{q}) = \frac{1}{\text{mes } \hat{T}} \int_{T} \operatorname{div} \underline{q} \, dx.$$

Proceeding as in each of the preceding examples, by patching together the maps $\mathcal{L}_{T}^{(1)}$, $T \in \mathcal{T}_{h}$, we define an operator $\mathcal{L}_{h}^{(1)}$ on $H^{1}(\Omega) \cap H^{2}(\mathcal{T}_{h})$ having as image

CHAPTER III

ABSTRACT THEORY OF APPROXIMATION OF SOLUTIONS OF PROBLEMS

WITH LAGRANGIAN MULTIPLIERS

8 - GENERAL REMARKS

In this chapter we develop an abstract theory for the approximation of solutions of problems with Lagrangian multipliers. The main purpose for the construction of such a theory is to unify the analysis of the mixed and hybrid methods presented in Chapter I, and hence to provide a framework for generalizations and extensions of the analysis to that of similar problems.

Even though these methods were conceived of in the context of approximation of problems of minimization of a quadratic functional with affine constraints, the properties of constrained optimization are not at all fundamental to the analysis of the problems studied here. On the contrary, more than one familiar with conforming approximations of minimization problems seems to have been misguided by this approach. In the exposition of the abstract theory given in this chapter, we will make no restrictions which allow the class of problems studied here to be interpreted as being minimization problems on an affine manifold. The expression "with Lagrangian multiplier" will be used in the same formal sense as is the expression "variational problem".

Let us now give the general framework for this study. Let W and M be Hilbert spaces with inner products $((\cdot,\cdot))_W$ and $((\cdot,\cdot))_M$ respectively. The norms associated with these inner products are denoted $||\cdot||_W$ and $||\cdot||_M$. Further, let $a(\cdot,\cdot)$ and $b(\cdot,\cdot)$ be bilinear forms defined on WxW and on WxM respectively. Given the linear forms $f(\cdot)$ defined on W and $g(\cdot)$ defined on M, we consider the following problem: find a pair (ϕ,λ) satisfying

$$\begin{cases} (\phi,\lambda) \in W \times M , \\ a(\phi,\psi) + b(\psi,\lambda) = f(\psi) & \text{for all } \psi \in W, \\ b(\phi,\mu) = g(\mu) & \text{for all } \mu \in M. \end{cases}$$

We introduce the product space X=W×M which is a Hilbert space with norm

Clearly one may also formulate (8,1) in the following manner: find $u=(\phi,\lambda)$ satisfying

(8.2)
$$\begin{cases} u \in X, \\ k(u,v) = \ell(v) \end{cases}$$
 for all $v \in X$,

where k(•,•) is a bilinear form on X×X given by

(8.3)
$$k(u,v) = a(\phi,\psi) + b(\psi,\lambda) + b(\phi,\mu) \qquad \text{for all } u = (\phi,\lambda) \in X \text{ and}$$
 for all $v = (\psi,\mu) \in X$,

and where $\ell(\cdot)$ is a linear form defined on X by

(8.4)
$$\ell(v) = f(\psi) + g(\mu) \qquad \text{for all } v = (\psi, \mu) \in X.$$

All questions of continuity aside, it is evident that the analysis of the variational problem (8.2) can not be carried out with the aid of the Lax-Milgram theorem since the bilinear form $k(\cdot,\cdot)$ is not X-elliptic: for each $v=(0,\mu)\in X$, we have k(v,v)=0. One may show, nevertheless, the following result:

THEOREM 8.1. Suppose that the bilinear form $k(\cdot, \cdot)$ is continuous on $X \times X$ and that it satisfies

i)
$$\inf \sup_{\{u \in X \ ; \ ||\ u\ ||_{X}=1\}} k(u,v)>0,$$

and

ii)
$$\sup_{\{u \in X ; \|u\|_{X}=1\}} k(u,v) > 0 \qquad \text{for all } v \in X \text{ with } v \neq 0.$$

Then for each continuous linear form $\ell(\cdot)$ on X, there exists a unique solution of problem (8.2).

This theorem is a special case of the theorem due to BABUSKA [1971] given in the following section. \Box

<u>Remark</u> 8.1. Hypothesis ii) is superfluous when the bilinear form $k(\cdot, \cdot)$ is symmetric because, in this case, i) implies ii).

9 - ANALYSIS AND APPROXIMATION OF THE SOLUTION OF A VARIATIONAL PROBLEM (THEORY OF BABUSKA)

We consider the two Hilbert spaces X and Y equipped with the norms $\|\cdot\|_X$ and $\|\cdot\|_Y$, and let $k(\cdot,\cdot)$ be a bilinear form on X×Y. We have

THEOREM 9.1. Suppose that k(•,•) is a continuous, bilinear form on X×Y which satisfies

i)
$$\inf \qquad \sup \qquad k(u,v) > 0 ,$$

$$\{u \in X ; ||u||_{X} = 1\} \quad \{v \in Y; ||v||_{Y} = 1\}$$

and

ii)
$$\sup_{\{u \in X \ ; \ ||u||_{X}=1\}} k(u,v) > 0 \qquad \text{for all } v \in Y \text{ with } v \neq 0.$$

Then for each continuous linear form $\ell(\bullet)$ on Y, the problem: find u satisfying

(9.1)
$$\begin{cases} u \in X, \\ k(u,v) = \ell(v) \end{cases}$$
 for all $v \in Y$,

has a unique solution.

 \underline{Proof} . Define the continuous linear operator K from X to Y by

$$((Ku,v))_{Y} = k(u,v)$$
 for all $u \in X$ and for all $v \in Y$.

Hypothesis i) implies that the operator K is injective and that its image is closed in Y; hypothesis ii) expresses the fact that the image of K is dense in Y.

Remark 9.1. Conditions i) and ii) together give a necessary and sufficient condition for problem (9.1) to have a unique solution given a continuous, linear form $\ell(\cdot)$. In the case X=Y, the ellipticity condition,

inf
$$k(v,v) > 0$$
, $\{v \in X ; ||v||_{X}=1\}$

obviously implies conditions i) and ii); thus it is a sufficient condition for problem (9.1) to have a unique solution. In this sense Theorem 9.1 generalizes the Lax-Milgram theorem.

In the next section we shall come back to and treat in detail the "particular case" where $X=W\times M$ and $k(\cdot,\cdot)$ is of the form (9.3), i.e. the case that motivates our study. First, however, we give a general result concerning the approximation of the solution of problem (9.1). Given finite dimensional subspaces X_h and Y_h of X_h and Y_h respectively, we seek the "approximate solution" $u_h \in X_h$, an approximation of the solution u_h of problem (9.1), defined by

(9.2)
$$\begin{cases} u_h \in X_h, \\ k(u_h, v_h) = \ell(v_h) \end{cases}$$
 for all $v_h \in Y_h$.

Applying Theorem 9.1 in the finite dimensional case, we obtain:

THEOREM 9.2. Suppose that $k(\cdot, \cdot)$ is a bilinear form on $X_h \times Y_h$ which satisfies

i)
$$\inf_{\{u_h \in X_h ; ||u_h||_{X}=1\}} \sup_{\{v_h \in Y_h ; ||v_h||_{Y}=1\}} k(u_h, v_h) > 0,$$

and

$$\begin{aligned} &\sup & k(u_h, v_h) > 0 & & \textit{for all } v_h \in Y \textit{ with } v_h \neq 0. \\ &\{u_h \in X_h \; ; \; ||u_h||_{X} = 1 \} \end{aligned}$$

Then for each linear form $\ell(\cdot)$ on Y_h , problem (9.2) has a unique solution.

Remark 9.2. The space X_h being finite dimensional, condition i) is equivalent to

$$\sup_{\{v_h\in Y_h\ ;\ ||\ v_h||\ Y^{=1}\}} k(u_h,v_h) > 0 \qquad \text{for all } u_h\in X_h \text{ with } u_h\neq 0.$$

Futhermore, when the hypotheses of Theorem 9.2. are realized, one has necessarily that

(9.3)
$$\dim Y_h = \dim X_h.$$

Conversely, when we know, a priori, that (9.3) is satisfied (for example when $X_h = Y_h$), it suffices to check that either one or the other of conditions i) and ii) of Theorem 9.2 holds.

Remark 9.3. It is important to note that it is not sufficient that hypotheses i) and ii) of Theorem 9.1 be satisfied in order that the corresponding conditions of Theorem 9.2 be also. Herein lies the charm of the analysis!

Remark 9.4. Even when X=Y and the bilinear form $k(\cdot, \cdot)$ is X-elliptic, it can be useful to consider examples for which (9.2) holds with $X_h \neq Y_h$. Such is, for instance, the case in a Petrov-Galerkin type approach to a diffusion-convection problem, cf, for example MIZUKAMI & HUGHES [1985].

We conclude this section with an estimation of the error committed in approximating the solution u of problem (9.1) by the solution u_h of problem (9.2):

THEOREM 9.3 Suppose that $K(\bullet, \bullet)$ is a continuous, bilinear form on $X \times Y$ and let K be a number such that

Suppose that the choice of subspaces X_h and Y_h is such that

$$(9.5) dim Y_h = dim X_h$$

and that there exists a number $\kappa>0$ such that

(9.6)
$$\sup_{\{v_h \in Y_h ; \|v_h\|_{Y}=1\}} k(u_h, v_h) \ge \kappa \|u_h\|_{X} \quad \text{for all } u_h \in X_h.$$

Then for each linear form $\ell(\bullet)$ on Y, if u satisfies (9.1) and u_h is the solution to (9.2), we have the error bound

(9.7)
$$|| u-u_h ||_X \le (1+\frac{K}{\kappa}) \inf_{\substack{w_h \in X_h}} || u-w_h ||_X.$$

This result is due to BABUSKA [1971]; cf. also BABUSKA & AZIZ [1972, Theorem 6.2.1, p. 186]); a simple demonstration in the case $X_h = Y_h$ can be found in RABIER & THOMAS [1985, Exercise 3.13], for example. For a given problem of the form (9.1), the constants K and κ which appear in (9.4) and (9.6), respectively, obviously depend on the choice of norm on X and on Y and on the choice of the subspaces X_h and Y_h - nothing prevents one from making the choice of the norms as a function of the choice of the finite dimensional subspaces. We will see some such examples later on.

For the moment we shall settle for a simple criterion for convergence when the index h, indicating the choice of finite dimensional subspace X_h as well as a real parameter, tends toward 0. The set of parameters will be denoted H.

THEOREM 9.4. Suppose that the norms $\|\cdot\|_X$ and $\|\cdot\|_Y$ on the spaces X and Y, respectively, are chosen independently of the parameter $h \in H$ and that condition (9.4) is satisfied. Suppose that for each $h \in H$, we have been able to choose the finite dimensional subpaces X_h and Y_h to be of the same dimension and such that condition (9.6) is realized with a constant $\kappa > 0$ independent of h. Suppose further that there exist a subspace X of X dense in X and a mapping T_h of X into X_h such that

$$\lim_{h\to 0} || w-r_h(w) ||_X = 0 \qquad \text{for all } w \in X.$$

Then, the variational approximation method converges in the sense that we have

$$\lim_{h\to 0} ||u-u_h||_{X} = 0,$$

where u is the solution of (9.1) and u_h the solution of (9.2).

The demonstration of this theorem is identical to that showing the convergence of the approximate solutions of an elliptic variational problem, cf. for example RAVIART & THOMAS [1983, Theorem 3.1-3].

10 - ANALYSIS AND APPROXIMATION OF THE SOLUTION OF A VARIATIONAL PROBLEM (THEORY OF BREZZI)

Equipped with the general results of the preceding section, we proceed with the analysis of a problem of the form (8.1) which we shall call a problem with Lagrangian multipliers: find a pair (ϕ, λ) satisfing

(10.1)
$$\begin{cases} (\phi,\lambda) \in W \times M, \\ a(\phi,\psi) + b(\psi,\lambda) = f(\psi) & \text{for all } \psi \in W, \\ b(\phi,\mu) = g(\mu) & \text{for all } \mu \in M. \end{cases}$$

In particular if a pair (ϕ, λ) satisfies (10.1) with $g(\cdot)=0$, then the first argument ϕ is a solution of the variational problem : find ϕ such that

(10.2)
$$\begin{cases} \phi \in V, \\ a(\phi, \psi) = f(\psi) \end{cases} \text{ for all } \psi \in V,$$

where V is the subspace of W defined by

(10.3)
$$V = \{ \psi \in W ; \forall \mu \in M, b(\psi,\mu) = 0 \}.$$

The following theorem is due to BREZZI [1974a, 1974b]:

THEOREM 10.1. Suppose that the bilinear form a(•,•) is continuous on W×W and V-elliptic, i.e.

i)
$$\inf \qquad a(\psi,\psi) > 0.$$

$$\{\psi \in V ; ||\psi||_{W}=1 \}$$

Suppose that the bilinear form b(•,•) is continuous on WxM and that it satisfies

ii)
$$\inf \sup_{\{\mu \in M ; \|\mu\|_{M}=1\}} \sup_{\{\psi \in W ; \|\psi\|_{W}=1\}} b(\psi,\mu) > 0.$$

Then for each pair of continuous linear forms $f(\cdot)$ on W and $g(\cdot)$ on M, problem (10.1) has a unique solution.

<u>Proof</u>. The original demonstration of this theorem is given in BREZZI [1974b]; a complete demonstration may also be found in GIRAULT & RAVIART [1979, Theorem 4.1] for example. We shall indicate here a variant of the demonstration of this fundamental result. Let V^{\perp} denote the subspace of W orthogonal to V; denote by ψ^0 and ψ^{\perp} the orthogonal projections of $\psi \in W$ in V and V^{\perp} , respectively. With this notation, problem (10.1) may be written: find ϕ^{\perp} , ϕ^0 , and λ such that

$$\begin{pmatrix}
\phi^{\perp} \in V^{\perp}, & \phi^{0} \in V, & \lambda \in M, \\
a(\phi^{\perp}, \psi^{\perp}) + a(\phi^{0}, \psi^{\perp}) + b(\psi^{\perp}, \lambda) & = f(\psi^{\perp}) & \text{for all } \psi^{\perp} \in V^{\perp}, \\
a(\phi^{\perp}, \psi^{0}) + a(\phi^{0}, \psi^{0}) & = f(\psi^{0}) & \text{for all } \psi^{0} \in V, \\
b(\phi^{\perp}, \mu) & = g(\mu) & \text{for all } \mu \in M.
\end{pmatrix}$$

Note that we have the equality

$$\sup \qquad b(\psi,\mu) = \sup b(\psi,\mu) \qquad \text{for all } \mu \in M,$$

$$\{\psi \in V^{\perp} ; ||\psi||_{W}=1\} \qquad \{\psi \in W ; ||\psi||_{W}=1\}$$

and that we have the characterization of V as a subset of W:

$$\psi \in V$$
 if and only if \sup $b(\psi,\mu) = 0$.
$$\{\mu \in M \; ; \; ||\mu||_{M} = 1\}$$

Using hypothesis ii) of Theorem 10.1 one may then establish with the aid of Theorem 9.1 that the linear operator B from M into V^{\perp} defined by

$$((\psi, B^*\mu))_W = b(\psi, \mu)$$
 for all $\psi \in W, \mu \in M$,

is an isomorphism of M onto V^{\perp} . Thus the linear operator B from V^{\perp} into M defined by

$$((B \psi, \mu))_{\mathbf{M}} = ((\psi, B^* \mu))_{\mathbf{W}}$$
 for all $\psi \in \mathbf{W}, \mu \in \mathbf{M}$

is an isomorphism of V^{\perp} onto M.

Finally, using hypothesis i) of Theorem 10.1 one may show with the aid of the Lax Milgram theorem that the linear operator A from V into V defined by

$$((A\phi, \psi))_{W} = a(\phi, \psi)$$
 for all $\phi \in V$, $\psi \in W$,

is an automorphism. Thus the linear system (10.4), to which is associated a matrix of operators of the form

$$\begin{pmatrix} \times & \times & B^* \\ \times & A & 0 \\ B & 0 & 0 \end{pmatrix}$$

has a unique solution.

Remark 10.1. Condition ii) of Theorem 10.1 is often called the inf-sup condition; some authors refer to it as the <u>Brezzi condition</u>, others as the <u>Brezzi-Babuska condition</u>. The current trend seems to be to use the expression <u>LBB condition</u>, where LBB stands for LADYZHENSKAYA [1949] - BABUSKA [1971] - BREZZI [1974].

Remark 10.2. A necessary and sufficient condition for problem (10.1) to have a unique solution $(\phi, \lambda) \in W \times M$ given the linear forms $f(\cdot)$ on W and $g(\cdot)$ on M is that the inf-sup condition be satisfied and that the operator $A \in \mathcal{L}(V;V)$ associated with the bilinear form $a(\cdot, \cdot)$ be an isomorphism. This result, through presented differently, can be found in GIRAULT & RAVIART [1985, Theorem 4.1].

With the goal of approximating the solution of (10.1), we take now finite dimensional subspaces of W and of M,

$$(10.5) W_h \subset W, M_h \subset M,$$

and define the discrete problem : find (ϕ_h, λ_h) satisfying

$$\begin{cases} (\phi_h, \lambda_h) \in W_h \times M_h, \\ a(\phi_h, \psi_h) + b(\psi_h, \lambda_h) = f(\psi_h) & \text{for all } \psi_h \in W_h, \\ b(\phi_h, \mu_h) = g(\mu_h) & \text{for all } \mu_h \in M_h. \end{cases}$$

As in the continuous case, we introduce the subspace V_h of W_h given by

(10.7)
$$V_h = \{ \psi_h \in W_h ; \forall \mu_h \in M_h, b(\psi_h, \mu_h) = 0 \}.$$

One should be careful to note that in general V_h is not a subspace of the space V defined in (10.3). Transposing to the finite dimensional case the results of Theorem 10.1, we obtain

THEOREM 10.2. Suppose that

- i) the bilinear form $a(\cdot, \cdot)$ is V_h -elliptic and
 - ii) the bilinear form b(•,•) satisfies

$$\inf \qquad \sup \qquad b(\psi_h,\mu_h) > 0. \\ \{\mu_h \in M_h \; ; \; ||\mu_h||_{M}^{=1}\} \qquad \{\psi_h \; \in W_h \; ; \; ||\psi_h||_{W}^{=1}\}$$

Then for each pair of linear forms $f(\cdot)$ on W_h and $g(\cdot)$ on M_h , problem (10.6) has a unique solution.

Remark 10.3. Condition ii) of Theorem 10.2 is called the discrete inf-sup condition or again the discrete LBB condition. It says that the choice of the subspace W_h of W and that of the subspace M_h of M cannot be made independently one

from the other, that there is a compatibility relation between the two subspaces. If we introduce the bases $(\psi_h^{(j)})_{j=1,\dots,\dim W_h}$ and $(\mu_h^{(i)})_{i=1,\dots,\dim M_h}$ of the subspaces W_h and M_h respectively, the discrete inf-sup condition may be interpreted as expressing the fact that the matrix $[B_h]$ whose i,j entry is $[B_h]_{i,j} = b(\psi_h^{(i)}, \mu_h^{(i)})$ has rank equal to the dimension of M_h :

$$rg[B_h] = dim M_h$$

In this form, the condition has been called the <u>rank condition</u> (cf. BABUSKA, ODEN & LEE, [1978]. Clearly a necessary condition for this condition to hold is that

$$\dim M_h \leq \dim W_h.$$

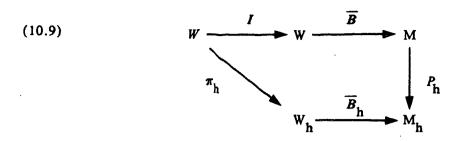
Remark 10.4. Obviously if the bilinear form $a(\cdot, \cdot)$ is V-elliptic and $V_h \subset V$ (as will be the case for dual mixed methods), then $a(\cdot, \cdot)$ is V_h -elliptic, or better yet if $a(\cdot, \cdot)$ is W-elliptic, so one no longer needs to have $V_h \subset V$ (as will be the case for dual hybrid and primal mixed methods), then $a(\cdot, \cdot)$ is V_h -elliptic. In other cases it is necessary to verify directly the V_h -ellipticity of $a(\cdot, \cdot)$, something which can be rather technical (this will be the situation for primal hybrid methods).

To verify the discrete inf-sup condition one can use the following result (cf. FORTIN, [1977]), the proof of which is immediate:

<u>THEOREM 10.3.</u> Suppose there is a Hilbert space $W \subset W$, an operator R_h from M_h into W satisfying

(10.8)
$$\frac{b(R_{h}\mu_{h},\mu_{h})}{||R_{h}\mu_{h}||_{W}} \geq C \sup_{\{\psi \in W; \||\psi||_{W}=1\}} b(\psi,\mu_{h}) \quad \text{for all } \mu_{h} \in M_{h},$$

for some $C \ge O$, and a continuous linear operator π_h from W into W_h such that the following diagram commutes:



where I is the inclusion of W into W, \overline{B} and \overline{B}_h are the mappings from W to M and from W_h to M_h , respectively, induced by $b(\cdot, \cdot)^{(1)}$ and P_h is the orthogonal projection of M onto M_h . Then the inf-sup condition of Theorem 10.1 implies the discrete inf-sup condition of Theorem 10.2.

We give an estimate of the error made in approximating the solution of problem (10.1) by that of problem (10.6). This estimate was obtained by BREZZI [1974a, 1974b].

<u>THEOREM 10.4.</u> Suppose that $a(\cdot, \cdot)$ and $b(\cdot, \cdot)$ are continuous bilinear forms on W×W and on W×M respectively. Let A and B be numbers such that

(10.10)
$$a(\phi,\psi) \leq A \|\phi\|_{W} \|\psi\|_{W} \qquad \text{for all } \phi \in W \text{ and } \psi \in W,$$

(10.11)
$$b(\psi,\mu) \leq B \|\psi\|_{W} \|\mu\|_{M} \qquad \text{for all } \psi \in W \text{ and } \mu \in M.$$

Suppose further that the subspaces W_h and M_h have been chosen such that there exist positive numbers α and β for which we have

(10.12)
$$a(v_h, v_h) \ge \alpha ||v_h||_W^{12} \qquad \text{for all } v_h \in V_h,$$

(10.13)
$$\inf \sup_{\{\mu_h \in M_h \ ; \ ||\mu_h||_{M}=1\}} b(\psi_h, \mu_h) \ge \beta.$$

Then there exists a constant C, depending only on A, B, α , and β , such that for each pair of linear forms $f(\cdot)$ on W and $g(\cdot)$ on M, if (ϕ,λ) satisfies (10.1) and (ϕ_h,λ_h) satisfies (10.5), we have

⁽¹⁾ The symbols B and B_h have been reserved for the restrictions of \overline{B} to V_h^{\perp} and of \overline{B}_h to V_h^{\perp} , respectively.

(10.14)
$$||\phi - \phi_h||_{W} = ||\lambda - \lambda_h||_{M} \le C \left\{ \inf_{\psi_h \in W_h} ||\phi - \psi_h||_{W} + \inf_{\mu_h \in M_h} ||\lambda - \mu_h||_{M} \right\}.$$

<u>Proof</u>. Let V_h^{\perp} be the subspace of W_h orthogonal to V. We introduce, as in the proof of Theorem 10.1, the operator $B_h \in \mathcal{L}(V_h^{\perp}, M)$ and its adjoint $B_h^* \in \mathcal{L}(M, V_h^{\perp})$. The discrete inf-sup condition (10.13) signifies that the operators B_h and B_h^* are invertible or more precisely that we have

$$||B_h^{-1}||_{\mathcal{L}(M_h^{\perp}, V_h^{\perp})} = ||B_h^{*-1}||_{\mathcal{L}(V_h^{\perp}, M_h^{\perp})} \le \frac{1}{\beta}.$$

Similary the condition of V_h -ellipticity (10.12) implies that the operator $A_h \in$ $\mathcal{L}(V_h$, V_h) associated with the bilinear form $a(\cdot,\cdot)$ is invertible and that we have

$$||A_h^{-1}||_{\mathscr{L}(V_h, V_h)} \leq \frac{1}{\alpha}.$$

It is now easy to check that the bilinear form defined with the aid of (9.3) on XxX, where X=WxM (and Y=X), satisfies (9.4) with a constant K depending only on A and B and for the choice $X_h = W_h \times M_h$ (and $Y_h = X_h$) satisfies (9.6) with a constant $\kappa>0$ which depends only on A, α and β . An application of Theorem 9.3 then leads us to the estimation (10.14).

Remark 10.5 . This error estimate is optimal in the product space X=W×M. This does not mean that the estimation of a single term, $||\phi-\phi_h||_W$ for example, is necessarily optimal. Say for instance that $g(\cdot) = 0$ and V_h is a subspace of V_h Then we obviously have

$$||\phi-\phi_h||_W \leq C \inf_{\psi_h \in V_h} ||\phi-\psi_h||_W.$$

This sort of refinement of the theory has been developed by FALK & OSBORN [1980]. We will point out such refinements in the examples treated hereafter.

We can deduce from Theorem 10.4 a convergence theorem similar to Theorem 9.4 for a family of pairs of subspaces (W_h, M_h) indexed by $h \in H$. On the other hand we point out that the approximation result of Theorem 10.4 is still of interest if the chosen norms or even the spaces W and M themselves change with the parameter h as will indeed be the case for hybrid methods. The error bound (10.14) holds with constant C independent of h as long as the constants A, B, α and β may be taken independently of h.

Remark 10.7 • Theorem 10.3 gives sufficient conditions for concluding that the inf-sup condition for $b(\cdot, \cdot)$ on W×M is inherited by $b(\cdot, \cdot)$ on W_hxM_h. However, given a family of pairs of subspaces (W_h, M_h) indexed by $h \in H$, to obtain convergence results (cf. Remark 10.6) we need to know that the discrete inf-sup condition holds uniformly in the sense that there exists β independent of h such that for each $h \in H$

$$\inf_{\{\mu_h \in M_h; \ ||\mu_h||_{M}=1\}} \sup_{\{\psi_h \in W_h; ||\psi_h||_{W}=1\}} b(\psi_h, \mu_h) \geq \beta > 0.$$

This will be the case if the norms of R_h and π_h of Theorem 10.3 are independent of h.

Remark 10.8. With the methods used here to analyze numerically a problem with Lagrangian multipliers of the form (10.1) whose solution is approximated by the solution of a finite dimensional problem of the form (10.6), it is easy to find sufficient conditions for analyzing a problem of the following form: find ϕ and λ satisfying

(10.16)
$$\begin{cases} \phi \in W , \lambda \in M , \\ a(\phi, \psi) + b(\psi, \lambda) = f(\psi) & \text{for all } \psi \in W , \\ b(\phi, \mu) + d(\lambda, \mu) = g(\mu) & \text{for all } \mu \in M. \end{cases}$$

This generalization introduces a supplementary bilinear form $d(\cdot, \cdot)$ on M×M which we suppose is continuous, and we let D be a constant such that

(10.17)
$$d(\lambda,\mu) \le D \|\lambda\|_{M} \|\mu\|_{M} \qquad \text{for all } \lambda \in M \text{ and } \mu \in M.$$

Given finite dimensional subspaces $W_h \subset W$ and $M_h \subset M$, the corresponding discrete problem associated with (10.16) is written as follows: find ϕ_h and λ_h satisfying

$$\begin{cases} \phi_h \in W_h, & \lambda_h \in M_h, \\ a(\phi_h, \psi_h) + b(\psi_h, \lambda_h) = f(\psi_h) & \text{for all } \psi_h \in W_h, \\ b(\phi_h, \mu_h) + d(\lambda_h, \mu_h) = g(\mu_h) & \text{for all } \mu_h \in M_h. \end{cases}$$

If in addition to the hypotheses of Theorem 10.4 we suppose that

$$\begin{array}{ll} \mathbf{a}(\psi_h,\psi_h) \!\!\ge\!\! 0 & \text{for all } \psi_h \!\!\in\!\! \mathbf{W}_h \ , \\ \mathbf{d}(\mu_h,\!\mu_h) \!\!\le\!\! 0 & \text{for all } \mu_h \!\!\in\!\! \mathbf{M}_h \ , \end{array}$$

and that one of the two forms $a(\cdot, \cdot)$ and $d(\cdot, \cdot)$ is symmetric, then problem (10.18) has a unique solution (ϕ_h, λ_h) . Moreover, if (ϕ, λ) denotes the solution of (10.16), an error bound analogous to that given by (10.14) may be obtained with a constant C depending only on A, B, D, α and β . Another generalization of Theorem 10.4 is obtained by assuming in addition to the hypotheses of Theorem 10.4 that there is a positive number δ , $0 < \delta \le 1$, such that

The conclusion is the same as in the above generalization except that here the constant C depends also on δ .

The demonstration of the first generalization relies on the fact that if H is a Hilbert space with inner product $((\ ,\))_H$ and $U\in \mathcal{L}(H,H)$ is positive in the sense that $((Ux,x))_H\geq 0$, for all $x\in H$, then I+U is invertible in $\mathcal{L}(H,H)$ and $||(I+U)^{-1}||_{\mathcal{L}(H,H)}\leq 1$. To show that the second generalization is valid one uses that if H is a Hilbert space and $U\in \mathcal{L}(H,H)$ has norm in $\mathcal{L}(H,H)$ less than $1-\theta$ for some θ in (0,1], then I+U is invertible in $\mathcal{L}(H,H)$ and $||(I+U^{-1})||_{\mathcal{L}(H,H)}\leq \frac{1}{\theta}$.

<u>Remark 10.9.</u> An extension of the above results of BREZZI [1974b] to the case of the following generalized saddle point problem:

$$\begin{cases} (\phi,\lambda) \in W \times M \ , \\ \\ a(\phi,\psi) + b_1(\psi,\lambda) = f(\psi) \end{cases} \qquad \text{for all } \psi \in W \ , \\ \\ b_2(\phi,\mu) = g(\mu) \qquad \text{for all } \mu \in M \ , \end{cases}$$

for which the bilinear forms $b_1(\cdot,\cdot)$ and $b_2(\cdot,\cdot)$ on W×M are distinct, can be deduced from the more extensive generalization given in NICOLAIDES [1982].

11 - EFFECTS OF NUMERICAL INTEGRATION

Under the hypotheses of Theorem 10.2, problem (10.6), the problem with Lagrangian multiplier, is reduced, once the bases of W_h and of M_h have been chosen, to the resolution of an invertible linear system of order N, with $N = \dim W_h + \dim M_h$. Several algorithms adapted to the resolution of such linear systems will be presented in Section 17. For the moment we remark only that even the formation of the linear system can turn out to be practically impossible or simply too expensive in calculation time. Thus we are led to replace problem (10.6) by the following problem: find (ϕ_h^*, λ_h^*) satisfying

$$\begin{pmatrix} (\phi_h^*, \lambda_h^*) \in W_h \times M_h, \\ a_h(\phi_h^*, \psi_h) + b_h(\psi_h, \lambda_h^*) = f_h(\psi_h) & \text{for all } \psi_h \in W_h, \\ b_h(\phi_h^*, \mu_h) = g_h(\mu_h) & \text{for all } \mu_h \in M_h; \end{pmatrix}$$

where $a_h(\cdot,\cdot)$ and $b_h(\cdot,\cdot)$ are bilinear forms on $W_h\times W_h$ and $W_h\times M_h$, respectively, and $f_h(\cdot)$ and $g_h(\cdot)$ are linear forms on W_h and M_h , respectively. When $g_h(\cdot)=0$, if (ϕ_h^*,λ_h^*) is the solution of (11.1), then the first argument ϕ_h^* is a solution of the variational problem: find a solution ϕ_h^* such that

(11.2)
$$\begin{cases} \phi_h^* \in V_h^*, \\ a_h(\phi_h^*, \psi_h) = f_h(\psi_h) \end{cases}$$
 for all $\psi_h \in V_h^*$;

where the subspace \boldsymbol{V}_{h}^{*} is defined by

(11.3)
$$V_{h}^{*} = \left\{ \psi_{h} \in W_{h}, \forall \mu_{h} \in M_{h}, b_{h}(\psi_{h}, \mu_{h}) = 0 \right\}.$$

<u>Remark 11.1</u> For all of the examples that we shall consider, the physical parameters of the modelled problem actually play a role only in the expression of the bilinear form $a(\cdot,\cdot)$. The form $b(\cdot,\cdot)$ is independent of these parameters; the same is true, a fortiori, for the subspace V defined by (10.3).

Futhermore, the spaces W_h and M_h are constructed in such a way that the exact calculation of the term $b(\psi_h, \mu_h)$ poses no practical difficulty since it always reduces to the calculation of an integral on an affine manifold of a polynomial function. In the following we shall point out (cf. Remark 11.2) where it can be of some interest to replace the bilinear form $b(\bullet, \bullet)$ by another $b_h(\bullet, \bullet)$.

Defining X_h to be the product space $W_h \times M_h$ and putting

(11.4)
$$\begin{cases} k_h(u_h, v_h) = a(\phi_h, \psi_h) + b(\psi_h, \lambda_h) + b(\phi_h, \mu_h) \\ and \\ \ell_h(v_h) = f_h(\psi_h) + g_h(\mu_h), \end{cases}$$

for $u_h = (\phi_h, \lambda_h)$ and $v_h = (\psi_h, \mu_h)$, we can write, as in Section 9, problem (11.1) in the form : find u_h^* satisfying

$$\begin{cases} u_h^* \in X_h, \\ k_h(u_h^*, v_h) = \ell_h(v_h). \end{cases}$$
 for all $u_h \in X_h$.

Then in a slightly more general setting, this leads us to study the approximation of the solution u of problem (9.2) by the solution u_h^* of the following problem: find u_h^* satisfying

(11.5)
$$\begin{cases} u_h^* \in X_h, \\ k_h(u_h^*, v_h) = \ell_h(v_h) \end{cases} \text{ for all } v_h \in Y_h,$$

where $k_h(\cdot, \cdot)$ is a bilinear form on $X_h \times Y_h$ and $\ell_h(\cdot)$ is a linear form on Y_h , with $X_h \subset X$ and $Y_h \subset Y$. The analysis is based on the following result:

<u>THEOREM 11.1.</u> Suppose that $k(\cdot, \cdot)$ is a bilinear form, continuous on $X \times Y$ and let K be a number such that (9.4) is satisfied. We suppose that the choice of subspaces X_h and Y_h is such that

$$\dim Y_h = \dim X_h$$

and that there exists a number $\kappa_h > 0$ such that

(11.6)
$$\sup_{\{v_h \in Y_h; \|v_h\|_{Y}=1\}} k_h(u_h, v_h) \ge \kappa_h \|u_h\|_{X} \qquad \text{for all } u_h \in X_h.$$

Then there exists a constant C, depending only on K and κ_h , such that if u is a solution of (9.2) and u_h is the solution of (11.5), we have

$$||u-u_{h}^{*}||_{X} \leq C \left\{ \inf_{w_{h} \in X_{h}} \left(||u-w_{h}||_{X} + \sup_{v_{h} \in Y_{h}} \frac{a(w_{h}, v_{h}) - a_{h}(w_{h}, v_{h})}{||v_{h}||_{Y}} \right) + \sup_{v_{h} \in Y_{h}} \frac{\ell(v_{h}) - \ell_{h}(v_{h})}{||v_{h}||_{Y}} \right\}.$$

This is, of course, in the case $a_h(\cdot,\cdot) = a(\cdot,\cdot)$ and $\ell_h(\cdot) = \ell(\cdot)$, the result (9.7) given in Theorem 9.3. It generalizes the results of STRANG [1972] developed for the elliptic variational problem. The demonstration may be obtained by simply adapting the demonstration of Strang's result given by CIARLET [1977, pp.186-187] for example.

We return now to the discrete problem (11.1), Theorem 10.4 may be generalized to the following result:

THEOREM 11.2. Suppose that $a(\cdot, \cdot)$ and $b(\cdot, \cdot)$ are continuous bilinear forms on W×W and on W×M respectively, and let A and B be numbers such that (10.10) and (10.11) are satisfied. Let A_b be a number such that

$$(11.8) a_h(\phi_h, \psi_h) \le A_h \|\phi_h\|_W \|\psi_h\|_W for all \phi_h \in W_h and for all \psi_h \in W_h.$$

Suppose further that there exists numbers $a_h>0$ and $\beta_h>0$ such that

(11.9)
$$a_h(v_h, v_h) \ge \alpha_h ||v_h||_W^2 \qquad \qquad \text{for all } v_h \in V_h^*,$$

and

(11.10) inf sup
$$b(\psi_h, \mu_h) \ge \beta_h$$
.
 $\{\mu_h \in M_h \; ; \; ||\mu_h||_{M} = 1\}$ $\{\psi_h \in W_h \; ; \; ||\psi_h||_{W} = 1\}$

Then there exists a constant C dependent only on A, B, A_h , α_h and β_h such that if (ϕ, λ) satisfies (10.1) and (ϕ_h^*, λ_h^*) is the solution of (11.1) then we have

$$\begin{array}{lll} \text{(II.II)} & & ||\phi-\phi_{h}^{*}||_{W} + & ||\lambda-\lambda_{h}^{*}||_{M} \leq \\ \\ \text{C} & \left\{ \inf_{\psi_{h} \in W_{h}} \left(||\phi-\psi_{h}||_{W} + \sup_{\eta_{h} \in W_{h}} \frac{a(\psi_{h},\eta_{h}) - a_{h}(\psi_{h},\eta_{h})}{||\eta_{h}||_{W}} \right) \\ \\ + & \inf_{\mu_{h} \in M_{h}} \left(||\lambda-\mu_{h}||_{M} + \sup_{\eta_{h} \in W_{h}} \frac{b(\eta_{h},\mu_{h}) - b_{h}(\eta_{h},\mu_{h})}{||\eta_{h}||_{W}} \right) \\ \\ + & \sup_{\eta_{h} \in W_{h}} \frac{f(\eta_{h}) - f_{h}(\eta_{h})}{||\eta_{h}||_{W}} + \sup_{\nu_{h} \in M_{h}} \frac{g(\nu_{h}) - g_{h}(\nu_{h})}{||\nu_{h}||_{M_{h}}} \right). \end{array}$$

We may obtain this result by applying Theorem 11.1. With the given hypotheses and the definition of the bilinear form $k_h(\cdot,\cdot)$ given in (11.4), condition (9.4) is satisfied with K depending only on A and B and condition (11.6) is satisfied with $\kappa_h>0$ depending only on A_h , α_h and β_h . The detailed demonstration of this type of result was given by THOMAS [1977, Chapter I, Paragraph 3].

Remark 11.2. The replacement of the bilinear form $b(\cdot, \cdot)$ by the bilinear form $b_h(\cdot, \cdot)$ entails the replacement of the subspace V_h defined by (10.7) by the subspace V_h^* defined by (11.3). The ellipticity hypothesis (11.9) can certainly be satisfied without (10.12)'s being satisfied, even in the case where $a_h(\cdot, \cdot) = a(\cdot, \cdot)$ on $W_h \times M_h$. Also, the discrete inf-sup condition (11.10) may hold without (10.13)'s holding. On the contrary, the converses are valid; (10.12) implies (11.9) and (10.13) implies (11.10). However, we stress that a careless choice of numerical integration technique can turn a well posed discrete problem of the form (10.6) into an ill posed discrete problem of the form (11.1). One should be aware that by replacing $b(\cdot, \cdot)$ by $b_h(\cdot, \cdot)$ one is actually changing the numerical method, and this must be done with some care.

CHAPTER IV

MIXED FINITE ELEMENT METHODS

12 - EXAMPLES OF PRIMAL MIXED METHODS FOR A DIRICHLET PROBLEM

A primal mixed variational formulation of the model problem, find a solution (p,u) of the system (4.4) - (4.5) with the homogeneous Dirichlet boundary condition u=0 on Γ , is obtained by characterizing the pair (p,u) as the solution of

$$(12.1) \qquad \left\{ \begin{array}{l} (\underline{p},u) \in \left(L^2(\Omega)\right)^n \times H^1_0(\Omega) \ , \\ \\ a(\underline{p},\underline{q}) + b(\underline{q},u) = 0 \qquad \qquad \text{for all } \underline{q} \in \left(L^2(\Omega)\right)^n, \\ \\ b(\underline{p},v) \qquad = -\int_{\Omega} fv \ dx \qquad \qquad \text{for all } v \in H^1_0(\Omega), \end{array} \right.$$

where

(12.2)
$$a(p,q) = \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} p_{j} q_{i} dx,$$

and

(12.3)
$$b(q,v) = -\int_{\Omega} q \cdot \operatorname{grad} v \, dx.$$

We are here in the context of Chapter III with

(12.4)
$$W = (L^2(\Omega))^n$$
, $\| \underline{q} \|_{W} = \| \underline{q} \|_{Q,\Omega}$,

and

(12.5)
$$M = H_0^1(\Omega), \qquad || v ||_M = |v|_{1,\Omega}.$$

With the given hypotheses concerning the coefficients a_{ij} , the bilinear form $a(\cdot, \cdot)$ is continuous on W × W and W-elliptic with constant of ellipticity $\alpha = \alpha_d$, cf. (4.8). The bilinear form $b(\cdot, \cdot)$ is continuous on W × M and we clearly have

(12.6)
$$|v|_{1,\Omega} = 1$$
 $\Rightarrow \sup_{\{g \in W; ||g||_{0,\Omega} = 1\}} b(g,v) = 1.$

Thus Theorem 10.1 guarantees that, for each $f \in L^2(\Omega)$, problem (12.1) has a unique solution.

In order to construct a primal mixed approximation to this solution, we introduce a finite dimensional subspace W_h of $(L^2(\Omega))^n$ and a finite dimensional subspace M_h of $H^1_0(\Omega)$. From Theorems 10.2 and 10.4 we deduce :

THEOREM 12.1 . Suppose that the subspaces W_h and M_h are compatible in the sense that the discrete inf-sup condition is satisfied :

(12.7)
$$\beta_{h} = \inf_{\{v_{h} \in M_{h} ; |v_{h}|_{1,\Omega} = 1\}} \sup_{\{q_{h} \in W_{h} ; ||q_{h}||_{0,\Omega} = 1\}} b(q_{h}, v_{h}) > 0.$$

Then the problem of finding a pair (p_h, u_h) satisfying

$$\left(\begin{array}{c} (\underline{p}_h,u_h) \in W_h \times M_h \\ \\ a(\underline{p}_h,\underline{q}_h) + b(\underline{q}_h,u_h) = 0 \\ \\ b(\underline{p}_h,v_h) \end{array} \right) = 0 \qquad \qquad \text{for all } \underline{q}_h \in W_h , \\ b(\underline{p}_h,v_h) = -\int_{\Omega} fv_h \, dx \qquad \qquad \text{for all } v_h \in M_h , \\ \end{array}$$

has a unique solution. Furthermore, for each family of pairs of subspaces $(W_h^{}, M_h^{})$ satisfying (12.7) uniformly in the sense that $\beta = \inf \beta_h^{} > 0$, there exists a constant C, independent of h, such that

Given a triangulation T_h by n-simplexes T of a domain $\bar{\Omega}$, supposed to be polyhedral, and a positive integer k, we take finite dimensional subspaces $W_h = W_h^{(k)}$ and $M_h = M_h^{(k)}$ given as follows:

(12.10)
$$W_h = \{q \in (L^2(\Omega))^n ; V T \in T_h, q_T \in (P_{k-1}(T))^n\},$$

and

(12.11)
$$M_h = \{ v \in H_0^1(\Omega) ; V T \in T_h, v_{|T} \in P_k(T) \}.$$

(With the notation (5.10) introduced in Chapter II we have $M_h = H_0^1(\Omega) \cap L_h^{(k)}$). Such a choice of subspaces is proposed in REDDY & ODEN (1975); cf. also BABUSKA, ODEN, & LEE (1977, Section 4).

THEOREM 12.2. For each positive integer k ,the subspaces \mathbb{W}_h and \mathbb{M}_h defined by (12.10) and (12.11) satisfy the discrete inf-sup condition (12.7). Moreover, if (T_h ; $h \in \mathcal{H}$) is a regular family of triangulations, there is a constant C, independent of h, such that if the solution (p,u) of (12.1) belongs to $(\mathbb{H}^k(\Omega))^n \times \mathbb{H}^{k+1}(\Omega)$ we have the error estimate

In virtue of Theorem 12.1, the demonstration of Theorem 12.2 consists essentially of establishing the existence of a constant $\beta>0$ such that for each function $v_h\in {}^M_h$, $v_h\neq 0$, there exists a function $\underline{q}_h\in {}^M_h$, $\underline{q}_h\neq 0$, such that

$$b(\underline{q}_h, v_h) \geq \beta \| \underline{q}_h \|_{0,\Omega} | v_h |_{1,\Omega}.$$

With the choice of subspaces given in (12.10) and (12.11), we clearly have that for each $v_h \in {}^M_h$, grad v_h belongs to W_h . The compatibility condition is

thus satisfied with $\beta=1$. The proof is now terminated with an application of classical results giving error bounds for interpolation in $(L^2(\Omega))^n$ and in $H_0^1(\Omega)$. (As $L_h^{(k)}$ u belongs to M_h , we may use Theorem 5.2).

<u>Remark 12.1</u>. When the coefficients a_{ij} are sufficiently regular, for example when $a_{ij} \in C^k$, the regularity hypothesis $u \in H^{k+1}(\Omega)$ implies the regularity $p \in (H^k(\Omega))^n$.

<u>Remark 12.2</u>. Even with no regularity hypothesis for the solution (p,u) one can show that the primal mixed method applied in the context of Theorem 12.2 is convergent:

$$\lim_{h\to 0} p_h = p \text{ in } (L^2(\Omega))^n, \qquad \lim_{h\to 0} u_h = u \text{ in } H^1_0(\Omega) ;$$

cf. Remark 10.6.

<u>Remark 12.3</u>. With $W_h = W_h^{(k)}$ as given in (12.10) and with $M_h = M_h^{(k')}$ as given in (12.11) but with k' not necessarily equal to k, the compatibility hypothesis is satisfied once $k' \le k$. Error estimates optimal in $W \times M$ are obtained when k' = k.

<u>Remark 12.4</u>. Let \hat{T} be the reference n-simplex and for each $T \in T_h$, let F_T denote an affine bijection of \hat{T} onto T with Jacobien J_T . The spaces W_h and M_h given by (12.10) and (12.11) may also be defined as follows:

$$(12.10 \text{ bis}) \ W_{h} = \{\underline{q} \in (L^{2}(\Omega))^{n} ; \ V \ T \in T_{h}, \ J_{T} \ DF_{T}^{-1} \ (\underline{q}_{\mid T^{0}} \ F_{T}) \in (P_{k-1}(\hat{T}))^{n} \} ,$$

It is with these definitions that the theory is extended to the case of curved finite elements where T is the image of \hat{T} under a bijection not required to be affine. However, we must note that for $v_h \in M_h$, grad v_h need no longer belong to W_h . To show the compatibility relation in this case we introduce for $v_h \in M_h$, the function $\underline{q}_h \in W_h$ defined on each $T \in T_h$ by

$$(\underline{q}_{h|T}) \circ F_{T} = \frac{1}{J_{T}} (DF_{T}) (DF_{T})^{t} (\underline{grad} v_{h|T}) \circ F_{T}$$

and apply the relations (6.23).

These primal mixed methods may be easily adapted to the case in which $T_{\rm h}$ is a triangulation by rectangles having sides parallel to the axes for n=2, by rectangular solids having faces parallel to the coordinate planes for n=3. For a regular family of such triangulations one may show:

THEOREM 12.3 . (Stated here only for the case n=2, T a rectangle). For each positive integer k, with $W_h = W_h^{(k)}$ and $M_h = M_h^{(k)}$ defined as follows:

(12.13)
$$W_h = \{q \in (L^2(\Omega))^2 ; V T \in T_h, q_T \in P_{k-1,k}(T) \times P_{k,k-1}(T)\}$$

and

(12.14)
$$M_h = \{ v \in H_0^1(\Omega) ; V T \in T_h, v_{|T} \in P_{k,k}(T) \},$$

the inf-sup condition, condition (12.7),is uniformly satisfied. Morever if the solution (p,u) of (12.1) belongs to $(\operatorname{H}^k(\Omega))^n \times \operatorname{H}^{k+1}(\Omega)$ and $(\operatorname{p}_h,\operatorname{u}_h)$ is the solution of (12.8) for W_h and M_h as in (12.13) and (12.14),there exists a constant C, independent of h, such that

(12.15)
$$\| p - p_h \|_{0,\Omega} + \| u - u_h \|_{1,\Omega} \le Ch^k (\| p \|_{k,\Omega} + \| u \|_{k+1,\Omega}).$$

13 - EXAMPLES OF DUAL MIXED METHODS FOR A DIRICHLET PROBLEM

A dual mixed variational formulation of the model problem, find a solution (p,u) of system (4.4) - (4.5) with the homogeneous Dirichlet boundary condition u=0 on Γ , is given by characterizing the pair (p,u) as the solution of

$$(13.1) \begin{cases} (\underline{p}, u) \in H(\operatorname{div}; \Omega) \times L^{2}(\Omega) , \\ \\ a(\underline{p}, \underline{q}) + b(\underline{q}, u) = 0 \end{cases} \qquad \text{for all } \underline{q} \in H(\operatorname{div}; \Omega) , \\ \\ b(\underline{p}, v) = -\int_{\Omega} fv \, dx \qquad \text{for all } v \in L^{2}(\Omega) , \end{cases}$$

where

(13.2)
$$a(\underline{p},\underline{q}) = \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} p_{j} q_{i} dx$$

and

(13.3)
$$b(q,v) = \int_{\Omega} v \, div \, q \, dx.$$

We are again in the context of Chapter III with now

(13.4)
$$W = H(\operatorname{div};\Omega), \qquad || g ||_{W} = || g ||_{H(\operatorname{div};\Omega)},$$

and

(13.5)
$$M = L^2(\Omega), \| v \|_{M} = \| v \|_{0,\Omega}.$$

The subspace V of functions $q \in W$ for which b(q,v) = 0 for all $v \in M$ is none other than

(13.6)
$$V = \{q \in H(div;\Omega) ; div q = 0\}.$$

With the given hypotheses concerning the coefficients a_{ij} , the bilinear form $a(\cdot, \cdot)$ is continuous on W × W and V-elliptic with constant of ellipticity $\alpha=\alpha_d$, cf. (4.8). The bilinear form $b(\cdot, \cdot)$ is continuous on W × M and the inf-sup condition is satisfied:

(13.7)
$$\inf_{\{v \in M ; \|v\|_{0,\Omega} - 1\}} \sup_{\{g \in W ; \|g\|_{H(\operatorname{div};\Omega)} = 1\}} b(g,v) > 0.$$

To establish (13.7) we shall use the auxillary problem. For each $v \in L^2(\Omega)$ denote by ϕ_v the unique solution in $H^1_0(\Omega)$, of $-\Delta\phi_v = v$ on Ω . Then $q_v = \operatorname{grad} \phi_v$ belongs to $H(\operatorname{div};\Omega)$ and $b(q_v,v) = \|v\|_{0,\Omega}^2$. Since $\|q_v\|_{H(\operatorname{div};\Omega)} \le \|\phi_v\|_{1,\Omega}^2 + \|v\|_{0,\Omega} \le C \|v\|_{0,\Omega}$, inequality (13.7) follows. Thus we can apply Theorem 10.1 to obtain the existence and uniqueness of the solution (p,u) of (13.1).

In order to construct a dual mixed approximation of the solution (p,u) of (13.1), we introduce a finite dimensional subspace W_h of $H(\text{div};\Omega)$ and a finite dimensional subspace M_h of $L^2(\Omega)$. From Theorems 10.2 and 10.4 we deduce:

 $\underline{\textit{THEOREM 13.1}}$. Suppose that the subspaces \mathtt{W}_h and \mathtt{M}_h are compatible in the sense that the discrete inf-sup condition is satisfied :

$$\beta_{h} = \inf_{\{v_{h} \in M_{h} ; \|v_{h}\|_{0,\Omega} = 1\}} \sup_{\{q_{h} \in W_{h} ; \|q_{h}\|_{H(\operatorname{div};\Omega)} = 1\}} b(q_{h}, v_{h}) > 0.$$

Then the problem of finding a pair (p_h, u_h) satisfying

$$\begin{cases} (\underline{p}_h, u_h) \in W_h \times M_h , \\ a(\underline{p}_h, \underline{q}_h) + b(\underline{q}_h, u_h) = 0 \\ b(\underline{p}_h, v_h) \end{cases} \qquad \text{for all } \underline{q}_h \in W_h , \\ b(\underline{p}_h, v_h) \qquad = -\int_{\Omega} f v_h \ dx \qquad \text{for all } v_h \in M_h , \end{cases}$$

has a unique solution. Furthermore, for each family of pairs of subspaces $(W_h^{},M_h^{})$ such that for each $q \in W_h$ we have

$$(13.10) \qquad (\int_{\Omega} v_h \operatorname{div} q_h \operatorname{dx} = 0, \text{ for all } v_h \in M_h) \qquad \Longrightarrow \qquad \operatorname{div} q_h = 0$$

and such that (13.8) is satisfied uniformly in the sense that $\beta=\inf\beta_h>0$, there exists a constant C, independent of h, such that

(13.11)
$$\| p - p_h \|_{H(\operatorname{div};\Omega)} + \| u - u_h \|_{0,\Omega} \le C \left\{ \inf_{\substack{q_h \in W_h \\ v_h \in M_h}} \| p - q_h \|_{H(\operatorname{div};\Omega)} + \inf_{\substack{q_h \in W_h \\ v_h \in M_h}} \| u - v_h \|_{0,\Omega} \right\}.$$

Hypothesis (13.10) expresses the fact that the space V_h ,

(13.12)
$$V_h = \{q_h \in W_h ; V v_h \in M_h, b(q_h, v_h) = 0\},$$

is a subspace of the space V defined by (13.6). This hypothesis assures the uniform ellipticity of $a(\cdot,\cdot)$ on V_h , where the norm on V_h is the $H(\text{div};\Omega)$ norm.

Given a triangulation T_h by n-simplexes T of a domain $\bar{\Omega}$, supposed to be polyhedral, and a positive integer k, we take finite dimensional subspaces $W_h = W_h^{(k)}$ and $M_h = M_h^{(k)}$ given as follows:

(13.13)
$$W_{h} = \{ \underline{q} \in H(\operatorname{div};\Omega) ; \forall T \in T_{h}, \underline{q}_{T} \in D_{k}(T) \},$$

and

(13.14)
$$M_h = \{ v \in L^2(\Omega) ; V T \in T_h, v_T \in P_{k-1}(T) \},$$

where $D_k(T)$ denotes the space of restrictions to T of the functions in the space D_k introduced in (6.5). (With the notation (6.17) introduced in Chapter II we have $W_h = E_h^{(k)}$). These examples were proposed and analyzed in RAVIART & THOMAS (1977) in the 2 dimensional case and in NEDELEC (1979) in the 3-dimensional case. In the literature the space $W_h \times M_h$ or $W_h^{(k)} \times M_h^{(k)}$ is referred to as the Raviart-Thomas space of index k (or sometimes k-1; cf. Remark (6.2)) or, in the 3-dimensional case, as the Raviart - Thomas - Nedelec space.

THEOREM 13.2. For each positive integer k, the subspaces W_h and M_h defined by (13.13) and (13.14) satisfy the discret inf-sup condition (13.8). Moreover, if $\{T_h; h \in H\}$ is a regular family of triangulations, there is a constant C independent of h such that if the solution (p,u) of (13.1) is such that $(p,u) \in (H^k(\Omega))^n \times H^k(\Omega)$ and div $p \in H^k(\Omega)$ and if (p_h,u_h) is the solution of (13.9) associated with (13.13) and (13.14), then we have the error estimate

$$(13.15) \| \underline{p} - \underline{p}_h \|_{H(\operatorname{div};\Omega)} + \| \underline{u} - \underline{u}_h \|_{0,\Omega} \le \operatorname{Ch}^k(\|\underline{u}\|_{k,\Omega} + \|\underline{p}\|_{k,\Omega} + \|\operatorname{divp}\|_{k,\Omega}).$$

<u>Proof</u>. With W_h and M_h as given in (13.13) and (13.14) we clearly have that div $\underline{q}_h \in M_h$ for each $\underline{q}_h \in W_h$ and thus that $V_h \subset V$. To establish the discrete inf-sup condition we apply Theorem 10.3. Take for subspace of $W = H(\operatorname{div};\Omega)$ the space $W = (H^1(\Omega))^n$ and for continuous linear operator π_h from W into W_h , the operator $\pi_h = E_h^{(k)}$, given by (6.2) together with (6.14). Then (6.15) guarantees that (10.9) is satisfied. Let R be any continuous linear map from $L^2(\Omega)$ into $(H^1(\Omega))^n$ satisfying

$$(13.16)$$
 div(Rv) = v,

and let R_h from M_h into $W = (H^1(\Omega))^n$ be defined by

$${}^{R}h^{v}_{h} = {}^{R}v_{h}$$
 for all $v_{h} \in M_{h}$.

We note that with sufficient regularity the function $q_v \in H(\operatorname{div};\Omega)$ used to establish the inf-sup condition (13.7) belongs to $(H^1(\Omega))^n$ and Rv could be taken to be q_v . To avoid having to invoke regularity results for Dirichlet problems in polyhedral domains which would necessitate certain restrictions on the geometry of Ω , it suffices to extend the functions $v \in L^2(\Omega)$ by zero to functions $\tilde{v} \in L^2(\Omega)$ where Ω is a ball containing Ω , to solve the Dirichlet problem $\Delta \phi_v = \tilde{v}$ in Ω , $\phi_v = 0$, and to define Rv to be the restriction to Ω of grad ϕ_v .

To verify (10.8) we note there is a constant $^{\text{C}}_{O}$ depending only on the ball $^{\text{O}}$ containing $^{\Omega}$, such that

$$\| \mathbf{R}_{\mathbf{h}} \mathbf{v}_{\mathbf{h}} \|_{1,\Omega} \le \mathbf{C}_{O} \| \mathbf{v}_{\mathbf{h}} \|_{0,\Omega}.$$

Thus
$$\frac{\frac{b(R_{h}v_{h},v_{h})}{\|R_{h}v_{h}\|_{W}} = \frac{\|v_{h}\|_{0,\Omega}^{2}}{\|R_{h}v_{h}\|_{1,\Omega}} \geq \frac{1}{C_{O}} \|v_{h}\|_{0,\Omega}.$$

Since $\sup_{\{\psi\in\mathbb{W}\ ;\ \|\psi\|_{\widetilde{W}}=1\}}$ $b(\psi,v_h)\leq\|\operatorname{div}\psi\|_{0,\Omega}\|v_h\|_{0,\Omega}\leq\|v_h\|_{0,\Omega}$, we have that (10.8) holds with constant $C=1/C_O$, and (13.8) now follows from Theorem 10.3.

Next suppose that $\{T_h; h\in H\}$ is a regular family of triangulations. We have just seen that R_h may be defined with norm independent of h, and Theorem 6.3 guarantees that for the operator $\pi_h = E_h^{(k)}$ there is a constant C independent of h such that

$$\|\pi_{h} g\|_{H(\operatorname{div};\Omega)} \le C \{h \|g\|_{1,\Omega} + \|\operatorname{div} g\|_{0,\Omega}\}.$$

Thus the discrete inf-sup condition is satisfied uniformly and we have the error bound (13.11). Then to estimate the error of approximation inf $\|p-q_h\|_{H(\operatorname{div};\Omega)}$, we take $q_h=E_h^{(k)}(p)$ and use (6.19) and (6.20).

<u>Remark 13.1</u>. When the coefficients a_{ij} are sufficiently regular, for example when $a_{ij} \in C^k$, the regularity hypothesis $p \in (\operatorname{H}^k(\Omega))^n$ and div $p \in \operatorname{H}^{(k)}(\Omega)$ is enough to give the error bound (13.15). If we suppose only that $p \in (\operatorname{H}^k(\Omega))^n$ we may easily adapt the last part of the demonstration of Theorem 13.2 to obtain

or with a more careful study we obtain

If in addition we suppose the regularity of the Dirichlet problem in Ω , a hypothesis which is satisfied once the open polyhedral domain Ω is convex, we have

(13.18)
$$\| u - u_h \|_{0,\Omega} \le Ch^k \| u \|_{k^*,\Omega}$$
 with $k^* = \max(k,2)$.

These results are demonstrated in FALK & OSBORN (1980, Section 3). \Box

<u>Remark 13.2</u>. Even with no regularity hypothesis for the solution (p,u) one may show that the dual mixed method applied in the context of Theorem 13.2 is convergent:

$$\lim_{h\to 0} p = p \qquad \text{in } H(\text{div};\Omega) , \qquad \lim_{h\to 0} u_h = u \quad \text{in } L^2(\Omega). \quad \square$$

Remark 13.3 . With $W_h = W_h^{(k)}$ as given in (13.13) and with $M_h = M_h^{(k')}$ as given in (13.4) but with k' not necessarily equal to k, the compatibility hypothesis is satisfied once $k' \le k$. Error estimates, optimal in $W \times M$, are obtained when k' = k. If $1 \le k' < k$ the subspace V_h is no longer a subspace of V; however, the method remains convergent.

<u>Remark 13.4</u>: let \hat{T} be the reference n-simplex and for each $T \in T_h$, let F_T denote an affine bijection of \hat{T} onto T with Jacobian J_T . The spaces W_h and M_h given by (13.13) and (13.14) may also be defined as follows:

$$(13.13 \text{ bis}) \ \textbf{W}_{h} = \{ \textbf{q} \in \textbf{H}(\textbf{div}; \Omega) \ ; \ \textbf{V} \ \textbf{T} \in \textbf{T}_{h}, \ \big| \ \textbf{J}_{T} \ \big| \ \textbf{DF}_{T}^{-1} \ (\textbf{q}_{\big| T} \ \textbf{o} \ \textbf{F}_{T}) \ \in \textbf{D}_{k}(\hat{\textbf{T}}) \} \ ,$$

It is with these definitions that the theory is extended to the case of curved finite elements where T is the image of \hat{T} under a C^1 bijection not required to be affine. To show the compatibility relation in this case, we construct for $v_h \in M_h$, the function $\underline{q}_h \in W_h$ such that

$$\operatorname{div} \, \operatorname{q}_{h} = \theta_{h} \operatorname{v}_{h},$$

where $\theta_h \in L^{\infty}(\Omega)$ is defined on each $T \in T_h$ by $\theta_h(x) = \sup_{y \in T} \frac{\left|J_T(y)\right|}{\left|J_T(x)\right|}$, for each $x \in T$.

<u>Remark 13.5</u>. In the context of Theorem 13.2 it is obvious that for $g_h \in W_h$ we have div $g_h \in M_h$. Moreover, the proof of Theorem 13.2 shows that for each $v_h \in M_h$ there exists a function $g_h = \pi_h R_h v_h \in W_h$ such that div $g_h = v_h$. Thus in this case we have the characterization

(13.19)
$$\mathsf{M}_{h} = \{ v_{h} \in L^{2}(\Omega) ; \exists q_{h} \in W_{h} \text{ with div } q_{h} = v_{h} \}.$$

Such is not the case for the generalizations described in Remarks 13.3 and 13.4. $\hfill\Box$

For the dual mixed method analyzed above, given a triangulation T_h of $\overline{\Omega}$ by n-simplexes, the vector functions in $W = H(\text{div};\Omega)$ are approximated in the space W_h defined by (13.13) which is the space $E_h^{(k)}$ defined by (6.17) in Chapter 2. An alternative method is obtained by instead approximating the vector functions by functions in the space $E_h^{\star(k)}$ also defined in Chapter 2. For each integer $k,k\geq 2$, take $W_h=W_h^{(k)}$ and $M_h=M_h^{(k)}$ as follows:

(13.20)
$$W_h = \{q \in H(div;\Omega) : V T \in T_h, q_{|T} \in D_k^*(T)\},$$

(13.21)
$$M_h = \{v \in L^2(\Omega) : V T \in T_h, v_T \in P_{k-2}(T)\},$$

where $D_k^*(T) = (P_{k-1}(T))^n$; cf. (7.1). These approximation spaces were introduced and the corresponding methods analyzed in BREZZI, DOUGLAS & MARINI (1985) in dimension n=2 and in BREZZI, DOUGLAS, DURAN & FORTIN (1987) in dimension n=3. In the literature the space $W_h \times M_h$ or $W_h^{(k)} \times M_h^{(k)}$ is referred to as the Brezzi-Douglas-Marini space of index k (or sometimes k-1, cf. Remark 6.2), or in the 3 dimensional case the Brezzi-Douglas-Duran-Fortin space.

As for the Raviart-Thomas (-Nedelec) spaces it is clear that if $\underline{p}_h \in W_h$, then div $\underline{p}_h \in M_h$, and one may show that the discrete inf-sup condition (13.8) is satisfied. In fact we again have the characterization (13.19). For a regular Dirichlet problem in Ω we have the error estimates

$$\begin{split} & \left\| \left\| \mathbf{p} - \mathbf{p}_{h} \right\|_{0,\Omega} \leq ch^{k} \left\| \mathbf{p} \right\|_{k,\Omega} , \\ & \left\| \left\| \mathbf{u} - \mathbf{u}_{h} \right\|_{0,\Omega} \leq ch^{k-1} \left\| \mathbf{u} \right\|_{k-1,\Omega} , \\ & \left\| \operatorname{div}(\mathbf{p} - \mathbf{p}_{h}) \right\|_{0,\Omega} \leq ch^{k-1} \left\| \operatorname{div} \mathbf{p} \right\|_{k-1,\Omega} . \end{split}$$

For each of the dual mixed methods presented above for the case where $T_{\rm n}$ is a triangulation by n-simplexes, there is a variant adapted to the case where $T_{\rm h}$ is a triangulation by rectangles with sides parallel to the axes when n=2 or by rectangular solids with sides parallel to the coordinate

planes when n=3. The description of these methods and the corresponding analyses may be found in the references cited above for the associated method for a triangulation by simplexes. Here we shall simply state for the case n=2 the analogue of Theorem 13.2 when $T_{\rm h}$ is a triangulation of $\bar{\Omega}$ by rectangles with sides parallel to the axes.

THEOREM 13.3 : For each positive integer k, the subspaces $W_h = W_h^{(k)}$ and $M_h = W_h^{(k)}$ defined by

$$(13.24) W_h = \{ \underline{q} \in H(\operatorname{div};\Omega) ; V T \in T_h, \underline{q}_{T} \in D_k^{\square}(T) \},$$

and

(13.25)
$$M_h = \{ v \in L^2(\Omega) ; V T \in T_h, v_{|T} \in P_{k-1,k-1}(T) \},$$

with $D_k^{\square} = P_{k,k-1} \times P_{k-1,k}$, satisfy the inf-sup condition (13.8). Moreover, if $(T_h; h \in H)$ is a regular family of triangulations, there is a constant C, independent of h, such that if the solution (p,u) of (13.1) is such that $(p,u) \in (H^k(\Omega))^n \times H^k(\Omega)$ and div $p \in H^k(\Omega)$ and if (p_h, u_h) is the solution of (13.9) associated with (13.24) (13.25), we have the error estimate

$$\|\mathbf{p} - \mathbf{p}_h\|_{\mathbf{H}(\mathbf{div};\Omega)} + \|\mathbf{u} - \mathbf{u}_h\|_{\mathbf{0},\Omega} \le \mathbf{Ch}^k(\|\mathbf{u}\|_{\mathbf{k},\Omega} + \|\mathbf{p}\|_{\mathbf{k},\Omega} + \|\mathbf{div}\mathbf{p}\|_{\mathbf{k},\Omega}). \square$$

Other dual mixed methods are presented and analysed in BREZZI, DOUGLAS, FORTIN, & MARINI (1987) and in NEDELEC (1986); cf. Remark 7.2.

14 - EXAMPLES OF MIXED METHODS FOR A NEUMANN PROBLEM

We consider the model problem, given $f \in L^2(\Omega)$ with $\int_{\Omega} f \, dx = 0$, find u (to within a constant) and p satisfying (4.4) and (4.5) with the homogeneous boundary condition $p \cdot \nu = 0$ on Γ .

A primal mixed variational formulation of this problem is obtained by characterizing the pair (p,u) as the solution of

$$(p,u) \in (L^2(\Omega))^n \times H^1(\Omega)/R ,$$

$$(14.1) \begin{cases} a(p,q) + b(q,u) = 0 & \text{for all } q \in (L^2(\Omega))^n , \\ b(p,v) & = -\int_{\Omega} \text{ for dx} & \text{for all } v \in H^1(\Omega)/R , \end{cases}$$

where the bilinear forms $a(\cdot, \cdot)$ and $b(\cdot, \cdot)$ are given by (12.2) and (12.3). Recalling the notation for the abstract problem of Section 10, here we put

(14.2)
$$W = (L^2(\Omega))^n$$
, $\|q\|_{W} = \|q\|_{0,\Omega}$

and

(14.3)
$$M = H^{1}(\Omega)/R$$
, $\|v\|_{M} = \|v\|_{1,\Omega}$.

With W_h and M_h denoting finite dimensional subspaces of $(L^2(\Omega))^n$ and $H^1(\Omega)/R$, respectively, we obtain the analogue of Theorem 12.1, with no modification of the statement, and also the analogue of Theorem 12.2 if we specify the subspaces W_h and M_h as follows:

(14.4)
$$W_h = (q \in (L^2(\Omega))^n ; V T \in T_h, q_T = (P_{k-1}(T))^n)$$

and

(14.5)
$$M_h = \{v \in H^1(\Omega)/R ; V T \in T_h, v_{|T} \in P_k(T)\},$$

for some $k \ge 1$. In the case of a triangulation by rectangles, with the appropriate modification of the subspaces W_h and M_h , the analogue of Theorem 12.3 is obtained.

Thus for primal mixed methods there is no difficulty in adapting the theory for the Dirichlet problem to obtain that for the Neumann problem. The construction of examples may even be easier for the Neumann problem as the functions of M_h need no longer satisfy, à priori, the zero trace condition on Γ . The situation for dual mixed methods is quite the opposite.

A dual mixed variational formulation of the homogeneous Neumann problem is obtained by characterizing the pair (\underline{p},u) as the solution of

$$(14.6) \begin{cases} (\underline{p},u) \in H_0(\operatorname{div};\Omega) \times L^2(\Omega)/\mathbb{R} , \\ \\ a(\underline{p},\underline{q}) + b(\underline{q},u) = 0 \\ \\ b(\underline{p},v) = -\int_{\Omega} \operatorname{fv} \, dx & \text{for all } \underline{v} \in L^2(\Omega)/\mathbb{R} , \end{cases}$$

where the bilinear forms $a(\cdot, \cdot)$ and $b(\cdot, \cdot)$ are, given by (13.2) and (13.3). For convenience, $L^2(\Omega)/R$ will be identified with the subspace of functions $v \in L^2(\Omega)$ such that $\int_{\Omega} v \ dx = 0$. Here the spaces W and M and their norms are defined by

(14.7)
$$W = H_0(\operatorname{div};\Omega),$$
 $\|g\|_{W} = \|g\|_{H(\operatorname{div};\Omega)},$

and

(14.8)
$$M = \{v \in L^2(\Omega) ; \int_{\Omega} v \, dx = 0\}, \qquad ||v||_{M} = ||v||_{0,\Omega}.$$

The inf-sup condition analogous to (13.7) is also established with the aid of the dual problem. For each $v \in M$ denote by ϕ_v the unique element of $H^1(\Omega)/R$ satisfying $\Delta\phi_v = v$ in Ω and $\frac{\partial\phi_v}{\partial\nu} = 0$ on Γ . The function $q_v = \operatorname{grad} \phi_v$ belongs to W and satisfies div $q_v = v$ in Ω .

With W_h and M_h denoting finite dimensional subspaces of $H_0(\text{div};\Omega)$ and M_h respectively, a demonstration similar to that given for Theorem 3.1 shows that if the discret inf-sup condition is satisfied, there exists a unique solution of the following problem : find (p_h, u_h) satisfying

$$\begin{cases} (\underline{p}_h, u_h) \in W_h \times M_h ; \\ a(\underline{p}_h, \underline{q}_h) + b(\underline{q}_h, u_h) = 0 & \text{for all } \underline{q}_h \in W_h ; \\ b(\underline{p}_h, v_h) = -\int_{\Omega} f v_h dx, & \text{for all } v_h \in M_h. \end{cases}$$

To show that the discret inf-sup condition holds for a pair of subspaces \mathbb{W}_h and \mathbb{W}_h , however, is a more delicate matter here than for the Dirichlet problem. It suffices, of course, to be able to associate to each $\mathbf{v}_h \in \mathbb{W}_h$ a $\mathbf{q}_h \in \mathbb{W}_h$ with div $\mathbf{q}_h = \mathbf{v}_h$, but this is made more complicated here by the requirement that $\mathbf{q}_h \in \mathbb{W}_h \subset \mathbb{W}$ have a vanishing normal trace on Γ . Nonetheless, one may show that the discrete inf-sup condition is satisfied when, for a triangulation T_h by n-simplexes, the subspaces \mathbb{W}_h and \mathbb{W}_h are defined as follows:

(14.10)
$$W_h = \{q \in H_0(div;\Omega) : V T \in T_h, q_T \in D_k(T)\}$$

and

(14.11)
$$M_h = \{ v \in L^2(\Omega) ; \int_{\Omega} v \, dx = 0 \text{ and } V T \in T_h, v_{\mid T} \in P_{k-1}(T) \},$$

for some positive integer k. We give the demonstration for the case $k\!=\!1$:

Let $v_h \in M_h$; thus v_h is constant on each n-simplex T and has zero average over Ω . Denote by Φ_h the space of functions $\phi_h \in H^1(T_h)$ with $\int_{\Omega} \phi_h \, \mathrm{d}x = 0$, with $\phi_h|_T$, for each $T \in T_h$, an affine function, and with ϕ_h continuous across the interior faces at the barycenters of the faces T'common to two n-simplexes of T_h . For each $v_h \in M_h$, introduce the non-conforming solution $\phi_h(v_h)$ of the dual Neumann problem : find ϕ_h satisfying

$$\begin{cases} \phi_{h} \in \Phi_{h} ; \\ \sum_{T \in T_{h}} \int_{T} \operatorname{grad} \phi_{h} \cdot \operatorname{grad} \psi_{h} dx = \int_{\Omega} v_{h} \psi_{h} dx, & \text{for all } \psi_{h} \in \Phi_{h}. \end{cases}$$

Let $\underline{q}_h = \underline{q}_h(v_h)$ be the function of $(L^2(\Omega))^n$ whose restriction to T, for each $T \in T_h$, coincides with the function of $D_1(T)$ having normal trace on ∂T equal to that of $\phi_h|_T$. One can verify that \underline{q}_h belongs to W_h and that div $\underline{q}_h = v_h$. One may further establish that, for a uniformly regular family of triangulations $\{T_h; h \in \mathcal{H}\}$, there exists a constant C, independent of h, such that

(14.13)
$$\|\psi_{h}\|_{0,\Omega} \le C \left(\sum_{T \in T_{h}} \|\psi_{h}\|_{1,T}^{2}\right)^{\frac{1}{2}},$$

a result which amounts to a generalization of the Friedrichs-Poincaré inequalities to the nonconforming case. One thus obtains the existence of a (new) constant C, independent of h, such that for the preceding construction we have

$$(14.14) \qquad \left\| g_{h}(v_{h}) \right\|_{H(\operatorname{div};\Omega)} \leq C \left\| v_{h} \right\|_{O,\Omega},$$

thereby showing that the discrete inf-sup condition is satisfied uniformly. Hence we may deduce that there is a constant C, independent of h, such that

$$\left\| \left\| \mathbf{p} - \mathbf{p}_{h} \right\|_{H(\operatorname{div};\Omega)} + \left\| \left\| \mathbf{u} - \mathbf{u}_{h} \right\|_{0,\Omega} \le \operatorname{Ch} \left(\left\| \mathbf{u} \right\|_{1,\Omega} + \left\| \mathbf{p} \right\|_{1,\Omega} + \left\| \operatorname{div} \mathbf{p} \right\|_{1,\Omega} \right).$$

To verify the discrete inf-sup condition we have avoided using the regularity of the Neumann problem on polyhedral domains. Thus the above argument may be extended to the case of a problem with a Dirichlet condition on part of the boundary and a Neumann condition on the remaining part. For this type of problem, an extension of the above result to the case k>1 may be found in RAVIART & THOMAS (1977a).

The uniform regularity hypothesis for the family of triangulations can be removed using the regularity of the solution of the Neumann problem for the Laplace equation in a polyhedral domain; cf. THOMAS (1987b).

15 - EXAMPLES OF MIXED FINITE ELEMENT METHODS FOR A ROBIN'S BOUNDARY CONDITION

In this section we discuss the adaptation of mixed formulations to the following model problem : for a given function $f \in L^2(\Omega)$ and function $q \in L^2(\Gamma)$, find a solution (p,u) of system (4.4) - (4.5) with the Robin's boundary condition

(15.1)
$$p \cdot \nu + tu = g$$
 on Γ ,

where t is a continuous, strictly positive function on Γ . (Thus t is bounded above and away from 0.)

A primal, mixed, variational formulation of this problem is obtained by characterizing the pair (p,u) as the solution of

$$(p,u) \in (L^2(\Omega))^n \times H^1(\Omega) ,$$

$$(15.2) \begin{cases} a(p,q) + b(q,u) = 0 & \text{for all } q \in (L^2(\Omega))^n , \\ b(p,v) + d(u,v) = -\int_{\Omega} fv \, dx - \int_{\Gamma} gv \, d\sigma & \text{for all } v \in H^1(\Omega) , \end{cases}$$

where the bilinear forms $a(\cdot, \cdot)$ and $b(\cdot, \cdot)$ are defined as for the Dirichlet and Neumann problems by (12.2) and (12.3) and where the bilinear form $d(\cdot, \cdot)$ is defined on $H^1(\Omega) \times H^1(\Omega)$ by

(15.3)
$$d(u,v) = -\int_{\Gamma} tuv \ d\sigma.$$

Problem (15.2) is the first example we have encountered of the generalization (10.16) developed in the abstract setting in Remark 10.7 with here

$$(15.4) \qquad W = (L^2(\Omega))^n$$

and

$$(15.5) \qquad \mathbf{M} = \mathbf{H}^{1}(\Omega).$$

The numerical analysis of a primal mixed finite element method derived from formulation (15.2) is not difficult when to a triangulation $T_{\hat{\mathbf{h}}}$ of $\hat{\Omega}$ by n-simplexes T are associated the finite dimensional subspaces

(15.6)
$$W_h = \{q \in (L^2(\Omega))^n ; V T \in T_h, q_T \in (P_{k-1}(T))^n\}$$

and

(15.7)
$$M_h = \{v \in H^1(\Omega) ; \forall T \in T_h, v_{|T} \in P_k(T)\},$$

where k is a positive integer . We remark nonetheless that in this situation the inf-sup condition can never be satisfied as the functions constant on each n-simplex belong to M_h .

With the dual mixed variational formulation, the solution (p,u) of system (4.4) - (4.5) with boundary condition (15.1) is characterized as the solution of

$$(15.8) \begin{cases} (\underline{p}, u) \in H(\operatorname{div}; \Omega) \times L^{2}(\Omega) , \\ a(\underline{p}, \underline{q}) + b(\underline{q}, u) = \int_{\Gamma} \frac{1}{t} g \underline{q} \cdot \underline{\nu} d\sigma & \text{for all } q \in H(\operatorname{div}; \Omega) , \\ b(\underline{p}, v) = -\int_{\Omega} fv dx & \text{for all } v \in L^{2}(\Omega) , \end{cases}$$

where the bilinear form $b(\cdot, \cdot)$ is defined as for the Dirichlet and Neumann problems by (13.3) and where the bilinear form $a(\cdot, \cdot)$ is now given by

(15.9)
$$a(\underline{p},\underline{q}) = \sum_{i,j=1}^{n} \int_{\Omega} A_{ij} p_{j} q_{i} dx + \int_{\Gamma} \frac{1}{t} \underline{p} \cdot \nu \underline{q} \cdot \nu d\sigma.$$

We thus obtain a problem of form (10.1) with now

(15.10)
$$W = H(\operatorname{div}; \Omega)$$

and

(15.11)
$$M = L^2(\Omega)$$
.

There exists a continuous linear operator R from $L^2(\Omega)$ into $H(\operatorname{div};\Omega)$ satisfying $\operatorname{div}(\mathbb{R} v) = v$ for each $v \in L^2(\Omega)$. (One may use, for example, the operator R from $L^2(\Omega)$ into $(\operatorname{H}^1(\Omega))^n$ constructed in the proof of Theorem 13.2). Hence the inf-sup condition is satisfied, the norms for the spaces W and M being given by $\|\cdot\|_W = \|\cdot\|_{H(\operatorname{div};\Omega)}$ and $\|\cdot\|_M = \|\cdot\|_{0,\Omega}$.

For a triangulation $T_{\rm h}$ of Ω by n-simplexes T, take for finite dimensional subspaces $W_{\rm h}$ and $M_{\rm h}$

(15.12)
$$W_{h} = \{ \underline{q} \in H(\operatorname{div};\Omega) ; \forall T \in T_{h}, \underline{q}_{T} \in D_{k}(T) \}$$

and

(15.13)
$$M_h = \{ v \in L^2(\Omega) ; V T \in T_h, v_{|T} \in P_{k-1}(T) \},$$

where k is a positive integer . One can show with no new difficulties beyond those already treated for the problem with Dirichlet boundary condition that the problem, find a pair (p_h, u_h) satisfying

$$\left(\begin{array}{c} (\underline{p}_h, \underline{u}_h) \in \underline{W}_h \times \underline{M}_h \ , \\ \\ a(\underline{p}_h, \underline{q}_h) + b(\underline{q}_h, \underline{u}_h) = \int_{\Gamma} \frac{1}{t} g \ \underline{q}_h \cdot \underline{\nu} \ d\sigma \quad \text{for all } \underline{q}_h \in \underline{W}_h \ , \\ \\ b(\underline{p}_h, \underline{v}_h) = -\int_{\Omega} f \underline{v}_h \ dx \qquad \qquad \text{for all } \underline{v}_h \in \underline{M}_h \ , \\ \end{array} \right)$$

admits a unique solution (p_h, u_h) . If $\{T_h; h \in H\}$ is a regular family of triangulations, the discrete inf-sup condition is satisfied uniformly and if further the solution (p,u) is sufficiently regular, then we have

(15.15)
$$\| \mathbf{p} - \mathbf{p}_{\mathbf{h}} \|_{H(\mathbf{div};\Omega)} + \| \mathbf{u} - \mathbf{u}_{\mathbf{h}} \|_{O,\Omega} = O(\mathbf{h}^{\mathbf{k}}).$$

16 - OTHER EXAMPLES OF MIXED FINITE ELEMENT METHODS

All of the examples presented above treat only elliptic problems governed by second order partial differential equations of the form

$$-\sum_{i,j=1}^{n} \frac{\partial}{\partial x_{i}} (a_{ij} \frac{\partial u}{\partial x_{j}}) - f.$$

The generalization of these methods to problems governed by an equation of the form

$$-\sum_{i,j=1}^{n} \frac{\partial}{\partial x_{i}} \left(a_{ij} \frac{\partial u}{\partial x_{j}}\right) + \sum_{j=1}^{n} b_{j} \frac{\partial u}{\partial x_{j}} + cu = f$$

poses no difficulty for the primal mixed formulation. The case is quite different for the dual mixed formulation. A mathematical analysis of this problem can be found in DOUGLAS & ROBERTS (1982) and (1985). A generalization of these methods to the quasilinear problem where the coefficients in the governing equation depend on u is analyzed in MILNER (1985). Several methods adapted to treat numerically the case where the the convection term is large in comparison with the diffusion term have been studied in JAFFRE (1984), JAFFRE & ROBERTS (1985), JOLY (1982), and THOMAS (1987a). Another generalization of mixed methods to problems governed by an equation of the form

$$-\sum_{i,j=1}^{n} \frac{\partial}{\partial x_{i}} \left(a_{ij} \frac{\partial u}{\partial x_{j}} \right) = f$$

with coefficients a_{ij} which can degenerate has been studied by LE ROUX (1982).

Estimations of the error made by a mixed method approximation have been given in norms other than the L²-norm and the H(div)-norm. Estimates for u-u_h in the L $^{\infty}$ -norm have been given in SCHOLZ (1977), (1983), DOUGLAS & ROBERTS (1985), and GASTALDI & NOCHETTO (1987). Negative norm estimates in H $^{-S}(\Omega)$, for s a nonnegative integer, are given for u-u_h, p-p_h, and div(p-p_h) in DOUGLAS & ROBERTS (1985). These estimates have been exploited in DOUGLAS & MILNER (1985) to obtain super convergence results.

Mixed methods have been "generalized" to 4^{th} order problems such as those governed by an equation of the form $\Delta^2 u = f$, where we write "generalized" with quotation marks to indicate that these methods were introduced, for the most part, before those presented for the 2^{nd} order problem and without the use of the general abstract theory for problems with Lagrangian multipliers. We cite in particular the work of BRAMBLE & FALK (1983); BREZZI & RAVIART (1978); CIARLET & RAVIART (1974), HELLAN (1967); HERRMANN (1967), (1983); JOHNSON (1973); MERCIER (1974); MIYOSHI (1973); QUARTERONI (1980a), SCAPOLLA (1980); and SCHOLZ (1978). A precise analysis of the best known of these methods can be found in FALK & OSBORN (1980). An excellent exposition is also given in CIARLET (1978 Chapter 7).

A generalization to problems governed by an equation of the form $(-\Delta)^m_{u=f}$ for an integer m,m ≥ 2 is proposed in BRAMBLE & FALK (1985).

17 - SOLUTION OF THE LINEAR SYSTEM

We conclude this chapter with some remarks concerning the solution of the algebraic system resulting from a mixed method or, in particular, from a dual mixed method. We shall restrict our attention to the case of the model problem (4.1), (4.2) as in Sections 12 through 14. We assume that the continuous problem is written as a problem with Lagrangian multipliers in the form (8.1), and that finite dimensional subspaces W_h and M_h of the spaces W_h and W_h of the spaces W_h

$$\begin{cases} \psi_{h} \in W_{h}, & \lambda_{h} \in M_{h}; \\ a(\phi_{h}, \psi_{h}) + b(\psi_{h}, \lambda_{h}) = (f, \psi_{h}), & \text{for all } \psi_{h} \in W_{h}; \\ b(\phi_{h}, \mu_{h}) & = (g, \mu_{h}) & \text{for all } \mu_{h} \in M_{h}; \end{cases}$$

and we further suppose that $a(\cdot, \cdot)$ is symmetric. Once bases W_h and M_h have been chosen we may write the algebraic problem associated with (17.1) as

(17.2)
$$\begin{pmatrix} A & B^* \\ B & O \end{pmatrix} \begin{pmatrix} \phi_h \\ \lambda_h \end{pmatrix} = \begin{pmatrix} f_h \\ g_h \end{pmatrix},$$

where A is a symmetric, dim $W_h \times \dim W_h$ matrix such that $(A\psi_h, \psi_h') = a(\psi_h, \psi_h')$

for all (ψ_h, ψ_h') in $W_h \times W_h$, and B is a dim $M_h \times$ dim W_h matrix such that $(B\psi_h, \mu_h) = b(\psi_h, \mu_h)$ for all (ψ_h, μ_h) in $W_h \times M_h$. The functions f_h and g_h are the orthogonal projections of f_h onto W_h and g_h onto M_h , respectively.

Theorem 10.2 guarantees that the matrix L,

$$(17.3) L = \begin{pmatrix} A & B^* \\ B & 0 \end{pmatrix},$$

of order N \times N, N = dim W_h dim M_h, is invertible; however, it is not positive definite. Hence direct solution of the system (17.2) is generally not feasible. On the other hand A is here positive definite. So, theoretically it is always possible to eliminate ϕ_h from the system,

(17.4)
$$\phi_{h} = A^{-1}(f_{h} - B^{*}\lambda_{h}),$$

to obtain

(17.5)
$$BA^{-1}B^* \lambda_h = -g_h + BA^{-1} f_h.$$

However, there is no reason, a priori, why A⁻¹ should be sparse.

For the primal mixed method, the space W is $(L^2(\Omega))^n$. So, among the degrees of freedom determining an element of W_h there is none needed to enforce a continuity requirement. Thus A (and hence A^{-1}) is block diagonal, and the structure of the linear system (17.5) is the same as that associated with the conforming problem (2.21).

For the dual mixed method the situation is different. Thus while we can not give here an exhaustive treatment of the solution of linear systems of the form (17.2), we would like to mention briefly some techniques that have been used when the system (17.2) results from a dual mixed method.

Remark 17.1 The algebraic system associated with the primal or dual hybrid methods discussed in the following chapter can, under restrictions similar to those indicated above, take the form (17.2). However, as with the primal mixed method, the form of A is such that the linear system that must be solved is that of a more standard problem, again as for a conforming problem for the dual hydrid method, as for a nonconforming problem for the primal hydrid method. Thus no section analogous to this one will be given in Chapter V.

Conjugate Gratient Method. Even for the dual mixed method where A^{-1} is full, the system (17.5) can be solved using conjugate gradient iteration as A^{-1} never actually has to be computed. Whenever $A^{-1}\psi$ is needed one can solve the system $A\psi = \psi$ using, for example, Gaussian elimination as A itself is sparse. However, conjugate gradient iteration without preconditioning is known to converge slowly, and in general there seems to be no obvious preconditioner, though in BROWN (1982) and in EWING & WHEELER (1983) several preconditioners are proposed and tested for the case of a triangulation by rectangles. Moreover, conjugate gradient iteration is known to be especially sensitive to round-off error.

It seems appropriate to mention that there has been much research concerning the use of the conjugate gradient method and its variants to solve matrix equations involving matrices that are not necessarily symmetric or positive definite. Thus one may consider applying such a method to solve equation (17.2) even though L is not positive definite. In JOLY (1984) several variants of the conjugate gradient method are described and their performance in solving equations of the form (17.2) coming from the implementation of a dual mixed method are compared. As pointed out in JOLY (1984), for such systems the role of the preconditioner is no longer simply to speed up the convergence but in fact to guarantee the convergence.

<u>Penalty Method</u>. This method consists of replacing (17.1) by a more regular problem : find $(\phi_{h,\epsilon},\lambda_{h,\epsilon})$ satisfying

$$\begin{cases} \phi_{h,\epsilon} \in W_h, & \mu_{h,\epsilon} \in M_h ; \\ \\ a(\phi_{h,\epsilon},\psi_h) + b(\psi_h\lambda_{h,\epsilon}) = (f,\psi_h), & \text{for all } \psi_h \in W_h ; \\ \\ b(\phi_{h,\epsilon},\mu_h) - \epsilon(\lambda_{h,\epsilon},\mu_h) = (g,\mu_h), & \text{for all } \mu_h \in M_h. \end{cases}$$

The corresponding algebraic problem then becomes

(17.7)
$$\begin{pmatrix} A & B^* \\ \\ B & -\epsilon I \end{pmatrix} \begin{pmatrix} \phi_{h,\epsilon} \\ \\ \lambda_{h,\epsilon} \end{pmatrix} = \begin{pmatrix} f_h \\ \\ g_h \end{pmatrix},$$

(17.8)
$$\lambda_{h,\epsilon} = \frac{1}{\epsilon} (B\phi_{h,\epsilon} - g_h),$$

to obtain

(17.9)
$$(\mathbf{A} + \frac{1}{\epsilon} \ \mathbf{B}^* \mathbf{B}) \ \phi_{\mathbf{h},\epsilon} = \mathbf{f}_{\mathbf{h}} + \frac{1}{\epsilon} \ \mathbf{B}^* \mathbf{g}_{\mathbf{h}}.$$

For our model problem the matrix $(A + \frac{1}{\epsilon} B^*B)$ is symmetric, positive definite, and sparse, hence (17.9) can be solved by known techniques. In BERCOVIER (1978, Theorem 3.1) such penalized problems are studied and it is shown that, under the more general hypotheses that i) and ii) of Theorem 10.2 and

(17.10)
$$\inf_{\{\psi_{h} \in W_{h}; \|\psi_{h}\|_{W_{h}} = 1\}} \{ \|a(\psi_{h}, \psi_{h})\| + \sup_{\{\mu_{h} \in M_{h}; \|\mu_{h}\|_{M_{h}} = 1\}} b(\psi_{h}, \mu_{h}) \} > 0$$

hold, there is a unique solution $(\phi_{h,\epsilon},\lambda_{h,\epsilon})$ of (17.6), and if (ϕ_h,λ_h) is the solution of (17.1) then there is a positive constant C, independent of ϵ , such that

$$(17.11) \qquad ||\phi_{h} - \phi_{h,\epsilon}||_{W_{h}} + ||\lambda_{h} - \lambda_{h,\epsilon}||_{M_{h}} \le C\epsilon.$$

The choice of ϵ may thus be made independently of h; though, as pointed out in BERCOVIER (1978), for each h there is an optimal ϵ .

This method is surely the simplest to implement, but it has one crucial flaw. The equation

$$\mathbf{B}\phi_{\mathbf{h}} = \mathbf{g}_{\mathbf{h}}$$

is not satisfied by the solution $\phi_{h,\epsilon}$. Thus the method, at least for our model problem, is mostly of historical interest.

Augmented Lagrangian Method. Under the hypothesis that $a(\cdot, \cdot)$ is symmetric and W_h -elliptic, (ϕ_h, λ_h) is a solution of (17.1) if and only if (ϕ_h, λ_h) is a saddle point of the Lagrangian Γ defined on $W_h \times M_h$ as follows:

(17.12)
$$\Gamma(\psi_h, \mu_h) = \frac{1}{2} a(\psi_h, \psi_h) + b(\psi_h, \mu_h) - f(\psi_h) - g(\mu_h).$$

The classical iterative method of Uzawa to obtain the saddle point of Γ is described by the algorithm :

Intialize:

Let
$$\lambda_h^{(0)}$$
 be an arbitrarily chosen element of M_h .

Calculate ϕ_h :

Once $\lambda_h^{(n)}$ is known, $\phi_h^{(n)}$ is defined to be the element of W_h minimizing

$$\Gamma(\psi_h, \lambda_h^{(n)})$$
; i.e. $\phi_h^{(n)}$ is the solution of

(17.13)
$$A\phi_{h}^{(n)} = f_{h} - B^{*}\lambda_{h}^{(n)}.$$

Advance λ_h :

Let $\lambda_h^{(n+1)}$ be defined by

(17.14)
$$\lambda_{h}^{(n+1)} = \lambda_{h}^{(n)} + \rho_{n} (B \phi_{h}^{(n)} - g_{h}).$$

This method in general converges quite slowly. The idea of the augmented Lagrangian method is to "penalize" the Lagrangian Γ to obtain a Lagrangian Γ for which Uzawa's method converges more rapidy. For r > 0, put

(17.15)
$$\Gamma_{\mathbf{r}}(\psi_{\mathbf{h}}, \mu_{\mathbf{h}}) = \Gamma(\psi_{\mathbf{h}}, \mu_{\mathbf{h}}) + \frac{\mathbf{r}}{2} ||\mathbf{B} \psi_{\mathbf{h}} - \mathbf{g}_{\mathbf{h}}||^{2}.$$

It is easy to see that (ϕ_h, λ_h) is a saddle point of Γ_r if and only if it is a saddle point of Γ . Thus this method is not really a penalty method as Γ does not have to tend toward ∞ to obtain the actual solution, and the augmented Lagrangian method does not have the defect of the earlier described penalty method. Applying Uzawa's method to Γ_r one obtains the algorithm:

Initialize:

Let $\lambda_h^{(0)}$ be an arbitrarily chosen element of M_h .

Calculate ϕ_h :

Once $\lambda_h^{(n)}$ is known, $\phi_h^{(n)}$ is defined to be the element of W_h minimizing

$$\Gamma_{\mathbf{r}}(\psi_{\mathbf{h}},\lambda_{\mathbf{h}}^{(\mathbf{n})})$$
; i.e. $\phi_{\mathbf{h}}^{(\mathbf{n})}$ is the solution of

(17.16)
$$\mathbf{A}_{r} \phi_{h}^{(n)} = (\mathbf{A} + r\mathbf{B}^{*}\mathbf{B}) \phi_{h}^{(n)} = f_{h} - \mathbf{B}^{*}\lambda_{h}^{(n)} + r\mathbf{B}^{*}\mathbf{g}_{h}.$$

Advance λ_h :

Let $\lambda_h^{(n+1)}$ be defined by

(17.17)
$$\lambda_{h}^{(n+1)} = \lambda_{h}^{(n)} + \rho_{n} (B \phi_{h}^{(n)} - g_{h}).$$

The augmented Lagrangian method originated with the work of HESTENES (1969) and POWELL (1969) and is studied extensively in FORTIN & GLOWINSKI (1983) where the following two results are shown (theorems 2.1 and 2.2) under the hypothesis that A is symmetric and positive definite:

- i) There is a positive constant α such that if $\alpha \leq \rho_n \leq 2r$ for each n, then the sequence $\phi_h^{(n)}$ converges to the solution ϕ_h .
- ii) If $\hat{\lambda}_h$ is the unique element of Im B satisfying (17.1) with ϕ_h , then the sequence $\lambda_h^{(n)}$ converges to $\hat{\lambda}_h + \lambda^\circ$, where λ° is the componant of $\lambda_h^{(0)}$ in Ker B*.

<u>Remark 17.2</u>. The algorithm of Uzawa, and hence that of the augmented Lagrangian, can be interpreted in the following manner: replace problem (17.1) by the virtual time dependent problem

and look for the solution (ϕ_h, λ_h) as the steady state solution of (17.18). Introducing an explicit time discretization with n^{th} time step $\rho_n = t^{n+1} - t^n$, one obtains

$$\begin{cases} a(\phi_h^{(n)}, \psi_h) + b(\psi_h, \lambda_h^{(n)}) = (f, \psi_h) & \text{for all } \psi_h \in W_h, \\ b(\phi_h^{(n)}, \mu_h) - (\frac{\lambda_h^{(n+1)} - \lambda_h^{(n)}}{\rho_n}, \mu_h) = (g, \mu_h) & \text{for all } \mu_h \in M_h, \end{cases}$$

which gives the algebraic system

$$\begin{cases} A\phi_h^{(n)} + B^* \lambda_h^{(n)} = f_h \\ \rho_n B\phi_h^{(n)} - (\lambda_h^{(n+1)} - \lambda_h^{(n)}) = \rho_n g_h \end{cases}$$

for Uzawa's method and

$$\begin{cases}
(A + rB^*B) & \phi_h^{(n)} + B^* \lambda_h^{(n)} = f_h + rB^*g_h \\
\rho_n & B & \phi_h^{(n)} - (\lambda_h^{(n+1)} - \lambda_h^{(n)}) = \rho_n & g_h
\end{cases}$$

for the augmented Lagrangian method.

<u>Remark 17.3</u>. Uzawa's algorithm can be viewed as a gradient type algorithm for the minimization of the functional J_* on M_h :

$$\begin{split} J_{*}(\mu_{h}) &= -\min_{\psi_{h} \in \mathbb{W}_{h}} \Gamma(\psi_{h}, \mu_{h}) \\ &= \frac{1}{2} (BA^{-1}B^{*}\mu_{h}, \mu_{h}) - (BA^{-1}f_{h}, \mu_{h}) + \frac{1}{2} (A^{-1}f_{h}, f_{h}) + (g_{h}, \mu_{h}), \end{split}$$

in that eliminating $\phi_h^{(n)}$ in (17.13) and (17.14),

$$\phi_{h}^{(n)} = A^{-1}(f_{h}-B^{*}\lambda_{h}^{(n)}),$$

we obtain the algorithm

Initialize:

Let $\lambda_h^{(0)}$ be an arbitrarily chosen element of M_h .

Advance λ_h :

Once $\lambda_h^{(n)}$ is known, put

$$\lambda_h^{(n+1)} = \lambda_h^{(n)} - \rho_n (BA^{-1}B^*\lambda_h^{(n)} - BA^{-1}f_h + g_h).$$

Thus the augmented Lagrangian algorithm can be viewed as a gradient type algorithm for the minimization of the dual functional J_{*r} on M_h :

$$\begin{split} &J_{*_{r}}(\mu_{h}) = -\min_{\psi_{h} \in W_{h}} \Gamma_{r}(\psi_{h}, \mu_{h}) \\ &= \frac{1}{2} \left(BA_{r}^{-1}B^{*}\mu_{h}, \mu_{h}\right) - \left(BA_{r}^{-1}(f_{h} + rB^{*}g_{h}), \mu_{h}\right) + \frac{1}{2} \left(A_{r}^{-1}(f_{h} + rB^{*}g_{h}) \right), \ f_{h} + rB^{*}g_{h}. \end{split}$$

However A_r^{-1} , like A^{-1} , is, a priori, full. Similarly, one may eliminate $\lambda_h^{(n)}$ to obtain an algorithm involving $\phi_h^{(n)}$, but again the matrix A_r^{-1} appears.

Uzawa's algorithm indicates no specific method for solving the equation determining $\phi_h^{(n)}$. As A and A_r are positive definite and sparse, either a direct method such as Gaussian elimination or an iterative method may be used to determine $\phi_h^{(n)}$ in (17.13) or (17.16). The Arrow-Hurwitz algorithm is an algorithm that may be interpreted as being a variant of the Uzawa algorithm in which the method for calculating $\phi_h^{(n)}$ is specified:

Initialize:

Let $\lambda_h^{(0)}$ and $\phi_h^{(0)}$ be arbitrarily chosen elements of M_h and W_h , respectively.

Advance ϕ_h :

Once
$$\lambda_h^{(n)}$$
 and $\phi_h^{(n)}$ are known, $\phi_h^{(n+1)}$ is defined by
$$\phi_h^{(n+1)} = \phi_h^{(n)} - \omega_n(A\phi_h^{(n)} + B^*\lambda_h^{(n)} - f_h).$$

Advance λ_h :

(17.20)
$$\lambda_{\mathbf{h}}^{(\mathbf{n}+1)} = \lambda_{\mathbf{h}}^{(\mathbf{n})} + \rho_{\mathbf{n}}(\mathbf{B}\phi_{\mathbf{h}}^{(\mathbf{n})} \mathbf{g}_{\mathbf{h}}).$$

Let $\lambda_h^{(n+1)}$ be defined by

Thus a variant of the augmented Lagrangian method is obtained by replacing A by A_r in (17.19).

Remark 17.4. If ω_n is taken to be equal to ρ_n , then the algorithm of Arrow-Hurwitz is also associated with the virtual time dependent problem

$$\begin{cases} a(\phi_h(t), \psi_h) + (\frac{d}{dt} \phi_h(t), \psi_h) + b(\psi_h, \lambda_h(t)) = (f, \psi_h) & \text{for all } \psi_h \in W_h, \\ b(\phi_h(t), \mu_h) + (\frac{d}{dt} \lambda_h(t), \mu_h) = (g, \mu_h) & \text{for all } \mu_h \in M_h, \end{cases}$$

of which the solution of (17.1) is the steady state solution.

Alternating-Direction Method. These methods can be used both in 2 dimensions and in 3 dimensions when the triangulation $T_{\rm h}$ of the domain Ω is a triangulation by rectangles. Here again one introduces a virtual time dependent problem to obtain the solution of the original problem as its steady state. The idea is, heuristically speaking, to reduce the solution of the two-(or three-) dimensional problem to the solution of one-dimensional problems by advancing half (or a third of) a time step in the x_1 -direction, then half (or a third of) a time step in the x_2 -direction (and in the three-dimensional case, a third of a time step in the x_3 -direction).

For example, in the two-dimensional case, for the Raviart-Thomas elements the subspace W_h of $H(\text{div};\Omega)$ in which the vectors are approximated is locally $P_{k+1,k} \times P_{k,k+1}$ and we suppose that the basis $\{w_j\}_{j=1}^{2m}$ has been chosen such that $W_{h,1} = \text{span} \ \{w_j\}_{j=1}^{m}$ is locally $P_{k+1,k} \times \{0\}$, and $W_{h,2} = \text{span} \ \{w_j\}_{j=m+1}^{2m}$ is locally $\{0\} \times P_{k,k+1}$. Thus the vector $w_h \in W_h$ may be written uniquely as $w_h = w_{h,1} + w_{h,2}$ where $w_{h,1} \in W_{h,1}$ and $w_{h,2} \in W_{h,2}$, and the matrices A and B are block diagonal:

$$A = \begin{pmatrix} A_1 & 0 \\ 0 & A_2 \end{pmatrix} , \qquad B = \begin{pmatrix} B_1 & 0 \\ 0 & B_2 \end{pmatrix} .$$

As in Remark 17.4 a fictitious time dependent problem is introduced:

$$\begin{cases} a(\phi_h(t), \psi_h) + b(\psi_h, \lambda_h(t)) = (f, \psi_h) & \text{for all } \psi_h \in W_h, \\ b(\phi_h(t), \mu_h) - (\frac{d}{dt} \lambda_h(t), \mu_h) = (g, \mu_h) & \text{for all } \mu_h \in M_h. \end{cases}$$

The two step iterative procedure is obtained by introducing an implicit time discretization as follows:

$$\begin{cases} a(\phi_{h,1}^{(n^*)} + \phi_{h,2}^{(n)}, \psi_h) + b(\psi_h, \lambda_h^{(n+\frac{1}{2})}) = (f, \psi_h) & \text{for all } \psi_h \in W_{h,1}, \\ b(\phi_{h,1}^{(n^*)} + \phi_{h,2}^{(n)}, \mu_h) - (\frac{\lambda_h^{(n+\frac{1}{2})} - \lambda_h^{(n)}}{\rho_n}, \mu_h) = (g, \mu_h) & \text{for all } \mu_h \in M_h, \end{cases}$$

$$\begin{cases} a(\phi_{h,1}^{(n^*)} + \phi_{h,2}^{(n+1)}, \psi_h) + b(\psi_h, \lambda_h^{(n+1)}) = (f, \psi_h) & \text{for all } \psi_h \in W_{h,2}, \\ b(\phi_{h,1}^{(n^*)} + \phi_{h,2}^{(n+1)}, \mu_h) - (\frac{\lambda_h^{(n+1)} - \lambda_h^{(n+\frac{1}{2})}}{\rho_n}, \mu_h) = (g, \mu_h) & \text{for all } \mu_h \in M_h, \end{cases}$$

which gives the algebraic systems

$$\begin{cases} A_1 \phi_{h,1}^{(n^*)} + B_1^* \lambda_h^{(n+\frac{1}{2})} &= f_{h,1} - A_2 \phi_{h,2}^{(n)} \\ \rho_n B_1 \phi_{h,1}^{(n^*)} - \lambda_h^{(n+\frac{1}{2})} &= \rho_n (g_h - B_2 \phi_{h,2}^{(n)}) - \lambda_h^{(n)}, \end{cases}$$

and

$$\begin{cases} A_2 \phi_h^{(n+1)} + B_2^* \lambda_h^{(n+1)} = f_{h,2} - A_1 \phi_{h,1}^{(n*)} \\ \rho_n B_2 \phi_h^{(n+1)} - \lambda_h^{(n+1)} = \rho_n (g_h - B_1 \phi_{h,1}^{(n*)}) - \lambda_h^{(n+\frac{1}{2})}, \end{cases}$$

From the first pair of the above equations $\lambda_h^{(n+\frac{1}{2})}$ can be eliminated to obtain

$$(\mathbf{A}_{1} + \rho_{n} \mathbf{B}_{1}^{*} \mathbf{B}_{1}) \ \phi_{\mathbf{h},1}^{(n^{*})} = \mathbf{f}_{\mathbf{h},1} - \rho_{n} \mathbf{B}_{1}^{*} \mathbf{g}_{\mathbf{h}} - (\mathbf{A}_{2} + \rho_{n} \mathbf{B}_{1}^{*} \mathbf{B}_{2}) \ \phi_{\mathbf{h},2}^{(n)} + \ \mathbf{B}_{1}^{*} \lambda_{\mathbf{h}}^{(n)},$$

and from the second pair, eliminating $\lambda_h^{(n+1)}$, one obtains

$$(\mathbf{A}_2 = \rho_n \mathbf{B}_2^* \mathbf{B}_2) \ \phi_{h,2}^{(n+1)} = \mathbf{f}_{h,2} - \rho_n \mathbf{B}_2^* \mathbf{g}_h - (\mathbf{A}_1 + \rho_n \mathbf{B}_2^* \mathbf{B}_1) \ \phi_{h,1}^{(n^*)} + \ \mathbf{B}_2^* \lambda_h^{(n+\frac{1}{2})}.$$

Thus the algorithm may be given as follows:

Initialize:

Let $\lambda_h^{(0)}$ be an arbitrarily chosen element of M_h .

The initialization $\phi_{h,2}^{(0)}$ of $\phi_{h,2}$ is obtained as the solution of

(17.21)
$$A_2 \phi_{h,2}^{(0)} = f_{h,2} - B_2^* \lambda_h^{(0)}$$

Advance a half time step in the x, direction:

Once $\lambda_h^{(n)}$ and $\phi_{h,2}^{(n)}$ are known,

$$(17.22) \qquad (A_1 + \rho_n B_1^* B_1) \ \phi_{h,1}^{(n^*)} = f_{h,1} - \rho_n B_1^* g_h - (A_2 + \rho_n B_1^* B_2) \ \phi_{h,2}^{(n)} + B_1^* \lambda_h^{(n)}.$$
 Let $\lambda_h^{(n+\frac{1}{2})}$ be defined by

(17.23)
$$\lambda_{h}^{(n+\frac{1}{2})} = \lambda_{h}^{(n)} + \rho_{n}(B_{1}\phi_{h,1}^{(n^{*})} + B_{2}\phi_{h,2}^{(n)} - g_{h}).$$

Advance a half time step in the x₂ direction:

Once
$$\lambda_h^{(n+\frac{1}{2})}$$
 and $\phi_{h,1}^{(n^*)}$ are known,

(17.24)
$$(A_2 + \rho_n B_2^* B_2^*) \phi_{h,2}^{(n+1)} = f_{h,2} - \rho_n B_2^* g_h - (A_1 + \rho_n B_2^* B_1^*) \phi_{h,1}^{(n^*)} + B_2^* \lambda_h^{(n+\frac{1}{2})}.$$
Let $\lambda_h^{(n+1)}$ be defined by

(17.25)
$$\lambda_{h}^{(n+1)} = \lambda_{h}^{(n+\frac{1}{2})} + \rho_{n}(B_{1}\phi_{h,1}^{(n^{*})} + B_{2}\phi_{h,2}^{(n+1)} - g_{h}).$$

Correct:

Finally to obtain
$$\phi_{h,1}^{(n+1)}$$
, solve

(17.26)
$$A_1 \phi_{h,1}^{(n+1)} = f_{h,1} - B_1^* \lambda_h^{(n+1)}.$$

This proceedure without the correction step was introduced by BROWN (1982), and as given here by DOUGLAS, DURAN & PIETRA (1987a), (1987b). It is an adaptation of the original alternating-direction procedure of PEACEMAN & RACHFORD (1955). Also in DOUGLAS, DURAN & PIETRA (1987a), (1987b) an alternating-direction scheme for the three-dimensional case is introduced. This procedure is based on the scheme of DOUGLAS (1962). As with the classical alternating directions iterative schemes, only limited convergence results have been obtained; cf. the above references.

For the Brezzi-Douglas-Marini elements there is no natural decomposition of W_h as for the Raviart-Thomas elements. Thus the above algorithm is not applicable. However, an alternating-direction procedure for these elements involving the introduction of two virtual time variables is given for 2 dimensions in DOUGLAS & PIETRA (1986) and for 3 dimensions in DOUGLAS, DURAN & PIETRA (1987a), (1987b) .

Mixed Hybrid Method . As pointed out earlier, the reason that the algebraic system derived from the dual mixed method poses special problems is that due to the degrees of freedom imposing the interelement continuity requirements on the elements ψ_h of W_h , the matrix A^{-1} is likely to be full. A technique often used to overcome this problem is to replace the dual mixed formulation by a mixed hybrid formulation,cf. FRAEIJS DE VEUBEKE (1965). The idea is to eliminate the interelement continuity requirements from the space W_h obtaining a space \hat{W}_h and to impose instead the desired continuity on the solution $\hat{\phi}_h \in \hat{W}_h$ via Lagrangian multipliers. Problem (17.1) is thus replaced by a problem of the following form: find $\hat{\phi}_h, \hat{\lambda}_h$, and η_h satisfying

$$\begin{pmatrix} \hat{\phi}_{h} \epsilon & \hat{W}_{h}, & \hat{\lambda}_{h} \epsilon & M_{h}, & \eta_{h} \epsilon & N_{h}, \\ & a(\hat{\phi}_{h}, \psi_{h}) + b(\psi_{h}, \hat{\lambda}_{h}) + c(\psi_{h}, \eta_{h}) & = (f, \psi_{h}) & \text{for all } \psi_{h} \epsilon & \hat{W}_{h}, \\ & b(\hat{\phi}_{h}, \mu_{h}) & = (g, \mu_{h}) & \text{for all } \mu_{h} \epsilon & M_{h}, \\ & c(\hat{\phi}_{h}, \tau_{h}) & = 0 & \text{for all } \tau_{h} \epsilon & N_{h}, \end{pmatrix}$$

where $N_h \subset \Pi$ $L^2(T')$ is the space of Lagrangian multipliers, and the bilinear $T' < T \in T_h$

form $c(\cdot, \cdot)$ on $\hat{W}_h \times N_h$ is such that the induced linear map C from W_h into N_h has as kernel the elements of \hat{W}_h which belong in fact to W_h . The linear system thus obtained is

(17.28)
$$\begin{pmatrix} \hat{A} & \hat{B}^* & C^* \\ \hat{B} & 0 & 0 \\ C & 0 & 0 \end{pmatrix} \begin{pmatrix} \hat{\phi}_h \\ \hat{\lambda}_h \\ \eta_h \end{pmatrix} = \begin{pmatrix} f_h \\ g_h \\ 0 \end{pmatrix},$$

where $\hat{\bf A}$ is now block diagonal as well as symmetric and positive definite, and $\hat{\phi}_{\bf h}$ and $\hat{\lambda}_{\bf h}$ are equal as functions to $\phi_{\bf h}$ and $\lambda_{\bf h}$, respectively. Eliminating $\hat{\phi}_{\bf h}$ from the system,

(17.29)
$$\hat{\phi}_{h} = \hat{A}^{-1} (f_{h} - \hat{B}^{*} \hat{\lambda}_{h} - C^{*} \eta_{h}),$$

we obtain

(17.30)
$$\begin{pmatrix} \hat{\mathbf{B}}\hat{\mathbf{A}}^{-1}\hat{\mathbf{B}}^{*} & \hat{\mathbf{B}}\hat{\mathbf{A}}^{-1}\mathbf{C}^{*} \\ \hat{\mathbf{C}}\hat{\mathbf{A}}^{-1}\hat{\mathbf{B}}^{*} & \hat{\mathbf{C}}\hat{\mathbf{A}}^{-1}\mathbf{C}^{*} \end{pmatrix} \begin{pmatrix} \hat{\lambda}_{h} \\ \eta_{h} \end{pmatrix} = \begin{pmatrix} \hat{\mathbf{B}}\hat{\mathbf{A}}^{-1}\mathbf{f}_{h}^{-g}\mathbf{h} \\ \hat{\mathbf{C}}\hat{\mathbf{A}}^{-1}\mathbf{f}_{h} \end{pmatrix} ,$$

and the space N_h and the form $c(\cdot, \cdot)$ have been constructed in such a way that this system is symmetric and positive definite. But now we further note that $\hat{B}\hat{A}^{-1}B^*$ is also block diagonal. Thus we have

(17.31)
$$\hat{\lambda}_{h} = (\hat{\mathbf{B}}\hat{\mathbf{A}}^{-1}\hat{\mathbf{B}}^{*})^{-1} \{\hat{\mathbf{B}}\hat{\mathbf{A}}^{-1}\mathbf{f}_{h}^{-1}\mathbf{g}_{h}^{-1}\hat{\mathbf{B}}\hat{\mathbf{A}}^{-1}\mathbf{C}^{*}\eta_{h}\}$$

and

This system is symmetric, positive definite, and sparse. Thus we can obtain η_h by solving (17.32), and then $\hat{\lambda}_h$ and $\hat{\phi}_h$ are given by (17.31) and (17.30), respectively. It is interesting to note that η_h often has a physical significance. For our model problem η_h represents the value of ϕ_h on the faces T' and in fact in certain cases can be used together with $\hat{\phi}_h$ to obtain a new approximation that converges more rapidly than $\hat{\phi}_h$ or equivalently ϕ_h ; cf. ARNOLD & BREZZI (1985) and BREZZI, DOUGLAS & MARINI (1985). Mixed hybrid methods will be taken up again in Section 21.

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CHAPTER V

HYBRID FINITE ELEMENT METHODS

18 - EXAMPLES OF PRIMAL HYBRID FINITE ELEMENT METHODS

We shall develop several examples of hybrid formulations for the model problem: find p and u satisfying the system of first order partial differential equations (4.4) - (4.5) and the Robin's boundary condition (15.1): $p \cdot \nu + tu = g$ on Γ , with t a continuous, strictly positive function on Γ . We shall indicate how to adapt the method to treat a problem with Dirichlet or Neumann boundary conditions.

To a triangulation T_h of the domain $\overline{\Omega}$, we have associated the spaces $H^1(T_h)$ and $H(\operatorname{div};T_h)$, cf. (5.11) and (6.16). The space $H^1(T_h)$ is a subspace of $L^2(\Omega)$ isomorphic to the product space Π $H^1(T)$, and the space $H(\operatorname{div};T_h)$ is a $T \in T_h$ subspace of $(L^2(\Omega))^n$ isomorphic to the product space Π $H(\operatorname{div};T)$. The hybrid $T \in T_h$

formulations, both primal and dual, will make use of a Lagrangian multiplier belonging to the space $L^2(\partial T_h)$ which is defined to be the product space

(18.1)
$$L^{2}(\partial T_{h}) = \prod_{T \in T_{h}} L^{2}(\partial T).$$

For the primal hybrid formulation, the subspace of multipliers will be

(18.2)
$$E(\partial T_h) = \left\{ \mu = (\mu_T)_{T \in T_h} \in L^2(\partial T_h) ; \exists q \in H(\operatorname{div}; \Omega) \right.$$
s.t. $\forall T \in T_h, \mu_T = q \cdot \nu_T \text{ on } \partial T \right\},$

where for this definition $\underline{q} \cdot \nu_T$ denotes the normal trace on ∂T of the restriction to T of the vector function \underline{q} . Thus $\underline{q} \cdot \nu_T$ is a priori an element of $H^{-\frac{1}{2}}(\partial T)$, cf. Remark 18.1 below.

Theorems 1.2 and 1.3 allow us to characterize $E(\partial T_h)$ as being the subspace

(18.3)
$$E(\partial T_h) = \left\{ \mu = (\mu_T)_{T \in T_h} \in L^2(\partial T_h) ; \forall T' = T_1 \cap T_2 \text{ with } \right.$$

$$T_1, T_2 \in T_h, \qquad \mu_{T_1} + \mu_{T_2} = 0 \text{ on } T' \left. \right\}.$$

For the dual hybrid formulation the subspace of multipliers will be

(18.4)
$$L(\partial T_h) = \left\{ \psi = (\psi_T)_{T \in T_h} \in L^2(\partial T_h) ; \exists v \in H^1(\Omega) \right.$$
s.t. $\forall T \in T_h, \quad \psi_T = v \text{ on } \partial T \right\},$

where in the expression " ψ_T = v on ∂T ", v denotes the trace on ∂T of the restriction to T of the function v. This trace is an element of $H^{\frac{1}{2}}(\partial T)$.

With Theorems 1.1 and 1.3 we may characterize $L(\partial T_h)$ as being the subspace

(18.5)
$$L(\partial T_h) = \left\{ \psi = (\psi_T)_{T \in T_h} \in \prod_{T \in T_h} H^{\frac{1}{2}}(\partial T), \forall T' = T_1 \cap T_2 \right\}$$
with $T_1, T_2 \in T_h, \psi_{T_1} = \psi_{T_2}$ on T'

<u>Remark 18.1</u>. In fact the natural space of Lagrangian multipliers for the primal hybrid formulation is

$$\begin{split} \mathbf{E}(\partial T_{\mathbf{h}}) &= \Big\{ \mu = (\mu_{\mathbf{T}})_{\mathbf{T} \in T_{\mathbf{h}}} \in \prod_{\mathbf{T} \in T_{\mathbf{h}}} \mathbf{H}^{-\frac{1}{2}}(\partial \mathbf{T}) \; ; \; \exists \; \mathbf{q} \in \mathbf{H} \; (\mathrm{div}; \Omega) \\ &\text{s.t.} \; \forall \; \mathbf{T} \in T_{\mathbf{h}} \; , \; \mu_{\mathbf{T}} = \; \mathbf{q} \cdot \nu_{\mathbf{T}} \; \text{in} \; \mathbf{H}^{-\frac{1}{2}}(\partial \mathbf{T}) \Big\}. \end{split}$$

The choice (18.2) while more restrictive is less cumbersome and leads to no restriction for examples of finite element methods.

Remark 18.2. The first examples of hybrid formulations given in Section 3 were described as having as Lagrangian multipliers functions defined on all Ω . Such a description is not in general well adapted for the numerical analysis of hybrid finite element methods. For hybrid formulations the Lagrangian

multipliers are defined only on ∂T_h . It is sometimes possible to extend the multipliers to functions defined on all Ω while taking into account complementary information when necessary. Such a procedure is a post processing technique.

A primal hybrid formulation of problem (4.4) - (4.5) with the Robin's boundary condition (15.1) is given as follows: find a pair (u,λ) satisfing

$$\begin{cases} (u,\lambda) \in H^1(T_h) \times E(\partial T_h), \\ a(u,v) + b(v,\lambda) = \int_{\Omega} fv \, dx & \text{for all } v \in H^1(T_h), \\ b(u,\mu) + d(\lambda,\mu) = -\sum_{T \in T_h} \int_{\partial T \cap \Gamma} \frac{1}{t} g\mu_T \, d\sigma & \text{for all } \mu \in E(\partial T_h), \end{cases}$$

where here

(18.7)
$$a(u,v) = \sum_{T \in T_h} \sum_{i,j=1}^n \int_{T} a_{ij} \frac{\partial u}{\partial x_j} \frac{\partial v}{\partial x_i} dx,$$

(18.8)
$$b(\mathbf{v},\mu) = -\sum_{\mathbf{T}\in T_{\mathbf{h}}} \int_{\partial \mathbf{T}} \mathbf{v}_{|\mathbf{T}} \mu_{\mathbf{T}} d\sigma,$$

(in this expression, $v_{|T}$ denotes the trace on ∂T of the restriction of v to T)

and

(18.9)
$$d(\lambda,\mu) = -\sum_{T \in T_h} \int_{\partial T \cap \Gamma} \frac{1}{t} \lambda_T \mu_T d\sigma.$$

It is easy to show that problem (18.6) has at most one solution. If the solution (p, u) of (15.2), which is also the solution of (15.8) and thus belongs to $H(\operatorname{div};\Omega)\times H^1(\Omega)$, satisfies the regularity hypothesis that $p\cdot\nu_T\in L^2(\partial T)$ for each $T\in T_h$, then the pair $(u,(p\cdot\nu_T)_{T\in T_h})$ is a solution of (18.6).

Thus we have

(18.10)
$$\lambda_{T} = p \cdot \nu_{T}$$
 on ∂T for each $T \in T_{h}$.

<u>Remark 18.3</u>. The spaces $H^1(T_h)$ and $L^2(\partial T_h)$ are supplied with Hilbert space structures in a natural way. With $W = H^1(T_h)$ and with $M = E(\partial T_h)$

having the norm induced by that of $L^2(\partial T_h)$, we note that the bilinear form $b(\cdot, \cdot)$ does not satisfy the inf-sup condition on W×M.

Toward the end of constructing an approximation (u_h, λ_h) of the solution (u, λ) of problem (18.6) we suppose that for each $T \in T_h$ we have two finite dimensional spaces $P_T \subset H^1(T)$ and $X_T \subset L^2(\partial T)$, and we put

(18.11)
$$W_h = \{ v \in H^1(T_h) ; \forall T \in T_h, v_{|_T} \in P_T \}$$

and

(18.12)
$$M_h = \{ \mu = (\mu_T)_{T \in T_h} \in E(\partial T_h) ; \forall T \in T_h, \mu_T \in X_T \}.$$

We then seek (u_h, λ_h) satisfying

$$\begin{cases} (\mathbf{u_h}, \lambda_h) \in \mathbf{W_h} \times \mathbf{M_h}, \\ \\ \mathbf{a}(\mathbf{u_h}, \mathbf{v_h}) + \mathbf{b}(\mathbf{v_h}, \lambda_h) = \int_{\Omega} \mathbf{f} \mathbf{v_h} \, \, \mathrm{d} \mathbf{x} & \text{for all } \mathbf{v_h} \in \mathbf{W_h}, \\ \\ \mathbf{b}(\mathbf{u_h}, \mu_h) + \mathbf{d}(\lambda_h, \mu_h) = -\sum_{T \in T_h} \int_{\partial T \cap T} \frac{1}{t} \, \mathbf{g} \mu_{hT} \mathrm{d} \sigma & \text{for all } \mu_h \in \mathbf{M_h}. \end{cases}$$

<u>Remark 18.4</u>. To treat the analogous problem with a Dirichlet boundary condition $u = \overline{u}$ on Γ , where \overline{u} is a given function in $H^{\frac{1}{2}}(\Gamma)$ we would keep the same choice of subspaces W_h and M_h and look for (u_h, λ_h) satisfying

$$\begin{cases} (u_h, \lambda_h) \in W_h \times M_h, \\ a(u_h, v_h) + b(v_h, \lambda_h) = \int_{\Omega} f v_h dx & \text{for all } v_h \in W_h, \\ b(u_h, \mu_h) = -\sum_{T \in T_h} \int_{\partial T \cap \Gamma} \overline{u} \mu_{hT} d\sigma & \text{for all } \mu_h \in M_h. \end{cases}$$

By contrast, to treat a homogeneous Neumann problem such as described in (14.1) or (14.6) we would introduce the subspace M_{0h} defined by

$$\mathbf{M}_{0h} = \left\{ \mu_{h} \in \mathbf{M}_{h} ; \forall T \in T_{h}, \mu_{hT} = 0 \text{ on } \partial T \cap \Gamma \right\}$$

and look for (u_h, λ_h) satisfying

$$\begin{cases} (u_h, \lambda_h) \in W_h \times M_{0h}, \\ a(u_h, v_h) + b(v_h, \lambda_h) &= \int_{\Omega} f v_h dx. & \text{for all } v_h \in W_h, \\ b(u_h, \mu_h) &= 0 & \text{for all } \mu_h \in M_{0h}. \end{cases}$$

The bilinear form $a(\cdot, \cdot)$ is not $H^1(T_h)$ -elliptic and will not be W_h -elliptic. However, we shall introduce a hypothesis to make it V_h -elliptic, where

(18.14)
$$V_h = \{ v_h \in W_h ; \forall \mu_h \in M_h, b(v_h, \mu_h) = 0 \}.$$

<u>THEOREM 18.1</u>. We suppose that the two subspaces W_h and M_h are compatible in the sense that the discrete inf-sup condition is satisfied:

(18.15)
$$\left\{ \mu_{h} \in M_{h} ; \forall v_{h} \in W_{h}, b(v_{h}, \mu_{h}) = 0 \right\} = \left\{ 0 \right\}.$$

We suppose further that for each interface $T' = T_1 \cap T_2$ with T_1 , $T_2 \in T_h$, there exists a function $\chi_T \in L^2(T')$ with $\int_{T'} \chi_{T'} d\sigma \neq 0$ such that the function $\Lambda = (\Lambda_T)_{T \in T_h}$ in $E(\partial T_h)$ defined by $\Lambda_{T_1} = \chi_{T'}$ on T', $\Lambda_{T_2} = -\chi_{T'}$ on T', and $\Lambda_{T} = 0$ on all other faces for all $T \in T_h$, belongs to the subspace M_h . Similarly we suppose that for each boundary face $T \subset \Gamma$ of an element $T_0 \in T_h$, there exists a function $\chi_{T'} \in L^2(T')$ with $\int_{T'} \chi_{T'} d\sigma \neq 0$ such that the function $\Lambda = (\Lambda_T)_{T \in T_h}$ in $E(\partial T_h)$ defined by $\Lambda_{T_0} = \chi_{T'}$ on T' and $\Lambda_{T} = 0$ on all other faces for all $T \in T_h$, belongs to the subspace M_h . Then problem (18.13) has a unique solution.

<u>Proof.</u> Each function $v_h \in V_h$ satisfying $a(v_h, v_h) = 0$ is a function whose restriction to T, for each $T \in T_h$, is constant. The first part of the additional hypothesis guarantees that such a function may be identified with a function constant on each component of Ω . The second part forces such a function to vanish on $\Gamma = \partial \Omega$. We conclude that the form $a(\cdot, \cdot)$ is V_h -elliptic and, in light of Remark 10.1, that the theorem follows.

It is important to note that the compatibility condition (18.15) that should be satisfied by the subspaces W_h and M_h given in the forms (18.11) and (18.12) is satisfied as soon as, for each $T \in \mathcal{T}_h$, the spaces P_T and X_T satisfy the local compatibility condition

(18.16)
$$\left\{\mu \in X_T ; \forall v \in P_T, \int_{\partial T} \mu v d\sigma = 0\right\} = \left\{0\right\}.$$

By contrast, the additional hypotheses of Theorem 18.1 demand that the functions belonging to V_h have a minimum amount of continuity at the interfaces of the triangulation. In the example that follows this continuity hypothesis will be trivially satisfied. (We will in fact choose χ_T , to be a constant function.)

THEOREM 18.2. Let Ω be an open polygonal domain in \mathbb{R}^2 triangulated by triangles T. Let k be a positive integer, and put

$$k = \begin{cases} k & \text{if } k \text{ is odd} \\ k + 1 & \text{if } k \text{ is even.} \end{cases}$$

For each triangle $T \in T_h$, define the subspace $X_T \subset L^2(\partial T)$ by

$$(18.17) X_{T} = D_{k}(\partial T)$$

(With the definition of the space $E_{\partial T_h}^{(k)}$ given by (6.30), we have $M_h = E_{\partial T_h}^{(k)}$). Suppose that the subspace $P_T \subset H^1(T)$ is chosen such that the following property is satisfied:

(18.18)
$$\forall w \in P_{k^*}(\partial T)$$
, $\exists v \in P_{T}$ s.t. $v = w$ on ∂T .

Then problem (18.13) has a unique solution.

<u>Proof.</u> The elements of the space $D_k(\partial T)$ defined in (6.32) are polynomials of degree $\leq k-1$ on each edge. It is clear that condition (18.16) is satisfied if each element of $P_{k+1}(\partial T)$ vanishing at the vertices of T is the trace of a function $v \in P_T$. Even though we have dim $P_k(\partial T) = \dim D_k(\partial T)$, we know (cf. RAVIART & THOMAS (1977b, lemma 4)) that the subspace of $D_k(\partial T)$ consisting of all those functions $\mu \in D_k(\partial T)$ satisfying $\int_{\partial T} \mu \, w \, d\sigma = 0$ for each $w \in P_k(\partial T)$ is trivial only if k is odd. It is of dimension one when k is even.

 $\overline{}$

Before giving an error estimate we need to make precise our choice of norms. For the sake of simplicity we suppose here that the triangulation is uniformly regular. To each function $v \in H^1(T_h)$ we associate the function $\delta v \in L^2(\partial T_h)$, constant on each edge $T' \in \partial T_h$, defined by

(18.19)
$$\widehat{\delta \mathbf{v}} = \begin{cases} \frac{1}{\text{mes } \mathbf{T}'} \int_{\mathbf{T}'} (\mathbf{v}|_{\mathbf{T}_1} - \mathbf{v}|_{\mathbf{T}_2}) \, d\sigma , & \text{if } \mathbf{T}' = \mathbf{T}_1 \cap \mathbf{T}_2; \, \mathbf{T}_1, \, \mathbf{T}_2 \in T_h \\ & \text{(the choice of } \mathbf{T}_1 \text{ being arbitrary)} \end{cases}$$

$$\frac{1}{\text{mes } \mathbf{T}' \mathbf{T}'} \int_{\mathbf{T}'} \mathbf{v}|_{\mathbf{T}} \, d\mathbf{T}, \qquad \text{if } \mathbf{T}' = \mathbf{T} \cap \mathbf{\Gamma}.$$

Once the hypotheses of Theorem 18.1 are satisfied, we clearly have $\delta v = 0$ if $v \in H^1_0(\Omega) \oplus V_h$. The space $L^2(\partial T_h)$ is given the norm

(18.20)
$$\|\mu\|_{0,\partial T_{\mathbf{h}}} = \left\{ \sum_{\mathbf{T} \in T_{\mathbf{h}}} \|\mu_{\mathbf{T}}\|_{0,\partial \mathbf{T}}^{2} \right\}^{\frac{1}{2}},$$

the space $H^1(T_h)$ the norm

(18.21)
$$[v]_{1,T_{h}} = \left\{ \sum_{T \in T_{h}} |v|_{1,T}^{2} + |h^{-1}| |\widehat{\delta v}|_{0,\partial T_{h}}^{2} \right\}^{\frac{1}{2}},$$

and the space $E(\partial T_h)$ the norm $|||\cdot|||_{\frac{1}{2},\partial T_h}$ defined in (6.36). One can show (by adapting, for example, the demonstration of Theorem V.4.3 in THOMAS (1977)) that the Poincaré-Friedrichs inequality,

(18.22)
$$\|\mathbf{v}\|_{0,\Omega} \le C \|\mathbf{v}\|_{1,T_h}$$
 for all $\mathbf{v} \in H^1(T_h)$,

is satisfied with a constant C independent of h. (Note that the two norms on $H^1(T_h)$, $v \longrightarrow \begin{bmatrix} v \end{bmatrix}_{1,T_h}$ and $v \longrightarrow ||| v |||_{1,T_h} = (\sum_{T \in T_h} (|v||_{1,T}^2 + h^{-2}||v||_{0,T}^2)^{\frac{1}{2}}$, are equivalent uniformly in h.) One may also verify that there exist constants c and C_k independent of h such that

(18.23)
$$|||\mu||_{-\frac{1}{2}}, \partial T_{h} \leq c h^{\frac{1}{2}} ||\mu||_{0,\partial T_{h}} \quad \text{for all } \mu \in L^{2}(\partial T_{h}),$$
 and

(18.24)
$$\|\mu_h\|_{0,\partial T_h} \le C_k^{-\frac{1}{2}} \|\|\mu\|\|_{-\frac{1}{2},\partial T_h} \text{ for all } \mu_h \in E_{\partial T_h}^{(k)}.$$

THEOREM 18.3. We suppose the hypotheses of Theorem 18.2 except here we fix for each $T \in T_h$

(18.25)
$$X_T = D_k(\partial T)$$
, $P_T = P_k^*$ (T).

Then there exists a constant C independent of h such that we have

$$[u-u_h]_{1,T_h} + ||| \lambda-\lambda_h |||_{-\frac{1}{2},\partial T_h} \le Ch^k \left\{ |u|_{k+1,\Omega} + |\underline{p}|_{k,\Omega} \right\}$$

once the solution (p, u) of (15.2) and of (15.8) belongs to $(H^k(\Omega))^2 \times H^{k+1}(\Omega)$; (λ is expressed as a function of p in (18.10).)

Proof. It is clear that the bilinear form a(•,•) satisfies

$$a(u,v) \le A \left[u\right]_{1,T_h} \left[v\right]_{1,T_h}$$
 for all $u, v \in H^1(T_h)$,

and

$$a(v_h, v_h) \ge \alpha_p \left[v_h\right]_{1, T_h}^2$$
 for all $v_h \in V_h$,

as $v_h \in V_h$ implies $\delta v_h = 0$. We can also establish the existence of a constant β , $\beta > 0$, such that

$$\sup_{\mathbf{v_h} \in \mathbf{W_h}} \frac{\mathbf{b}(\mathbf{v_h}, \mu)}{\left[\mathbf{v_h}\right]_{1,\Omega}} \geq \beta ||| \mu_h |||_{-\frac{1}{2}, \partial T_h} \quad \text{for all } \mu_h \in \mathbf{M_h}.$$

Finally we show the existence of a constant Bo such that

$$b(v,\mu) \leq B_0 \quad [v]_{1,T_h} \quad ||| \; \mu \; |||_{-\frac{1}{2}}, \partial T_h$$
 for all $v \in H^1(T_h)$ and all $\mu \in E_0(\partial T_h)$,

where $E_0(\partial T_h)$ is the subspace of functions $\mu \in E(\partial T_h)$ such that $\int_T \mu \, d\sigma = 0$ on each edge $T' \subset \Gamma$. Now, for a Dirichlet boundary condition or for a Neumann boundary condition, using a variant of Theorem 10.4, we can establish the existence of a constant C, independent of h, such that

$$[u-u_h]_{1,T_h} + |||\lambda-\lambda_h|||_{-\frac{1}{2},\partial T_h} \leq C \left\{ \inf_{v_h \in W_h} [u-v_h]_{1,T_h} + \inf_{\{\mu_h \in M_h; \lambda-\mu_h \in E_0(\partial T_h)\}} |||\lambda-\mu_h|||_{-\frac{1}{2},\partial T_h} \right\},$$

and then using (6.37) deduce the inequality (18.26).

To handle the case of Robin's boundary condition we introduce the direct sum decomposition of $M = E(\partial T_h)$ in the form $M = M_0 \oplus \overline{M}$, where M_0 is the subspace of all $\mu \in M$ such that for each $T \in T_h$ with $T \cap \Gamma \neq \phi$, $\mu_T = 0$ on each edge T'contained in the boundary Γ . Problem (18.6), respectively (18.13), may be put in the abstract form (10.1), respectively (10.6), posed in $(H^1(T_h) \times \overline{M}) \times M_0$, with for bilinear forms

$$(u,\bar{\lambda}),(v,\bar{\mu}) \longrightarrow a(u,v) + b(v,\bar{\lambda}) - b(u,\bar{\mu}) - d(\bar{\lambda},\bar{\mu})$$

and

$$(\mathbf{v}, \bar{\boldsymbol{\mu}}), \boldsymbol{\mu_0} \longrightarrow b(\mathbf{v}, \boldsymbol{\mu_0})$$
.

We then deduce as above

$$\begin{split} \left[\mathbf{u} - \mathbf{u}_{\mathbf{h}} \right]_{1,T_{\mathbf{h}}} &+ || & \tilde{\lambda} - \tilde{\lambda}_{\mathbf{h}} ||_{0,\partial T_{\mathbf{h}}} + |||\lambda_{0} - \lambda_{0\mathbf{h}}|||_{-\frac{1}{2},\partial T_{\mathbf{h}}} \\ &\leq C \|\mathbf{h}^{\mathbf{k}} \| \left\{ \|\mathbf{u}\|_{\mathbf{k}+1,\Omega} + \|\mathbf{p}\|_{\mathbf{k},\Omega} \right\}. \end{split}$$

<u>Remark 18.5</u>. If the integer k is even, then k = k+1, and the choice $P_T = P_{k+1}(T)$ is not optimal. We can obtain the same results by taking P_T to be the space generated by $P_k(T)$ together with an appropriately choosen element of $P_{k+1}(T)$. For more details see RAVIART & THOMAS (1977b).

The primal hybrid finite element method corresponding to the subspaces M_h and W_h associated with the spaces X_T and P_T of (18.25) by (18.11) and (18.12) is convergent for each positive integer k, in the following sense: we can show, in the absence of any hypothesis concerning the regularity of the solution u, that

$$\lim_{h \to 0} || u - u_h ||_{0,\Omega} = 0.$$

Approximation methods based on a principle analogous to that which was the basis for the development of primal hybrid finite element methods have been studied by BABUSKA (1973), BRAMBLE (1981), FALK (1976), and PITKÄRANTA (1979), (1980a), (1981). For a problem with a nonhomogeneous Dirichlet boundary condition $u = \bar{u}$ on Γ , we dualize the boundary condition using a multiplier $\lambda \in H^{-\frac{1}{2}}(\Gamma)$. The discrete inf-sup condition is established by means of compactness arguments; thus the results, notably existence and uniqueness, are assured only for "h sufficiently small".

19 PRIMAL HYBRID METHODS AND NONCONFORMING METHODS

We consider the model Dirichtlet problem

$$\left\{ \begin{array}{lll} -\ \Delta u \ = \ f & & \text{in } \ \Omega \ , \\ \\ u \ = \ \bar{u} & & \text{on } \ \Gamma \ , \end{array} \right.$$

where f is a given function in $L^2(\Omega)$, \bar{u} a given function in $H^{\frac{1}{2}}(\Gamma)$, and Ω a bounded, open domain in \mathbb{R}^2 with polygonal boundary Γ . For a triangulation $T_{
m h}$ of Ω by triangles, the simplest primal hybrid finite element method for solving this problem is given as follows: find (u_h, λ_h) satisfying

$$\begin{cases} (u_h, \lambda_h) \in W_h \times M_h , \\ a(u_h, v_h) + b(v_h, \lambda_h) = \int\limits_{\Omega} f v_h dx & \text{for all } v_h \in W_h , \\ b(u_h, \mu_h) = -\sum\limits_{T \in T_h} \int\limits_{\partial T \cap \Gamma} \bar{u} \mu_{hT} d\sigma & \text{for all } \mu_h \in M_h , \end{cases}$$
 with

with

$$\begin{aligned} \mathbf{W}_{\mathbf{h}} &= \left\{ \begin{array}{l} \mathbf{v} \in \ \mathbf{H}^{1}(T_{\mathbf{h}}) \ ; \ \forall \ \mathbf{T} \in T_{\mathbf{h}} \ , \ \mathbf{v}_{\mid \mathbf{T}} \in \mathbf{P}_{1}(\mathbf{T}) \right\} \ , \\ \\ \mathbf{M}_{\mathbf{h}} &= \left\{ \begin{array}{l} \mu = (\mu_{\mathbf{T}})_{\mathbf{T} \in T_{\mathbf{h}}} \in E(\partial T_{\mathbf{h}}) \ ; \ \forall \ \mathbf{T} \in T_{\mathbf{h}} \ , \ \mu_{\mathbf{T}} \in \mathbf{D}_{1}(\partial \mathbf{T}) \end{array} \right\}, \\ \\ \mathbf{a}(\mathbf{u}, \mathbf{v}) &= \sum_{\mathbf{T} \in T_{\mathbf{h}}} \int_{\mathbf{T}} \operatorname{grad} \ \mathbf{u} \cdot \operatorname{grad} \ \mathbf{v} \ d\mathbf{x} \ , \\ \\ \mathbf{b}(\mathbf{v}, \mu) &= -\sum_{\mathbf{T} \in T_{\mathbf{h}}} \int_{\partial \mathbf{T} \cap \mathbf{T}} \mu_{\mathbf{T}} \ \mathbf{v}_{\mid \mathbf{T}} \ d\sigma. \end{aligned}$$

Eliminating the Lagrangian multiplier, we obtain a characterization of uh as the solution of the variational problem

(19.2)
$$\begin{cases} u_h \in V_h^{\bar{u}}, \\ a(u_h, v_h) = \int_{\Omega} f v_h dx \end{cases}$$
 for all $v_h \in V_h$,

with

$$\mathbf{V}_{\mathbf{h}}^{\tilde{\mathbf{u}}} = \left\{ \mathbf{v}_{\mathbf{h}} \in \mathbb{W}_{\mathbf{h}} \; ; \; \forall \; \boldsymbol{\mu}_{\mathbf{h}} \in \mathbb{M}_{\mathbf{h}} \quad b(\mathbf{v}_{\mathbf{h}}, \boldsymbol{\mu}_{\mathbf{h}}) = -\sum_{\mathbf{T} \in \boldsymbol{T}_{\mathbf{h}}} \int_{\partial \mathbf{T} \cap \; \boldsymbol{\Gamma}} \bar{\mathbf{u}} \; \boldsymbol{\mu}_{\mathbf{h} \mathbf{T}} \; \mathrm{d}\boldsymbol{\sigma} \right\},$$

and

$$V_h = \left\{ v_h \in W_h ; \forall \mu \in M_h, b(v_h, \mu_h) = 0 \right\}.$$

It is a simple matter to check that V_h^u is the set of functions in W_h that are affine on each triangle $T \in T_h$, are continuous at the midpoint of each edge T' common to two triangles T_1 and T_2 of T_h and have for value at the midpoint of each boundary edge $T' \subset \Gamma$ the average value of \bar{u} on T'. The set of midpoints of the three sides of a triangle is P_1 -unisolvent. Thus a function $v_h \in V_h$ is uniquely determined by its values at the midpoints of the edges common to two triangles of T_h . It is not difficult to give a basis of V_h associated with this choice of degrees of freedom for a function in V_h . Thus, problem (19.2) can be solved directly. Such a procedure is said to be a nonconforming finite element method since V_h is not a subspace of $V = H_0^1(\Omega)$, the space of test functions for the primal variational formulation.

If the solution u belongs to $H^2(\Omega)$ we can obtain, using the results of the previous section, the error bound

(19.3)
$$\left\{ \sum_{T \in T_h} |u - u_h|^2_{1,T} \right\}^{\frac{1}{2}} + ||u - u_h||_{0,T} \le C \ h \ |u|_{2,\Omega} ,$$

an estimate that is classically obtained in a direct manner, cf. STRANG (1972) or STRANG & FIX (1973, Chapter 4.2).

After having calculated the solution of (19.2), to obtain the solution of (19.1) we have only to determine $\lambda_h \in M_h$ satisfying

(19.4)
$$b(v_h, \lambda_h) = \int_{\Omega} fv_h dx - a(u_h, v_h) \qquad \text{for each } v_h \in W_h,$$

i.e. the functions λ_{hT} , constant on each edge of T, are determined by

$$\int_{\partial T} v \ \lambda_{hT} \ d\sigma = \int_{T} fv \ dx - \int_{T} grad \ u_{h} \cdot grad \ v \ dx \ for each \ v \in P_{1}(T).$$

Thus once the nonconforming approximation u_h of u has been calculated, a postprocessing procedure furnishes directly an approximation of the flux of grad u across each edge T. If T is the edge common to T_1 and $T_2 \in T_h$, then the values λ_{hT_1} and λ_{hT_2} obtained from (19.4) satisfy

$$\lambda_{hT_1} + \lambda_{hT_2} = 0$$
, on $T' = T_1 \cap T_2$.

The flux across T' from T_1 into T_2 is minus that from T_2 into T_1 .

More elaborate examples of nonconforming finite elements derived from primal hybrid finite element methods for second order elliptic problems on an open set Ω in \mathbb{R}^2 may be found in RAVIART & THOMAS (1977b). For those examples the continuity of the functions of V_h is imposed at the Gauss points of the edges T. The example treated above, where the functions v_h are locally affine, is readily generalized to the 3-dimensional case, where Ω is an open set of \mathbb{R}^3 . However the situation quickly becomes more difficult for higher orders of approximation. An example with k=2 is given in FORTIN (1985).

In fact nonconforming finite element methods were first developed for elliptic problems more complicated than the model problem above, e.g.:

- 1) a fourth order elliptic problem where the space V is a subspace of $(H^2(\Omega))^2$ and is thus made up of functions of class C^1 on $\bar{\Omega}$. Examples and analyses of such methods may be found in CIARLET (1974), (1978, pp. 362-380), CIAVALDINI & NEDELEC (1974), and ZHONG-CI SHI (1984a).
- 2) Stokes problem formulated as a variational problem on a subspace of the space V of divergence free functions in $(H^1(\Omega))^n$. Analyses of such methods are given in CROUZEIX & RAVIART (1973), FORTIN (1981), HECHT (1981) and TEMAM (1977, pp. 172-181).
- 3) a problem of linear elasticity formulated as a variational problem on a subspace of $V = (H^1(\Omega))^n$. The motivation in this case is to obtain models less rigid than those provided by conforming methods, cf. in particular PIAN (1971) (1972), FRAEIJS DE VEUBEKE (1974b). The most popular noncorforming finite element for this type of problem is Wilson's brick (WILSON et al. (1973)). An analysis of the method may be found in LESAINT (1976), LESAINT & ZLAMAL (1980), ZHONH-CI SHI (1984b) or as a particular case of primal hybrid methods in THOMAS (1977,Ch. VI). Other interpretations using hybrid methods have been proposed in IRONS (1972), PIAN & TONG (1986).

We conclude this section by pointing out that the problem of finding necessary and sufficient conditions for convergence of a nonconforming finite element method has been the subject of many discussions between theoreticians and engineers, cf. in particular the response of IRONS & LOIKKANEN (1983) to STUMMEL (1979), (1980a), and (1980b). When a nonconforming method can be

interpreted as a primal hybrid method, possibly with reduced numerical integration, the study of its convergence is based in particular on the verification of the discrete inf-sup condition.

20 EXAMPLES OF DUAL HYBRID FINITE ELEMENT METHODS

Let T_h be a triangulation of the domain $\overline{\Omega}$. Associated to this triangulation are the spaces H (div; T_h) defined by (6.16) and $L(\partial T_h)$ introduced in (18.4) and (18.5). Given a function $f \in L^2(\Omega)$ denote by $Q^f(T_h)$ the affine manifold defined by

(20.1)
$$Q^{f}(T_{h}) = \{q \in H(div; T_{h}) ; \forall T \in T_{h}, div (q_{|T}) + f_{|T} = 0 \},$$

or more precisely defined as the space of vector functions $\mathbf{q} \in H(\operatorname{div};T_h)$ which satisfy $\int\limits_{\Omega} \mathbf{q} \cdot \mathbf{grad} \ \mathbf{v} \ d\mathbf{x} = \int\limits_{\Omega} \operatorname{fvdx}$, for each function $\mathbf{v} \in H^1(\Omega)$ whose trace on ∂T_h is trivial. In particular, $Q^0(T_h)$ is the vector space of functions $\mathbf{q} \in H(\operatorname{div};T_h)$ for which $\mathbf{q}_{|T}$ is divergence free for each $\mathbf{T} \in T_h$.

A dual mixed hybrid formulation of our model problem (4.4), (4.5) with Robin's boundary condition (15.1) is given as follows: find a pair (p, ϕ) satisfying

with here

(20.3)
$$a(p,q) = \sum_{i,j=1}^{n} \int_{T} A_{ij} p_{j} q_{i} dx,$$

(20.4)
$$b(\mathbf{q}, \psi) = -\sum_{\mathbf{T} \in T_{\mathbf{h}}} \int_{\partial \mathbf{T}} \psi_{\mathbf{T}} \mathbf{q} \cdot \mathbf{v}_{\mathbf{T}} d\sigma,$$

(in this last expression, $q \cdot \nu_T$ denotes the normal trace on ∂T of the restriction of q to T), and

(20.5)
$$d(\phi,\psi) = -\sum_{\mathbf{T}\in T_{\mathbf{h}}} \int_{\partial \mathbf{T}\cap \mathbf{\Gamma}} \mathbf{t} \ \phi_{\mathbf{T}} \psi_{\mathbf{T}} \ d\sigma.$$

After noting that, for each $T \in T_h$, the normal trace on the boundary of T is a surjective mapping of $\{q \in H(\operatorname{div};T) : \operatorname{div} q = 0\}$ onto $\{\mu \in L^2(\partial T) : \int_{\partial T} \mu \, d\sigma = 0\}$, it is easy to verify that problem (20.2) has at most one solution. One shows that if the solution (p,u) of (15.2) and of (15.8) belongs to $(H(\operatorname{div};\Omega) \cap H(\operatorname{div};T_h)) \times H^1(\Omega)$, (this is the same regularity hypothesis as that required for the existence of a solution to the primal hybrid problem (18.6)), then the pair $(p,u) \cap T \in T_h$ is a solution of (20.2). One thus has

(20.6)
$$\phi_{\mathbf{T}} = \mathbf{u} \quad \text{on } \partial \mathbf{T}$$
 for all $\mathbf{T} \in T_{\mathbf{h}}$.

<u>Remark 20.1.</u> If we equip the spaces $H(\text{div};T_h)$ and $L(\partial T_h)$ with Hilbert space structures, then $b(\cdot,\cdot)$ considered as a bilinear form on W×M, with $W = Q^0(T_h)$ a subspace of $H(\text{div};T_h)$ and $M = L(\partial T_h)$, does not satisfy the inf-sup condition. (This situation is analogous to the case of the primal hybrid formulation, cf. Remark 18.3.)

To construct an approximation of the solution (p,ϕ) of problem (20.2), we take for each $T\in T_h$ two finite dimensional subspaces, $Q_T\subset H(\text{div};T)$ and $Y_T\subset H^{\frac{1}{2}}(\partial T)$. To a given function $f\in L^2(\Omega)$, we associate the function $f_h\in L^2(\Omega)$ whose restriction to T, for each $T\in T_h$, is the orthogonal projection in $L^2(T)$ of $f_{\mid T}$ into the finite dimensional space $\{\text{div } q: q\in Q_T\}$. We then put

$$(20.7) W_h^f = \left\{ \underline{q} \in H(\operatorname{div}; T_h) ; \forall T \in T_h, \underline{q}_{|T} \in Q_T \right\} \cap Q^{f_h}(T_h),$$

$$(20.8) W_h = W_h^0 = \left\{ \ \ \underline{q} \ \in H(\text{div}; T_h) \ ; \ \forall \ T \in T_h, \ \ \underline{q}_{|T} \in Q_T \right\} \ \cap Q^{\circ}(T_h) \ ,$$

and

$$(20.9) M_h = \left\{ \psi = (\psi_T)_{T \in T_h} \in L(\partial T_h) ; \forall T \in T_h, \psi_T \in Y_T \right\}.$$

We seek a pair (p_h, ϕ_h) satisfying

$$\left(\begin{array}{c} (\underline{p}_h,\phi_h) \in W_h^f \times M_h \ , \\ \\ a(\underline{p}_h,\underline{q}_h) \ + \ b(\underline{q}_h,\phi_h) \ = \ 0 \\ \\ b(\underline{p}_h,\psi_h) \ + \ d(\phi_h,\psi_h) \ = \ -\sum_{T \in T_h} \int_{\partial T \cap \ \Gamma} \ g \ \psi_{hT} \ d\sigma \quad \text{ for all } \psi_h \in M_h. \end{array} \right)$$

<u>Remark 20.2</u>. The analogous problem with a Neumann boundary condition $p \cdot \nu = g$, where g is a given function in $L^2(\Gamma)$ satisfying $\int_{\Gamma} f \, dx + \int_{\Gamma} g \, d\sigma = 0$, is formulated as follows: find a pair (p_h, ϕ_h) satisfying

$$\begin{cases} (\underline{p}_h, \phi_h) \in W_h^f \times M_h \ , \\ \\ a(\underline{p}_h, \underline{q}_h) + b(\underline{q}_h \ , \ \phi_h) = 0 & \text{for all } \underline{q}_h \in W_h \ , \\ \\ b(\underline{p}_h, \psi_h) = -\sum_{T \in T_h} \int_{\partial T \cap \Gamma} g \, \psi_{hT} \, d\sigma & \text{for all } \psi_h \in M_h \ , \end{cases}$$

where W_h^f , W_h , and M_h are as defined in (20.7), (20.8), and (20.9).

To treat the homogeneous Dirichlet problem we seek a solution (p_h, ϕ_h) of the following problem :

$$\begin{cases} (\underline{p}_h, \phi_h) \in W_h^f \times M_{0h} ,\\ \\ a(\underline{p}_h, \underline{q}_h) + b(\underline{q}_h, \phi_h) = 0 \\ \\ b(\underline{p}_h, \psi_h) = 0 \end{cases} \qquad \text{for all } \underline{q}_h \in W_h ,\\ \\ b(\underline{p}_h, \psi_h) = 0 \qquad \qquad \text{for all } \psi_h \in M_{0h} ,$$

where M_{0h} is the subspace ot M_h defined by

(20.11)
$$M_{0h} = \{ \psi_h \in M_h : \forall T \in T_h, \psi_{hT} = 0 \text{ on } \partial T \cap \Gamma \}.$$

By construction of the function f_h , the affine manifold W_h^f is not empty. Applying Theorem 10.2 with Remark 10.7 we obtain with no difficulty the following theorem:

<u>THEOREM 20.1</u> Problem (20.10) has a unique solution if and only if the subspaces W_h and M_h are compatible in the sense that the discrete inf-sup condition is satisfied:

(20.12)
$$\left\{ \psi_{h} \in M_{0h} ; \forall q_{h} \in W_{h}, b(q_{h}, \psi_{h}) = 0 \right\} = \left\{ 0 \right\},$$

where M_{0h} is the subspace of M_h defined in (20.11).

The global condition (20.12) is satisfied once, for each T \in T_h , the choice of spaces Q_T and Y_T satisfy the local compatibility condition

$$\left\{\psi\in Y_{T} \;\;;\;\;\forall\;\; q\in Q_{T} \;\;\text{with div}\;\; q=0\;\;,\;\; \int_{\partial T} \psi\; q\cdot \nu_{T} \;\;\mathrm{d}\sigma=0\right\}=\; P_{0}(\partial T).$$

Next we develop an example where the above theory is applied.

THEOREM 20.2 . Let Ω be an open polygonal domain in \mathbb{R}^2 triangulated by triangles T. Let k be a positive integer, and put

$$k^* = \begin{cases} k & if k \text{ is odd }; \\ k+1 & if k \text{ is even.} \end{cases}$$

For each triangle $T \in T_h$ define the subspace $Y_T \subset H^{\frac{1}{2}}(\partial T)$ by

$$(20.14) Y_{T} = P_{k}(\partial T).$$

(With the definition of the space $L_{\partial T_h}^{(k)}$ given by (5.14), we have $M_h = L_{\partial T_h}^{(k)}$). Suppose that the space $Q_T \subset H(\text{div},T)$ is such that the following property is satisfied:

$$(20.15) \qquad \forall \mu \in D_{k}^{*}(\partial T) \ , \ \exists \ \underline{q} \in Q_{t} \ s.t. \ \underline{q} \cdot \underline{\nu}_{T} = \mu \ on \ \partial T.$$

Then problem (20.10) has a unique solution.

To establish this theorem it clearly suffices to show that the local compatibility condition (20.13) is satisfied. A demonstration can be found in THOMAS (1976, lemma 4.1). We point out that the subspace of all $\psi \in P_k(\partial T)$ which satisfy $\int_{\partial T} \psi \, \mu \, d\sigma = 0$ for each $\mu \in D_k(\partial T)$ with $\int_{\partial T} \mu \, d\sigma = 0$ reduces to $P_0(\partial T)$ only when k is odd. Otherwise, when k is even, it is a subspace of dimension 2.

We shall give error estimates using the norm $||\cdot||_{0,\Omega}$ on $(L_2(\Omega))^2$ for vector functions (with this method, on each $T \in T_h$, div $(p-p_h) = -(f-f_h)$ and is considered to be known or estimated a priori) and the norm $|||\cdot|||_{\frac{1}{2},\partial T_h}$ given in (5.19) for the traces on ∂T_h of functions in $H^1(\Omega)$. We suppose we have a uniformly regular family of triangulations, where the hypothesis of uniformness is made simply for ease of exposition. We note the existence of constants c and C_k independent of h such that

$$(20.16) ||\psi||_{0,\partial T_h} \leq ch^{\frac{1}{2}} |||\psi|||_{\frac{1}{2},\partial T_h} for all \psi \in H^{\frac{1}{2}}(\partial T_h) ,$$

and

$$(20.17) \qquad |||\psi_{\mathbf{h}}|||_{\frac{1}{2},\partial T_{\mathbf{h}}} \leq C_{\mathbf{k}} h^{-\frac{1}{2}} ||\psi_{\mathbf{h}}||_{0,\partial T_{\mathbf{h}}} \qquad \text{for all } \psi_{\mathbf{h}} \in L_{\partial T_{\mathbf{h}}}^{(k)}.$$

<u>THEOREM 20.3</u>. We suppose the hypotheses of Theorem 20.2 except here we fix for each $T \in T_h$

(20.18)
$$Y_T = P_k(\partial T)$$
, $Q_T = D_k^*(T)$.

Then there exists a constant C independent of h such that we have

$$(20.19) \qquad || \ \underline{p} - \underline{p}_h \, ||_{0, T_h} + ||| \phi - \phi_h \, |||_{\frac{1}{2}, \partial T_h} \leq C \, h^k \left\{ |\underline{p}|_{k, \Omega} + |u|_{k+1, \Omega} \right\}$$

once the solution (p,u) of (15.2) and of (15.8) belongs to $(H^k(\Omega))^2 \times H^{k+1}(\Omega)$; (\$\phi\$ is expressed as a function of u in (20.6)).

<u>Proof</u>. To apply the general theory developed in Chapter 3 we begin by making the change of variables \underline{p} to $\underline{p}-E_h^{(k)}\underline{p}$, and \underline{p}_h to $\underline{p}_h-E_h^{(k)}\underline{p}$. Thus we have

$$\bar{p} - E_h^{(k)} \bar{p} \in Q^0(T_h), \qquad \bar{p}_h - E_h^{(k)} \bar{p} \in W_h \subset Q^0 (T_h),$$

where $E_h^{(k)}p$ is the equilibrium interpolant of order k described in Section 6. The continuity properties and W_h -ellipticity with constants independent of h relative to the norm $\|\cdot\|_{0,\Omega}$ are obvious. Further, we can establish the existence of constants B and $\beta_0 > 0$, independent of h, such that

$$\begin{array}{ll} b(\underline{q},\psi) \leq B \mid\mid \underline{q}\mid\mid_{0,\Omega} \mid\mid\mid \psi\mid\mid\mid_{\frac{1}{2},\partial T_{h}} & \text{for all } \underline{q} \in Q^{0}(T_{h}) \\ & \text{and all } \psi \in L(\partial T_{h}) \end{array},$$

and

$$\sup_{\substack{q_h \in W_h}} \frac{b(\underline{q_h}, \psi_h)}{||\underline{q_h}||_{0,\Omega}} \ge \beta_0 \quad [\psi_h]_{\frac{1}{2},\partial T_h} \quad \text{for all } \psi_h \in M_h,$$

where the semi norm $[\cdot]_{\frac{1}{2},\partial T_h}$ is defined by

$$\left[\psi \right]_{\frac{1}{2}, \partial T_{\mathbf{h}}} = \left\{ \sum_{\mathbf{T} \in T_{\mathbf{h}}} \left| \left| \left| \psi_{\mathbf{T}} \right| - \frac{1}{\text{mes } \partial \mathbf{T}} \right| \left(\int_{\partial \mathbf{T}} \psi_{\mathbf{T}} \, d\sigma \right) \left| \left| \left| \left| \frac{1}{2}, \partial \mathbf{T} \right| \right|^{\frac{1}{2}}. \right.$$

Using the Poincaré-Friedrichs inequality we can show that $[\cdot]_{\frac{1}{2},\partial T_h}$ is a norm on M_{0h} equivalent (uniformly in h) to the norm $|||\cdot|||_{\frac{1}{2},\partial T_h}$. Thus there exists $\beta>0$, independent of h, such that

$$\sup_{\substack{q_h \in W_h}} \frac{b(\underline{q_h}, \psi_h)}{||q_h||_{0,\Omega}} \ge \beta |||\psi_h|||_{\frac{1}{2},\partial T_h} \qquad \text{for all } \psi_h \in M_{0h} ,$$

The theorem then follows easily.

<u>Remark 20.3</u>. When the integer k is even, we can obtain the same result while taking the space Q_T to be the space generated by $D_k(T)$ together with an appropriately chosen element of $D_{k+1}(T)$. For details and other examples cf. THOMAS (1976).

The dual hybrid finite element method corresponding to the choice of subspaces M_h and W_h associated with (20.18) by (20.8) and (20.9) is convergent for each integer $k \ge 1$ in the following sense: (with no regularity hypothesis for the solution, one can show that):

$$\lim_{h\to 0} \quad || \ \underline{p} - \underline{p}_h ||_{0,\Omega} = 0 .$$

The principle of the dual hybrid finite element methods given above for a second order elliptic problem corresponds to that developed by specialists in structural mechanics and termed assumed stress hybrid finite element model. We cite first of all the work of PIAN (1964), (1971), (1972), (1983); cf. also PIAN & TONG (1969), (1986), TONG (1983), and the recent articles of PIAN and his collaborators where can be found a treatment of nonaffine-equivalent finite elements. Other examples are given in ATLURI (1971), ATLURI, TONG & MURAKAWA (1983), PUNCH & ATLURI (1984), SPILKER & MUNIR (1980a), (1980b) and WOLF (1972a), (1972b) amoung many other publications on the subject. For a mathematical analysis of such methods we refer to the work of BREZZI (1975), (1977), of KIKUCHI (1973) and of QUARTERONI (1979).

Remark 20.4. The construction of the examples given above followed this strategy: given a space M_h , find a space W_h large enough for the compatibility condition to be satisfied. Under reasonable hypotheses we obtain asympotic orders of optimal error (relative to the given W_h). The construction of examples proposed by engineers often follows another strategy: given two

spaces W_h and M_h^* , the choice of spaces being consistent with the physics of the problem, we note the existence of "mechanisms" (... the homogeneous problem does not have the trivial solution as unique solution). One then tries to eliminate these mechanisms, that is to reduce M_h^* to a space M_h small enough for the compatibility condition with W_h to be satisfied. This technique does not lead to satisfactory asymptotic orders of convergence.

21 - HYBRIDIZATION OF EQUILIBRIUM METHODS

The dual hybrid finite element method just constructed yields - when the compatibility hypothesis (18.12) is satisfied - an approximation p_h of p_h satisfying div p_h + f_h = 0 in the interior of each finite element T; however, the reciprocity of the normal traces on the interfaces is not realized pointwise along the entire interface. In particular, the solution p_h is not in $H(\text{div};\Omega)$.

A variant introduced by FRAEIJS DE VEUBEKE (1965) yields an approximated solution $\underline{p}_h \in H(\text{div};\Omega)$ satisfying div $\underline{p}_h + f_h = 0$ in Ω . The mathematical analysis can be found in THOMAS (1977, Chapter VIII), cf. also RAVIART & THOMAS (1979). These results are completed by a postprocessing technique developed in ARNOLD & BREZZI (1985).

We thus reconsider the approximation of the solution (p,ϕ) of the model problem with the dual hybrid formulation (20.2) - (20.5). On the one hand, for each $T \in T_h$, we take a finite dimensional subspace $Q_T \subset H(\text{div};T)$ and as before, cf. (20.7), (20.8), define the spaces

$$(21.1) W_h^f = \left\{ \underline{q} \in H(\text{div}; T_h) ; \forall T \in T_h, \underline{q}_{|T} \in Q_T \right\} \cap Q^{f_h} (T_h)$$

and

(21.2)
$$W_h = W_h^0 = \{ q \in H(\text{div}; T_h) : \forall T \in T_h, q_{|T} \in Q_T \} \cap Q^0(T_h).$$

On the other hand, for each $T \in T_h$, we put

$$(21.3) Y_{\mathbf{T}} = \left\{ \psi \in L^{2}(\partial \mathbf{T}) ; \exists \mathbf{q} \in Q_{\mathbf{T}} \text{ s.t. } \mathbf{q} \cdot \nu_{\mathbf{T}} = \psi \right\}$$

and define

(21.4)
$$M_{h} = \left\{ \psi = (\psi_{T})_{T \in T_{h}} \in L^{2}(\partial T_{h}) ; \forall T \in T_{h}, \psi_{T} \in Y_{T} \text{ and } \right.$$

$$\forall T' = T_{1} \cap T_{2} \text{ with } T_{1}, T_{2} \in T_{h}, \psi_{T_{1}} = \psi_{T_{2}} \text{ on } T' \right\}.$$

We seek (p_h, ϕ_h) satisfying (cf. (20.10))

$$\left(\begin{array}{l} (\underline{p_h}, \phi_h) \in W_h^f \times M_h \ , \\ \\ a(\underline{p_h}, \underline{q_h}) \ + \ b(\underline{q_h}, \phi_h) \ = \ 0 \qquad \qquad \text{for all } \underline{q_h} \in W_h, \\ \\ b(\underline{p_h}, \psi_h) \ + \ d(\phi_h, \psi_h) \ = \ -\sum_{T \in T_h} \int_{\partial T \cap \Gamma} g \ \psi_{hT} \ d\sigma \qquad \text{for all } \psi_h \in M_h. \end{array} \right)$$

Remark 21.1 . Such an approximation method is not a dual hybrid approximation in the sense of the preceeding section because the spaces Y_T are not subspaces of $H^{\frac{1}{2}}(\partial T)$, and thus the space M_h of multipliers is not a subspace of the space $L(\partial T_h)$ defined in (18.4) or (18.5). Whereas examples of finite dimensional subspaces of $L^2(\partial T)$ and of $H^{-\frac{1}{2}}(\partial T)$ are identical, the examples of finite dimensional subspaces of $L^2(\partial T)$ and of $H^{\frac{1}{2}}(\partial T)$ are fundamentally different: the functions in $H^{\frac{1}{2}}(\partial T)$ which are locally continuous are globally continuous on ∂T .

It is clear that for each choice of spaces $Q_T \subset H(\text{div};T)$ with Y_T associated to Q_T by (21.3), problem (21.5) has a unique solution (p_h,ϕ) .

Suppose that the coefficient t>0 is constant on each face T'contained in the boundary Γ of Ω - (this hypothesis is made solely for simplicity of exposition). Then the elimination of the multiplier ϕ_h in (21.5) yields a characterization of p_h as the solution of the variational problem

with

$$(21.7) V_{h}^{f} = \left\{ \underline{q} \in H(\operatorname{div};\Omega); \operatorname{div} \underline{q} + f = 0 \text{ on } \Omega \right\}$$

$$\cap \left\{ \underline{q} \in \left(L^{2}(\Omega) \right)^{n} ; \forall T \in T_{h}, \underline{q}_{|T} \in Q_{T} \right\} ,$$

$$(21.8) V_{h} = V_{h}^{0} = \left\{ \underline{q} \in H(\operatorname{div};\Omega) ; \operatorname{div} \underline{q} = 0 \text{ on } \Omega \right\}$$

$$\cap \left\{ \underline{q} \in \left(L^{2}(\Omega) \right)^{n} ; \forall T \in T_{h}, \underline{q}_{|T} \in Q_{T} \right\} .$$

The approximation p_h of p thus obtained is an equilibrium approximation (cf. Section 2); i.e. p_h is a vector function in $H(\text{div};\Omega)$ that satisfies the equilibrium equation

(21.9)
$$\operatorname{div} \, p_h + f_h = 0$$
 in Ω .

Conversely, to an equilibrium formulation (21.6) - (21.8) can be associated a hybrid method defined by (21.1) - (21.5). Contrary to what one might suspect at first glance, it is simpler to solve numerically (21.5) than (21.6). The numerical solution of (21.5) can be carried out by first eliminating the principal variable \underline{p}_h . (This elimination is performed at the level of the n-simplex T). Then to calculate ϕ_h , one must solve a linear system of the same structure as the one which results from a nonconforming finite element method, cf., in particular, FRAEIJS DE VEUBEKE (1974a).

Remark 21.2 . When the given function f is trivial, or sufficiently simple that one can numerically reduce the problem to the case f = 0, to solve (19.6) (equilibrium formulation) it suffices to know how to construct a basis for V_h . In dimension n=2, using stream functions we can treat directly problem (21.6), cf. HASLINGER & HLAVACEK (1975), (1976a), (1976b), HLAVACEK (1980), and HLAVACEK & KRIZEK (1984). We have not chosen this presentation because it is not easily generalized in dimension n=3 to the case where f is arbitrary and the boundary condition is a Robin's condition.

When T_h is a triangulation of $\overline{\Omega}$ by n-simplexes T, if we take for spaces Q_T the spaces $D_k(T)$, where k is a given integer, $k \ge 1$, the described method is to take for spaces Y_T the spaces $D_k(\partial T)$. Thus here we have

$$(21.10) W_h^f = \left\{ \underline{q} \in H(\operatorname{div}; T_h) ; \forall T \in T_h, \underline{q}_{|T} \in D_k(T) \right\} \cap Q^{fh}(T_h),$$

and

$$(21.11) W_h = W_h^0 = \left\{ \underline{q} \in H(\text{div}; T_h) ; \forall T \in T_h , \underline{q}_{|_T} \in (P_{k-1}(T))^n \right\} \cap Q^0(T_h) ,$$

where for the characterization of W_h we have used the fact that each vector function in D_k that is divergence free belongs in fact to $(P_{k-1}(T))^n$. The associated space V_h is then none other than

(21.12)
$$V_h = \left\{ \underline{q} \in H(\operatorname{div};\Omega) ; \text{ div } \underline{q} = 0 \text{ on } \Omega \text{ and s.t. } \forall T \in T_h, \underline{q}_{|T} \in (P_{k-1})^n \right\},$$

i.e. with the notation of (6.17)

$$V_h = E_h^{(k)} \cap \left\{ \underline{q} \in \left(L^2(\Omega)\right)^n \; ; \; \text{div } \underline{q} = 0 \; \text{ on } \Omega \; \right\}.$$

For simplicity of exposition, to give error estimates we suppose that the family of triangulations T_h is uniformly regular. With $Q_T = D_k(T)$ and $Y_T = D_k(\partial T)$ for each $T \in T_h$, there is no difficulty in showing that there is a constant β_0 > 0, independent of h, such that

$$\sup_{\substack{q_h \in W_h}} \frac{\frac{b(q_h, \psi_h)}{||q_h||_{0,\Omega}} \ge \beta_0 h^{-\frac{1}{2}} \left[\psi_h\right]_{0,\partial T_h} \quad \text{for all } \psi_h \in M_h ,$$

where the semi-norm $\begin{bmatrix} \cdot \end{bmatrix}_{0,\partial T_h}$ is defined by

$$\left[\psi \right]_{0,\partial T_{\mathbf{h}}} = \left\{ \sum_{\mathbf{T} \in T_{\mathbf{h}}} ||\psi_{\mathbf{T}} - \frac{1}{\text{mes}\partial \mathbf{T}} \left(\int_{\partial \mathbf{T}} |\psi_{\mathbf{T}}| \, \mathrm{d}\sigma \right) ||_{0,\partial \mathbf{T}}^{2} \right\}^{\frac{1}{2}} .$$

Thus we have

<u>THEOREM 21.1</u> . Suppose that, for each $T \in T_h$,

(21.13)
$$Q_T = D_k(T), Y_T = D_k(\partial T).$$

Then there exists a constant C, independent of h, such that if (p_h, ϕ_h) is the solution of (21.5), then

$$(21.14) || p - p_h ||_{0,\partial T_h} + h^{-\frac{1}{2}} || \phi - \phi_h ||_{0,\partial T_h} \le Ch^k \{ |p|_{k,\Omega} + |u|_{k+1,\Omega} \},$$

provided, of course, that the solution (p,u) of (15.2) and of (15.8) belongs to $(H^k(\Omega))^n \times H^{k-1}(\Omega)$. (ϕ is expressed as a function of u in (20.6)).

Once the solution (p_h, ϕ_h) of problem (21.5), with the spaces Q_T and Y_T as given in (21.13), has been obtained, we can easily calculate the function $u_h \in L^2(\Omega)$ whose restriction to each $T \in T_h$ is the polynomial of degree at most k-1 defined by

$$\int_{T} u_{h} \operatorname{div} \, \underline{q} \, dx = -\sum_{i,j=1}^{n} \int_{T} A_{ij} \, p_{hj} \, q_{i} \, dx + \int_{\partial T} \phi_{hT} \, \underline{q} \cdot \underline{\nu}_{T} \, d\sigma$$

$$for \, each \, \underline{q} \in D_{k}(T).$$

Using that (p_h, ϕ_h) satisfies (21.5) we remark that this function u_h satisfies

$$\begin{split} &\int\limits_{\Omega} u_h \ \text{div} \ \ \underline{q}_h \ \ \text{dx} \ = \ -\sum\limits_{i,j=1}^n \ \int\limits_{\Omega} A_{ij} \ p_{hj} \ q_{hi} \ \ \text{dx} \ + \ \int\limits_{\Gamma} \frac{1}{t} \ (g - \underline{p}_h \cdot \underline{\nu}) \ \underline{q}_h \cdot \underline{\nu} \ \ \text{d}\sigma \end{split}$$
 for each $\underline{q}_h \in H(\text{div};\Omega)$ such that $\underline{q}_{h|T} \in D_k(T)$ for each $T \in T_h$.

The pair (p_h, u_h) is the unique solution of problem (15.14) (dual mixed formulation) obtained with (15.12) and (15.13). Thus we know u_h is a approximation of u of order k relative to the norm $||\cdot||_{0,\Omega}$.

Now for each $T \in T_h$, we have both the approximation $\phi_{hT} \in D_k(\partial T)$ of $\phi_T = u_{|\partial T}$ and the approximation $u_{h|T} \in P_{k-1}(T)$ of $u_{|T}$. These two approximations provide complementary information about $u_{|T}$ and it is reasonable to try to use all of this information to obtain a better approximation of $u_{|T}$. We take an example in dimension n=2 when k is odd. The relations

$$\int_{\partial T} (u_h^* - \phi_{hT}) \psi d\sigma = 0 \qquad \text{for all } \psi \in D_k(\partial T)$$

$$\int_{T} (u_h^* - u_h) v dx = 0 \qquad \text{for all } v \in P_{k-3}(T) \text{ (if } k \ge 3)$$

define a unique function $u_h^* \in L^2(\Omega)$ whose restriction to each triangle $T \in T_h$ is a polynomial of degree at most k. Under the standard regularity hypotheses, one can show (ARNOLD & BREZZI (1985, Theorem 2.2)) the order of approximation is thus improved:

$$\| \| \mathbf{u} - \mathbf{u}_{\mathbf{h}}^* \|_{0,\Omega} = \theta(\mathbf{h}^{k+1}).$$

and

CHAPTER VI

EXTENSIONS AND VARIATIONS

22 - OTHER EXAMPLES OF MIXED AND HYBRID METHODS

The mixed and hybrid finite element methods given in the preceding sections were developed almost exclusively for second order elliptic problems. Formulations of such methods for fourth order elliptic problems were noted. The mathematical framework constructed for the numerical analysis of these methods for approximating the solutions of partial differential aquations apparently ignores the physical origin of the model. We have not even sought to exploit the supplementary properties that might result from the fact in most cases we are concerned with the solutions of contrained minimization problems. This does not however imply that the elaboration of such approximation techniques can be successfully carried out without taking into account the physical properties of the underlying problem. There is no lack of examples in computational mechanics in the recent scientific literature. We give here only a brief overview.

Stokes' problem is a fundamental problem of fluid mechanics. The primary unknown u is the velocity vector for the fluid displacement. It can be characterized as being the solution of a minimization problem in a subspace of the space

$$V = \{v = (v_i) ; v_i \in H^1(\Omega), 1 \le i \le n, s.t. \text{ div } v = 0\}.$$

The Lagrangian multiplier associated with the constraint that the velocity field be divergence free is identifield with a pressure. The approximation of the solution of this problem has been and will continue to be the object of many publications. The state of the art in this domain can be found in the book of GIRAULT & RAVIART (1986) wich has recently appeared, and it does not seem opportune to duplicate here the bibliography of this work.

The system of linear elasticity constitutes another priviledged domain for the application of the theory of mixed and hybrid methods. The variational formulation of the particle displacement being taken as primal formulation, the dual formulation called the equilibrium formulation, is a variational formulation of the stress equation. This equilibrium formulation is a minimization principle in a subspace of the space

$$S = \{ \underset{\approx}{\sigma} = (\sigma_{ij}) ; \sigma_{ij} \in L^{2}(\Omega), \sigma_{ij} = \sigma_{ji}, 1 \le i, j \le n, \text{ div } \underset{\approx}{\sigma} = f \}.$$

With some simplifying restrictions concerning the data and geometry of the problem under consideration, it is possible to use Airy functions to reduce the study of the equilibrium formulation to that of the two-dimensional biharmonic problem. The degree of sophistication necessary for constructing equilibrium finite element models is more easily understood with this analogy. Examples can be found in JOHNSON & MERCIER (1978), (1979), KRIZEK (1982), (1983), and KRIZEK & NEITTAAMAKI (1986). Analyses of mixed and hybrid formulations for the system of elasticity are proposed in AMARA & THOMAS (1978), (1979), ARNOLD, BREZZI, & DOUGLAS (1985), ARNOLD, DOUGLAS, & GUPTA (1984), PITKARANTA & STENBERG (1983), and STENBERG (1986),(1987a), (1987b). The principal source of difficulty in the approximation of the stress formulation comes from the necessity of respecting, rigorously or not, the symmetry relations that should be satisfied by the stresses. The article of ARNOLD & FALK (1987) proposes a new formulation with one obtains distinct approximations of σ_{ij} and σ_{ij} :

Other recent articles concerning the numerical treatment of problems of solid mechanics with mixed and hybrid methods are referenced in the bibliography. While some of these articles may appear to lack a certain degree of mathematical rigor, experience has shown that the examples described therein can not, nevertheless, be neglected.

Another domain where quite a few models using mixed or hybrid finite elements have been constructed is the subdomain of fluid mechanics dealing with fluid flow in porous media and in particular petroleum reservoir engineering. These reservoir models are based on nonlinear equations too complex to be described here in a few lines. We refer the interested reader to the book of CHAVENT & JAFFRE (1986), where in Chapter V some models directly inspired by the mixed methods developed above may be found. Other models are described and analyzed in CHAVENT, COHEN, DUPUY, JAFFRE & RIBERA (1984), DOUGLAS (1984), DOUGLAS, EWING, & WHEELER (1983), DOUGLAS & ROBERTS (1983), DARLOW, EWING & WHEELER (1984), EYMARD, GALLOUET, & JOLY (1987), JAFFRE & ROBERTS (1985), and WHEELER & GONZALES (1984).

We conclude this section by mentioning that mixed finite elements have also been used for the equations of electromagnetism: NEDELEC (1980), BENDALI (1984), BOSSAVIT (1987).

23 - EXTENSIONS AND VARIATIONS OF THE THEORY

The methods of Lagrangian multipliers that have been presented here concern the approximation of elliptic problems. These methods have extended to the study of eigenvalue problems; cf. CANUTO (1978), MERCIER, OSBORN, RAPPAZ, & RAVIART (1981). They have also been extended to the study of problems of evolution: parabolic problems as in JOHNSON & THOMEE (1981) and in THOMEE (1984), hyperbolic problems as in CANUTO (1981), in DOUGLAS & GUPTA (1986), and in QUARTERONI (1980b).

Certain mixed methods have been successful only because of the use of reduced numerical integration. Some examples can be found in JOHNSON & PITKARANTA (1982), MALKUS & HUGUES (1978), NOOR & PETERS (1983), SHIMODAIRA (1985). Other examples of reduced integration are proposed in RAVIART & THOMAS (1977b) to obtain the interpretation of certain primal hybrid methods as nonconforming methods.

In the preceding studies the approximation of elliptic problems of constrained minimization has been approached exclusively with the technique of dualization of the constraints using Lagrangian multipliers. The technique of penalization of the constraints constitutes an alternative way of attacking these minimization problems. These techniques can fortunately be combined. The resulting technique is called the augmented Lagrangian technique for which we refer the reader to the book of FORTIN & GLOWINSKI (1983). Even though the penalization technique was not initially retained in the mixed and hybrid formulations of the problem, it can nonetheless be useful for the resolution of the linear system resulting from the mixed hybrid formulation; cf. BERCOVIER (1978), FELIPPA (1986); cf. also Section 17.

Other finite element methods such as nodal methods can be interpreted as being mixed or hybrid methods; cf. HENNART (1985), (1986), and HENNART, JAFFRE, & ROBERTS (1986). The so called Trefftz methods can also be understood as being variantes of the dual mixed hybrid formulations; cf. JIROUSEK & LAN GUEZ (1986) and ZIELINSKI & ZIENKIEWICZ (1985).

As an extension of the classical abstract theory of Ritz-Galerkin, with the development of the theory of internal approximation of variational problems which has been seen to be fundamental for the analysis of "standard" finite element methods, the abstract theory of Babuska-Brezzi for the approximation of problems with Lagrangian multipliers has become popular because of its applications to the analysis of mixed and hybrid finite element methods. Other

applications of the Babuska-Brezzi theory have been proposed more recently; cf; for example the analysis of spectral methods proposed in the articles of BERNARDI, CANUTO, & MADAY (1987), and of BERNARDI, MADAY, & METIVET (1987) or the analysis of modal decomposition methods in the articles of STOLARSKI & BELYTSCHKO (1986a), (1986b).

Domain decomposition methods for the solution of partial differential equations are concurrently enjoying a resurgence in popularity for their use in constructing algorithms well adapted to use on multiprocessing machines. The use of mixed and hybrid formulations seem adequate for this type of method; cf. WHEELER & GLOWINSKI (1987) for a first example in this direction.

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