# Jim Jr. Douglas <br> Todd Dupont <br> Peter Percell <br> Ridgway Scott <br> A family of $C^{1}$ finite elements with optimal approximation properties for various Galerkin methods for 2nd and 4th order problems 

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# A FAMILY OF C ${ }^{1}$ FINITE ELEMENTS WITH OPTIMAL APPROXIMATION PROPERTIES FOR VARIOUS GALERKIN METHODS FOR 2ND AND 4TH ORDER PROBLEMS (*) () 

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#### Abstract

Families of $C^{1}$ piecewise polynomial spaces of degree $r \geqq 3$ on triangles and quadrilaterals in two dimensions are constructed, and approximation properties of the families are studied. Examples of the use of the families in Galerkin methods for 2 nd and 4 th order elliptic boundary value problems on arbitrarily shaped domains are given. The approximation properties on the boundary are such that the rate of convergence of the Galerkin methods is the optimal rate determined by the degree $r$ of the piecewise polynomial space.

Résumé. - En dimension deux, on construit des familles d'espaces de classe $C^{1}$, formés de polynômes de degré $r \geqq 3$ par morceaux, sur des triangles et des quadrilatères, et on étudie les propriétés d'approximation de ces familles. On en donne des exemples d'application à des méthodes de Galerkin pour les problèmes aux limites elliptiques du $2^{e}$ et du $4^{e}$ ordre posés sur des domaines de forme arbitraire. Les propriétés d'approximation de la frontière sont telles que le taux de convergence des méthodes de Galerkin est le taux optimal, déterminé par le degré r de l'espace des polynômes par morceaux.


## INTRODUCTION

In finite element approximation of fourth order elliptic boundary value problems and in some Galerkin methods for second order problems, $C^{1}$ piecewise polynomial spaces are required in order to satisfy conformity. In addition, if the boundary is curved, essential boundary conditions must be closely approximated in a negative Sobolev boundary norm [2]. These two requirements rule out the use of standard triangular and quadrilateral elements, and the purpose of this note is to present families of macroelements that do possess the necessary smoothness and boundary condition approximation

[^0]properties. An additional bonus of the families is that they have minimal smoothness: e.g., the standard 21 degree of freedom $C^{1}$ quintic [34] is $C^{2}$ at vertices, whereas the elements described here comprise the full space of $C^{1}$ piecewise polynomials, with no additional constraints.

A basic idea of the Galerkin methods treated here is to interpolate the boundary conditions rather than attempt to satisfy them exactly, without imposing a penalty on the nonsatisfaction of boundary conditions. This idea was first studied in [2] and then amplified in [4] and [30]; in these papers, only second order Dirichlet problems using $C^{0}$ piecewise polynomial spaces on triangles were considered. Here, we treat fourth order problems, as well as general second order problems using Galerkin techniques [13], [15], [16] that require $C^{1}$ spaces. To achieve $C^{1}$ continuity while retaining enough degrees of freedom at the boundary to have the necessary approximation properties, macroelements based on the Clough-Tocher [12] and Fraeijs de Veubeke-Sander [18], [27] macrocubics are used.

The families of macroelements studied here were first considered in [32], and further developed in [21], [25], although accurate approximation of essential boundary conditions was not treated. In these works, the convergence parameter is the degree $r$ of the piecewise polynomials. The point of view in the present paper is that $r$ is fixed and the convergence parameter is the mesh size $h$.

## 1. FAMILIES OF MACROELEMENTS

In this section we present a family of $C^{1}$ triangular macroelements (also see [32]) for which we shall prove useful approximation properties. The family contains an element of degree $n$ for each $n \geqq 3$ beginning with the well known cubic Clough-Tocher element [12], [8]. We shall discuss several modifications of this basic family and also briefly discuss a very similar family of quadrilateral macroelements which begins with the cubic Fraeijs de VeubekeSander macroelement [18], [27], [11].

First we give some general definitions. Following Ciarlet [9], we define a finite element to be a triple $(K, F, \Sigma)$ such that
(a) $K \subset \mathbf{R}^{2}$ is a compact region having the restricted cone property,
(b) $F$ is a finite dimensional vector space of real-valued functions on $K$, and
(c) $\Sigma$ is a finite set of linear functionals $\varphi_{i}, 1 \leqq i \leqq N$, called the degrees of freedom of the finite element, which are defined on a vector space of functions containing $F$ and which have the property that for any real numbers $\alpha_{i}$, $1 \leqq i \leqq N$, there exists a unique function $f \in F$ which satisfies

$$
\varphi_{i}(f)=\alpha_{i}, \quad 1 \leqq i \leqq N
$$

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A nodal finite element is a finite element ( $K, F, \Sigma$ ) for which each degree of freedom is a functional which picks out the value or the value of some derivative at a point, called a node, in $K$. For $n \geqq 0$ and $U \subset \mathbf{R}^{2}$, let $P_{n}(U)$, denote the space of restrictions to $U$ of polynomials (in the coordinates of $\mathbf{R}^{2}$ ) of total degree not greater than $n$. The degree of a finite element $(K, F, \Sigma)$ is the largest integer $n$ such that $P_{n}(K) \subset F$ or $-\infty$ if $F$ does not contain $P_{n}(K)$ for any $n \geqq 0$.


Figure 1. - A macrotriangle $T$.

Now we define the family of triangular macroelements. For $n \geqq 3$, an element of degree $n$ in the family of triangular macroelements is a nodal finite element $\left(T, S_{n}(T), \Sigma_{n}\right)$ defined as follows:
(a) $T \subset \mathbf{R}^{2}$ is a macrotriangle; i.e., a triangle $T$ triangulated by three subtriangles $T_{1}, T_{2}$ and $T_{3}$ (see fig .1 ). Note that the point $b$ is allowed to be anywhere in the interior of $T$;
(b) $\mathrm{S}_{n}(T)=\left\{f \in C^{1}(T): f \mid T_{i} \in P_{n}\left(T_{i}\right), 1 \leqq i \leqq 3\right\}$, where the vertical bar denotes restriction of a function;
(c) The degrees of freedom $\Sigma_{n}$ are

1. the value and gradient (i.e. $\partial / \partial x$ and $\partial / \partial y$ ) at the exterior vertices,
2. the value at $n-3$ distinct points in the interior of each exterior edge of $T$,
3. the normal derivative at $n-2$ distinct points in the interior of each exterior edge of $T$,
and if $n \geqq 4$,
4. the value and gradient at the interior vertex,
5. the value at $n-4$ distinct points in the interior of each interior edge of $T$,
6. the normal derivative at $n-4$ distinct points in the interior of each interior edge of $T$, and
7. the value at $(1 / 2)(n-4)(n-5)$ distinct points in the interior of each $T_{i}$ chosen so that if a polynomial of degree $n-6$ vanishes at those points, then it vanishes identically.

For an example of one of these macroelements, seefig. 2 in which a dot denotes a value, a circle denotes a gradient and a dash denotes a normal derivative.


Figure 2. - The degrees of freedom $\Sigma_{6}$
For use in the proof of the following theorem, let $\lambda_{i}, 1 \leqq i \leqq 3$, be the (unique) affine function on $\mathbf{R}^{2}$ such that $\lambda_{i}(b)=1$ and $\lambda_{i}$ vanishes at the exterior vertices of $T_{i}$. Furthermore, let $\lambda$ be the function on $T$ defined by

$$
\lambda\left|T_{i}=\lambda_{i}\right| T_{i}, \quad 1 \leqq i \leqq 3
$$

Note that $\lambda$ is well defined and continuous on $T$ because $\lambda_{i}$ and $\lambda_{i+1}$ agree on the edge $E_{i+2}=T_{i} \cap T_{i+1}$ since both are zero at $a_{i+2}$ and one at $b$. (Here and through the following proof, subscripts referring to triangles are counted modulo 3).

Theorem 1: $\left(T, S_{n}(T), \Sigma_{n}\right), n \geqq 3$, is a well defined nodal finite element of degree $n$.

Proof: It is easy to see by simply counting that, for all $n \geqq 3$, the number of degrees of freedom, denoted by $\# \Sigma_{n}$, is $(3 / 2)\left(n^{2}-n+2\right)$. On the other hand, the following argument due to Strang (see [29], [31] for example) shows that

$$
\operatorname{dim} S_{n}(T) \geqq \frac{3}{2}\left(n^{2}-n+2\right)
$$

Let $R_{n}(T), n \geqq 3$, denote the space of all (discontinuous) piecewise polynomials of degree $n$ on the macrotriangle $T$. Note that $R_{n}(T)$ is isomorphic to $\prod_{j=1}^{3} P_{n}\left(T_{j}\right)$, so
$\operatorname{dim} R_{n}(T)=(3 / 2)(n+1)(n+2)$. Consider $S_{n}(T)$ to be a subspace of $R_{n}(T)$ consisting of functions satisfying certain constraints imposed across the interior edges of $T$. In general, it takes $2 n+1$ constraints ( $n+1$ on values and $n$ on normal derivatives) to make two polynomials of degree $n$ agree in a $C^{1}$ fashion along an edge. However, if a function in $R_{n}(T)$ satisfies the $2 n+1$ constraints along two of the interior edges of $T$, then it automatically has a uniquely defined value and gradient at the interior vertex of $T$. Thus it takes the satisfaction of no more than $3(2 n+1)-3=6 n$ constraints to ensure that a function in $R_{n}(T)$ is actually in $S_{n}(T)$. Therefore,

$$
\operatorname{dim} S_{n}(T) \geqq \operatorname{dim} R_{n}(T)-6 n=\frac{3}{2}\left(n^{2}-n+2\right)
$$

To show that $\left(T, S_{n}(T), \Sigma_{n}\right)$ is a well defined finite element we must show that $\Sigma_{n}$ is a basis for the dual space of $S_{n}(T)$, which we denote $\left(S_{n}(T)\right)^{\prime}$. But we have just shown that

$$
\# \Sigma_{n} \leqq \operatorname{dim} S_{n}(T)=\operatorname{dim}\left(S_{n}(T)\right)^{\prime}
$$

so it is enough to show that $\Sigma_{n}$ spans $\left(S_{n}(T)\right)^{\prime}$. This is equivalent to showing that if all the degrees of freedom for $f \in S_{n}(T)$ are zero, then $f$ vanishes identically.

When the degrees of freedom are all zero, in particular those of types 1,2 and 3, $f$ and $\nabla f$ vanish on the boundary of $T$. Thus $f=p \lambda^{2}$, where $p_{i}=p \mid T_{i}$ is a polynomial of degree $n-2$ and $p$ is continuous (because $f$ and $\lambda$ are continuous and because $\lambda$ does not vanish in the interior of $T$ ). Since $f \in C^{1}(T), \nabla f$ is well defined on $E_{i-1}$ and hence can be expressed there by both

$$
2 p \lambda \nabla \lambda_{i}+\lambda^{2} \nabla p_{i} \quad \text { and } \quad 2 p \lambda \nabla \lambda_{i+1}+\lambda^{2} \nabla p_{i+1}
$$

Therefore

$$
\begin{equation*}
2 p \nabla\left(\lambda_{i+1}-\lambda_{i}\right)+\lambda \nabla\left(p_{i+1}-p_{i}\right)=0 \quad \text { on } E_{i-1} \tag{1}
\end{equation*}
$$

Since $\lambda\left(a_{i-1}\right)=0$ and $\nabla\left(\lambda_{i+1}-\lambda_{i}\right) \neq 0$ (because the lines $\lambda_{i}=0$ and $\lambda_{i+1}=0$ cannot be parallel) it follows that $p\left(a_{i-1}\right)=0$. Now, suppose for a moment that $n=3$. (See also [8] and [26] for this case.) Then each $p_{i}$ is a linear polynomial which vanishes at the exterior vertices $a_{i-1}$ and $a_{i+1}$ of the triangle $T_{i}$, so $p_{i}$ vanishes on the entire exterior edge of $T_{i}$. Hence $p_{i}=c_{i} \lambda_{i}$ for some $c_{i} \in \mathbf{R}$, so $f_{i}=c_{i} \lambda_{i}^{3}$. Thus $\nabla f_{i}(b)=3 c_{i} \nabla \lambda_{i}(b)$. Since the vectors $\nabla \lambda_{i}(b), 1 \leqq i \leqq 3$, are linearly independent and $f$ is differentiable at $b$ [i.e., $\nabla f_{i}(b)=\nabla f_{i+1}(b)$ ], it follows that all the $c_{i}^{\prime}$ s are zero; i. e., $f$ vanishes identically. This finishes the case $n=3$. From now on we suppose that $n \geqq 4$. Then $p(b)=0$, since $f(b)=0$ and $\lambda(b)=1$, and $\nabla p_{i}(b)=0$ since

$$
0=\nabla f(b)=\left(2 p \lambda \nabla \lambda_{i}+\lambda^{2} \nabla p_{i}\right)(b)=\nabla p_{i}(b)
$$

Thus $p \mid E_{i-1}$ is a polynomial of degree $n-2$ which vanishes along with its derivative at $b$ and is zero at $a_{i-1}$ and the $n-4$ nodes of type 5 on $E_{i-1}$; consequently, $p \mid E_{i-1}$ is identically zero. Therefore,

$$
p_{i}=q_{i}\left(\lambda_{i+1}-\lambda_{i}\right)\left(\lambda_{i}-\lambda_{i-1}\right)
$$

and

$$
f_{i}=q_{i} \lambda_{i}^{2}\left(\lambda_{i+1}-\lambda_{i}\right)\left(\lambda_{i}-\lambda_{i-1}\right),
$$

where $q_{i}$ is a polynomial of degree $n-4$, since $p_{i}=0$ on the two interior edges of $T_{i}$ and since $\lambda_{i+1}-\lambda_{i}$ vanishes on one of these edges while $\lambda_{i}-\lambda_{i-1}$ vanishes on the other. Also, using (1) and the fact that $p=0$ on $E_{i-1}$, we see that $\nabla p_{i}=\nabla p_{i+1}$ on $E_{i-1}$. But on $E_{i-1}$,

$$
\nabla p_{i}=q_{i}\left(\lambda_{i}-\lambda_{i-1}\right) \nabla\left(\lambda_{i+1}-\lambda_{i}\right)
$$

and

$$
\nabla p_{i+1}=q_{i+1}\left(\lambda_{i-1}-\lambda_{i+1}\right) \nabla\left(\lambda_{i+1}-\lambda_{i}\right) ;
$$

so since $\lambda_{i}=\lambda_{i+1}$ on $E_{i-1}$, we find that $q_{i}=-q_{i+1}$ on $E_{i-1}$. In particular, $q_{i}(b)=-q_{i+1}(b)$, and

$$
q_{i}(b)=-q_{i+1}(b)=q_{i+2}(b)=-q_{i}(b)
$$

which implies that $q_{i}(b)=0$. At this point we are finished if $n=4$ because then $q_{i}$ is a constant polynomial ( $c f$. [26]). If $n \geqq 5$, then $\nabla f_{i}=0$ at the $n-4$ nodes of type 6 on $E_{i \pm 1}$ because $f=p \lambda^{2}=0$ on $E_{i \pm 1}$. Since

$$
\nabla f_{i}=q_{i} \lambda_{i}^{2}\left(\lambda_{i}-\lambda_{i \pm 1}\right) \nabla\left(\lambda_{i \mp 1}-\lambda_{i}\right) \quad \text { on } E_{i \pm 1}
$$

it follows that $q_{i}$ vanishes at the same $n-4$ nodes on $E_{i+1}$ and $E_{i-1}$. As it also vanishes at $b$ and is a polynomial of degree $n-4$, it vanishes identically on $E_{i+1}$ and $E_{i-1}$. If $n=5$, this means that $q_{i} \equiv 0$, since it is linear. For $n \geqq 6$, this means that

$$
q_{i}=r_{i}\left(\lambda_{i+1}-\lambda_{i}\right)\left(\lambda_{i}-\lambda_{i-1}\right)
$$

and

$$
f_{i}=r_{i} \lambda_{i}^{2}\left(\lambda_{i+1}-\lambda_{i}\right)^{2}\left(\lambda_{i}-\lambda_{i-1}\right)^{2}
$$

where $r_{i}$ is a polynomial of degree $n-6$. Finally, since $f_{i}$ is zero at the nodes of type 7 , it follows that $r_{i}$ is also zero at these nodes and hence $r_{i}$ must vanish identically because it has degree $n-6$. //

Remark 1: Let $\alpha_{i}$ : $\left[0, L_{i}\right] \rightarrow T_{i}, 1 \leqq i \leqq 3$, be a parametrization by arclength of the exterior edge of $T_{i}$ and let $\partial / \partial n_{i}$ denote differentiation normal to this edge. If for some $i, 1 \leqq i \leqq 3$, the degrees of freedom along the exterior edge of $T_{i}$ of
types 2 and 3 in the set $\Sigma_{n}$ are replaced by
2.

$$
f \rightarrow \int_{0}^{L_{i}} p_{j}(s) f\left(\alpha_{i}(s)\right) d s, \quad 0 \leqq j \leqq n-4
$$

and
3. $\quad f \rightarrow \int_{0}^{L_{i}} q_{j}(s) \frac{\partial f}{\partial n_{t}}\left(\alpha_{i}(s)\right) d s, \quad 0 \leqq j \leqq n-3$,
where $\left\{p_{j}\right\}$ (resp. $\left\{q_{j}\right\}$ ) is a basis for the space of polynomials in one variable of degree $n-4$ (resp. $n-3$ ), then the result is a new (non-nodal) finite element of degree $n$, which we denote $\left(T, S_{n}(T), \tilde{\Sigma}_{n}\right)$. This follows from the proof of theorem 1 because the number of degrees of freedom is unchanged and because the degrees of freedom of types $1, \tilde{2}$ and $\tilde{3}$ again uniquely determine $f$ and $\nabla f$ on an exterior edge of $T$ when $f \in S_{n}(T)$. Finite elements using degrees of freedom of type $\tilde{2}$ were studied by Blair [4] in the context of second order problems. If the sets $\left\{p_{j}\right\}$ and $\left\{q_{j}\right\}$ are chosen to be orthogonal polynomials, a hierarchical structure may be achieved [25]. //

Remark 2: In both the original element ( $\left.T, S_{n}(T), \Sigma_{n}\right)$ and the modified element ( $T, S_{n}(T), \tilde{\Sigma}_{n}$ ), the exterior edges of $T$ need not be straight but can bè smooth curves which are just $C^{1}$ close to being straight. Using the notation introduced above, we say that the exterior edge of $T_{i}$ is $C^{1}$ close to being straight if

$$
\sup _{0 \leq s \leq L_{i}}\left\{\left|\lambda_{i}\left(\alpha_{i}(s)\right)\right|,\left|\frac{d}{d s}\left[\lambda_{i}\left(\alpha_{i}(s)\right)\right]\right|\right\}
$$

is small. It follows immediately from this definition that a $C^{1}$ small perturbation of a straight edge moves points on the edge and normals to the edge only slightly. Thus if $T$ has exterior edges which are sufficiently $C^{1}$ close to being straight, then the degrees of freedom $\Sigma_{n}$ or $\tilde{\Sigma}_{n}$ are close enough to those for the corresponding element for the macrotriangle having straight edges and the same vertices so that the degrees of freedom still determine uniquely a function $f \in S_{n}(T)$. (The distance between two degrees of freedom is measured in the dual space of $C_{b}^{1}$, the Banach space of bounded $C^{1}$ functions with bounded first derivatives.) //

Remark 3: The elements in theorem 1 have the virtue that they may be pieced together to form $C^{1}$ functions. For this reason, they may be referred to as " $C$ " finite elements". To see why this is so, let $T^{1}$ and $T^{2}$ be two macro-triangles that share (only) a common edge $E$ (the vertices of $E$ are required to be vertices of both $T^{1}$ and $T^{2}$ ). Let $\Sigma_{n}^{i}$ be degrees of freedom defined on $S_{n}\left(T^{i}\right), i=1,2$, that
are consistent, i.e., the nodal points of type 2 and 3 on E for $\Sigma_{n}^{1}$ and $\Sigma_{n}^{2}$ pairwise coincide. Let $f_{i} \in S_{n}\left(T^{i}\right), i=1,2$, be such that the degrees of freedom of types $1-3$ of $f_{1}$ and $f_{2}$ pairwise coincide, and let $f$ be defined on $T^{1} \cup T^{2}$ by $f \mid T^{i}=f_{i}, i=1,2$. Then $f \in C^{1}\left(T^{1} \cup T^{2}\right)$. (Proof: Let $\tau_{i} \subset T^{i}, i=1,2$, be the subtriangle having $E$ as an edge. The polynomial $P=f_{1}\left|\tau_{1}-f_{2}\right| \tau_{2}$ vanishes to second order at the vertices of $E$ and at $n-3$ other points on $E$, and since the degree of $P$ is $n, P \mid E \equiv 0$. Similar reasoning shows that the normal derivative of $P$ vanishes on $E$. Thus $P$ vanishes to second order on $E$, and this means that $f$ is $C^{1}$.)

The elements in remark 1 can also be pieced together to give $C^{1}$ functions, requiring only a matching of the orientations of $E$. Elements with curved edges as in remark 2 will arise only when the curved edge lies on the boundary of a domain; thus the problem of piecing together elements across a curved edge is avoided. Elements of all three types can be attached to each other across straight edges by matching the types of degrees of freedom on the shared edges. In fact, while the degrees of freedom of types 2 and 3 , or $\tilde{2}$ and $\tilde{3}$, will be required in section 2 on the boundary of the domain, other finite elements may be used in the interior, with the appropriate matching. For $n \geqq 5$, there is a well known (cf. [34] for the case $n=5$ ) $C^{1}$ finite element of degree $n$, which we shall denote by ( $\tau, P_{n}(\tau), \bar{\Sigma}_{n}$ ), such that $\tau$ is an ordinary triangle and the degrees of freedom $\bar{\Sigma}_{n}$ are

1. the value plus all first and second derivatives at each vertex,
2. the value at $n-5$ distinct points on each edge,
3. the normal derivative at $n-4$ distinct points on each edge, and
4. the value at $(1 / 2)(n-4)(n-5)$ distinct points in the interior of $\tau$ chosen so that if a polynomial of degree $n-6$ vanishes at those points, then it vanishes identically.

The transition from this element to the one in theorem 1 is via two unsymmetric macroelements, denoted ( $\left.T, S_{n}^{\prime}(T), \Sigma_{n}^{\prime}\right)$ and ( $\left.T, S_{n}^{\prime \prime}(T), \Sigma_{n}^{\prime \prime}\right)$, for which $S_{n}^{\prime}(T)$ and $S_{n}^{\prime \prime}(T)$ are proper subspaces of $S_{n}(T)$ consisting of functions which have second derivatives at certain exterior vertices of $T$, and whose degrees of freedom when $n=6$ are presented in figure 3 . In this figure a second circle around a node means that all the second derivatives at that point are degrees of freedom. We shall not write out the degrees of freedom $\Sigma_{n}^{\prime}$ and $\Sigma_{n}^{\prime \prime}$ in detail because the pattern should be clear. Proofs that these elements are well defined can be given along the same lines as in the proof of theorem 1. Note that remarks 1 and 2 hold for these elements. //

We now quickly present the family of $C^{1}$ quadrilateral macroelements, discussing just the differences between this family and the family of triangular
macroelements. For $n \geqq 3$, an element of degree $n$ in the family of quadrilateral macroelements is a nodal finite element $\left(Q, S_{n}(Q), \Lambda_{n}\right)$ such that
(a) $Q \subset \mathbf{R}^{2}$ is a convex quadrilateral triangulated by the four triangles $Q_{i}$, $1 \leqq i \leqq 4$, obtained by drawing in the diagonals of $Q$,


Figure 3. - Triangular transition elements of degree 6.


Figure 4. - The degrees of freedom $\Lambda_{6}$
(b) $S_{n}(Q)=\left\{f \in C^{1}(Q): f \mid Q_{i} \in P_{n}\left(Q_{i}\right), 1 \leqq i \leqq 4\right\}$, and
(c) after the obvious change of $T^{\prime}$ s to $Q^{\prime} \mathrm{s}$, the description of $\Lambda_{n}$ is the same as the description of $\Sigma_{n}$, except that (when $n \geqq 4$ ) an extra normal derivative at a point along just one of the interior edges of $Q$ is added to the degrees of freedom of type 6 .

The proof that $\left(Q, S_{n}(Q), \Lambda_{n}\right)$ is a well defined nodal finite element of degree $n$ is essentially the same as the proof of theorem 1 with only two changes of any significance needed. The first change is necessary because the interior vertex in the quadrilateral macroelement is singular in the sense of [29]. Counting shows that $\# \Lambda_{n}=2\left(n^{2}-n+2\right)$, while the methods used in the proof of theorem 1 only show that

$$
\operatorname{dim} S_{n}(Q) \geqq 2\left(n^{2}-n+2\right)-1
$$

However, in order to be able to proceed as before we need to know that $\# \Lambda_{n} \leqq \operatorname{dim} S_{n}(Q)$. In [29], it is shown that this is true since one of the constraints on normal derivatives one expects to be necessary is actually redundant in the


Figure 5. - Quadrilateral transition element of degree 6.
presence of the singular vertex. The only other change needed in the proof comes when one is trying to show that the $q_{i}$ 's are zero at the interior vertex. The alternating sign approach used before breaks down for the quadrilateral macroelements because an even number of edges abut the interior vertex. However, it is easy to see that the proof is rescued by the extra degree of freedom of type 6 in $\Lambda_{n}$ available when $n \geqq 4$.

Remarks 1 and 2 clearly apply to this family of quadrilateral macroelements, and the ideas discussed in remark 3 may be extended to the present situation as follows. First note that triangular and quadrilateral macroelements of the same degree are compatible and can be used together in the same domain $\Omega$ because the degrees of freedom associated with exterior edges are the same for both types of element. One may thus use the standard $C^{1}$ triangular element ( $\tau, P_{n}(\tau), \bar{\Sigma}_{n}$ ), $n \geqq 5$, as before, in the interior of $\Omega$ and a layer of quadrilateral macroelements
one element thick along $\partial \Omega$. Only one type of boundary layer quadrilateral macroelement is required for the transition, and this is depicted in figure 5 for $n=6$. Transition elements of other degrees may be easily constructed using the ideas discussed earlier. Note that the transition element in figure 5 may be used as a boundary element directly without having an additional layer of elements of the type depicted in figure 4, if desired.

## 2. APPROXIMATION PROPERTIES

In this section we prove approximation properties for the finite elements introduced in section 1. We start by fixing some notation to be used throughout the remainder of this section.

For any domain $\Omega \subset \mathbf{R}^{2}$ having the restricted cone property and for any smooth $\operatorname{arc} \Gamma \subset \mathbf{R}^{2}$ parametrized by arclength, denoted by $s$, let

$$
(v, w)_{\Omega}=\int_{\Omega} v(x) w(x) d x
$$

be the (real) inner product in $L^{2}(\Omega)$ and let

$$
\langle v, w\rangle_{\Gamma}=\int_{\Gamma} v(s) w(s) d s
$$

be the inner product in $L^{2}(\Gamma)$. Furthermore, if $m$ is a non-negative integer, let $[.]_{m, \Omega}$ denote the seminorm defined by

$$
[v]_{m, \Omega}^{2}=\sum_{|\alpha|=m}\left(\mathrm{D}^{\alpha} v, \mathrm{D}^{\alpha} v\right)_{\Omega}
$$

let $\|\cdot\|_{m, \Omega}$ be the norm for the Sobolev space $H^{m}(\Omega)$ defined by

$$
\|v\|_{m, \Omega}^{2}=\sum_{j=0}^{m}[v]_{, \Omega}^{2} .
$$

When $m$ is a nonintegral positive real number, let $\bar{m}$ denote the integral part of $m$ and define the semi-norm $[.]_{m, \Omega}$ by

$$
[v]_{m, \Omega}^{2}=\sum_{|x|=\bar{m}} \iint_{\Omega \times \Omega} \frac{\left|D^{\alpha} v(x)-D^{\alpha} v(y)\right|^{2}}{|x-y|^{2+2(m-\bar{m})}} d x d y
$$

define the norm $\left\|\|_{m, \Omega}\right.$ for $H^{m}(\Omega)$ by

$$
\|v\|_{m, \Omega}^{2}=\|v\|_{m, \Omega}^{2}+[v]_{m, \Omega}^{2}
$$

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Let $\|\cdot\|_{-m . \Omega}$ denote the dual norm for the space $H^{-m}(\Omega)=\left(H^{m}(\Omega)\right)^{\prime}$ defined by

$$
\|v\|_{-m, \Omega}=\sup \left\{\frac{\left|(v, \varphi)_{\Omega}\right|}{\|\varphi\|_{m, \Omega}}: \varphi \in H^{m}(\Omega), \varphi \neq 0\right\}
$$

Similarly, for $m$ a non-negative integer, let $|\cdot|_{m, \Gamma}$ denote the norm for $H^{m}(\Gamma)$ defined by

$$
|v|_{m, \Gamma}^{2}=\sum_{j=0}^{m}\left\langle\frac{\partial^{j} v}{\partial s^{j}}, \frac{\partial^{j} v}{\partial s^{j}}\right\rangle_{\Gamma} .
$$

For nonintegral positive real $m$, define the norm $|\cdot|_{m, \Gamma}$ for $H^{m}(\Gamma)$ by

$$
|v|_{m, \Gamma}^{2}=|v|_{\tilde{m}, \Gamma}^{2}+\iint_{\Gamma \times \Gamma} \frac{\left|v^{(m)}(s)-v^{(m)}(t)\right|^{2}}{d(s, t)^{1+2(m-\tilde{m})}} d s d t
$$

where $v^{(m)}(s)=\partial^{\bar{m}} v / \partial s^{\bar{m}}(x(s))$ with $x(s)$ parametrizing $\Gamma$ and $d(s, t)$ is the minimum arc length on $\Gamma$ between $x(s)$ and $x(t)$. Let

$$
|v|_{-m, \Gamma}=\sup \left\{\frac{\left|\langle v, \varphi\rangle_{\Gamma}\right|}{|\varphi|_{m, \Gamma}}: \varphi \in H^{m}(\Gamma), \varphi \neq 0\right\}
$$

be the norm for $H^{-m}(\Gamma)=\left(H^{m}(\Gamma)\right)^{\prime}$. Finally, for $m$ a non-negative integer, let $\{.\}_{m, \Omega}$ denote the $C^{m}(\Omega)$ norm defined by

$$
\{v\}_{m, \Omega}=\sup \left\{\left|D^{\alpha} v(x)\right|: x \in \Omega, 0 \leqq|\alpha| \leqq m\right\}
$$

Let $\Omega$ be a fixed, bounded domain in $\mathbf{R}^{2}$ with $C^{\infty}$ boundary $\partial \Omega$ parametrized by arclength. A triangulation of $\Omega$ will mean a collection $\mathscr{I}=\left\{T_{1}, \ldots, T_{N}\right\}$ of polygons each of which is a triangle, macrotriangle, or convex quadrilateral such that $\bar{\Omega}=\bigcup_{j=1}^{N} T_{j}$ and such that distinct polygons intersect at most in a common vertex or a common edge. By a simple triangle in such a triangulation, we mean either a triangle in $\mathscr{I}$ (but not a macrotriangle), a subtriangle of a macrotriangle, or a traingle obtained by subdividing a convex quadrilateral by drawing its diagonals. Simple triangles are assumed to have straight edges except that an edge between two vertices on $\partial \Omega$ is assumed to be contained in $\partial \Omega$.

Now suppose that $\Omega$ is provided with a family of triangulations $\mathscr{I}_{h}$, $0<h \leqq h_{0} \leqq 1$, as described above, such that
(a) if $T \in \mathscr{I}_{h}$, then diam $(T) \leqq h$, and
(b) the family is non-degenerate in the sense that there exists a number $\rho>0$ such that if $T$ is a polygon in the family and $\tau \subset T$ is a simple triangle, then $\tau$ is starlike with respect to a disk whose diameter is $\rho$ diam $(T)$.

Consider for a moment an arbitrary finite element ( $K, F, \Sigma$ ) with degrees of freedom $\varphi_{1}, \ldots, \varphi_{N}$ that are continuous on $H^{m}(K)$. We associate with ( $K, F, \Sigma$ ) an interpolation operator

$$
\Pi: \quad H^{m}(K) \rightarrow F
$$

by letting $\Pi v$ be the unique function in $F$ such that

$$
\varphi_{i}(\Pi v)=\varphi_{i}(v), \quad i=1, \ldots, N .
$$

Let $\delta=\operatorname{diam}(K)$ and let $\hat{x}=x / \delta$ and $\tilde{K}=\{\hat{x}: x \in K\}$.
For any function $f$ defined on $K$, set $\hat{f}(\hat{x})=f(x)$, and define an operator $\hat{\Pi}$ on $H^{m}(\hat{K})$ by

$$
\hat{\Pi} \hat{v}=\Pi v, \quad v \in \mathrm{H}^{m}(K)
$$

Note that $\hat{\Pi}$ is the interpolation operator associated with the finite element ( $\hat{K}, \hat{F}, \hat{\Sigma}$ ), and that, for any of the finite elements discussed in section 1 , including that of remark 1, $\hat{\Sigma}$ has the same description as $\Sigma$ (because dilations-as opposed to arbitrary linear maps - preserve orthogonality).

Now, suppose $h_{0}$ is small enough (see remark 2 of section 1) so that each polygon $T_{h, j} \in \mathscr{I}_{h}, 1 \leqq j \leqq N_{h}$, can be given the structure $\left(T_{h, j}, F_{h, j}, \Sigma_{h, j}\right)$ of one of the finite elements of section 1 of some fixed degree $r \geqq 3$ for $0<h \leqq h_{0}$. Let $\Pi_{r, h, j}$ be the interpolation operator associated with ( $\left.T_{h, j}, F_{h, j}, \Sigma_{h, j}\right)$ and let $m_{0}$ be a real number such that each $\varphi \in \Sigma_{h, j}$ is continuous on $H^{m_{0}}\left(T_{h, j}\right)$ for all $j=1, \ldots$, $N_{h}$ and $0<h \leqq h_{0}$. We assume that the following uniformity condition holds:

$$
\begin{equation*}
\left\|\hat{\Pi}_{r, h, \vec{j}} \hat{v}\right\|_{2, \hat{r}_{m},} \leq C\|\hat{v}\|_{m_{0} \cdot \hat{T}_{m},} \tag{c}
\end{equation*}
$$

for $1 \leqq j \leqq N_{h}, 0<h \leqq h_{0}$, and $\hat{v} \in H^{m_{0}}\left(\hat{T}_{h, j}\right)$. (Here and throughout the rest of this section, unless otherwise stated, $C$ denotes a generic constant which may be different in different places and which depends only on $\Omega, r, h_{0}$, and the families of triangulations and interpolation operators.) The existence of such a constant $C$ follows from Sobolev's inequality ( $c f$. Grisvard [19]), with $m_{0}$ any real number greater than one plus the highest order of derivative occurring as a degree of freedom (hence note that $m_{0}>2$ always). The condition can be satisfied with $C$ independent of $h$ because the family of triangulations is non-degenerate: there is room in the normalized (hatted) polygons for degrees of freedom that are bounded away from degenerate configurations. Examples of degeneracy are coalescence of two value nodes or colinearity of three value nodes in the interior of a simple triangle associated with an element of degree 7 . Now define a global interpolation operator $\Pi_{r, h}$ on $H^{m_{0}}(\Omega)$ by

$$
\left(\Pi_{r, h} v\right) \mid T_{h, j}=\Pi_{r, h, j}\left(v \mid T_{h, j}\right), \quad j=1, \ldots, N_{h} .
$$

Let $S_{r, h}$ be the image of $H^{m_{0}}(\Omega)$ via the mapping $\Pi_{r, h}$. We assume that the family of elements is consistent (see remark 3 of section 1 ), namely, we assume that
(d) $S_{r, h} \subset C^{1}(\bar{\Omega})$.

Since each function in $S_{r, h}$ is a polynomial in each simple triangle in $\mathscr{I}_{h}$, condition (d) implies that $S_{r, h} \subset H^{2}(\Omega)$.

Remark 4: When each element in $\mathscr{I}_{h}$ is one of the macroelements of degree $r$ of section 1 (except for the transition macroelements described in figures 3 and 5), then $S_{r, h}$ is the space of all $C^{1}$ piecewise polynomials of degree $r$ on the triangulation $\bar{I}_{h}$ of $\Omega$ obtained by considering all simple triangles in $\mathscr{I}_{h}$ separately. As shown in [23], when $r \geqq 5$, this space has a nodal basis, but of a more complicated variety than considered in section 1: the macroelements provide a simpler nodal parametrization of $S_{r, h}$. In the case $r \geqq 5$, the theorems to follow may be generalized to hold for $S_{r, h}\left(\overline{\mathscr{I}}_{h}\right)$ for an arbitrary triangulation $\overline{\mathscr{I}}_{h}$ of $\Omega$, i. e., one not coming from a macro-triangulation, by changing the nodal parameters in [23] at the boundary appropriately. The essential change required is algebraic, as the analysis of this section involves only the simple triangles. //

For the remainder of this section, we make the following

Assumption: For $0<h \leqq h_{0}$, the family of triangulations $\mathscr{I}_{h}$ of $\Omega$ satisfies conditions $(a)$ and $(b)$ and the associated family of interpolation operators $\Pi_{r, h}$, $r \geqq 3$, satisfies conditions (c) and (d).

Theorem 2: Let $q$ and $m$ be real numbers such that $0 \leqq q \leqq 2$ and $m_{0} \leqq m \leqq r+1$, and let $v \in H^{m}(\Omega)$. Then

$$
\left\|v-\Pi_{r, h} v\right\|_{q, \Omega} \leqq C h^{m-q}[v]_{m, \Omega} .
$$

Proof: Let $T_{h, j} \in \mathscr{I}_{h}$. From the fact that $\hat{\Pi}_{r, h, j}$ preserves polynomials of degree $r$ and the uniformity condition (c), it follows that, for any $\varphi \in P_{r}\left(\hat{T}_{h, j}\right)$,

$$
\left\|\hat{v}-\hat{\Pi}_{r, h, j} \hat{v}\right\|_{q, \hat{T}_{h, j}}=\left\|(\hat{v}-\varphi)-\hat{\Pi}_{r, h, j}(\hat{v}-\varphi)\right\|_{q, \hat{T}_{n, j}} \leqq C\|\hat{v}-\varphi\|_{m, \hat{T}_{h, j}}
$$

The following version (proved in [17] and [20]) of the standard BrambleHilbert lemma (see [5] and [10]) is now required:

Lemma 1: Given $0<\rho<1$ and a positive real number $m$, there exists a constant $C=C(\rho, m)$ such that if $K \subset \mathbf{R}^{2}$ is a domain having diameter at most one that is
starlike with respect to a disc of diameter $\rho$ and if $v \in H^{m}(K)$, then

$$
\inf _{\varphi \in P_{\bar{m}}(K)}\|v-\varphi\|_{m, K} \leqq C[v]_{m, K}
$$

where $\bar{m}$ is the greatest integer less than $m$. //
Applying the lemma,

$$
\left\|\hat{v}-\hat{\Pi}_{r, h, j} \hat{v}\right\|_{q, T_{h, j}} \leqq C[\hat{v}]_{m . T_{h},} .
$$

Since diam $\left(T_{h, j}\right) \leqq h$, the homogeneity of the norms implies that

$$
\left\|v-\Pi_{r, h} v\right\|_{q, \hat{r}_{h},} \leqq C h^{m-q}[v]_{m, \grave{T}_{h, j}} \cdot \quad / /
$$

From now on, we shall make frequent use of the following easily proved inequalities: if $K \subset \mathbf{R}^{2}$ is a region such that diam $(K) \leqq \delta \leqq 1$ and $K$ contains a disc of diameter $\rho \delta, \rho>0$, then for $\varphi \in P_{r}(K)$ and any $m$,

$$
\begin{equation*}
\{\varphi\}_{m, K} \leqq C \delta^{-1}\|\varphi\|_{m, K} \tag{2}
\end{equation*}
$$

and

$$
\begin{equation*}
\|\varphi\|_{m, K} \leqq C \delta^{j-m}\|\varphi\|_{J, K}, \quad 0 \leqq j \leqq m \tag{3}
\end{equation*}
$$

where $C$ depends only on $\rho$ and $r$. The following lemma will be used in the proofs of the remaining theorems of this section (for a proof, see the appendix of [30], or use the trace theorem [22] with a little care).

Lemma 2: Let $B\left(\mathscr{I}_{h}\right)$ be the collection of all simple triangles in $\mathscr{I}_{h}$ having an edge on $\partial \Omega$ and let $v \in H^{m}(\Omega)$ with $m_{0} \leqq m \leqq r+1$. Then

$$
\left(\sum_{\tau \in B\left(\xi_{n}\right)}\left\|\Pi_{r, h} v\right\|_{m-1, \tau}^{2}\right)^{1 / 2} \leqq C h^{1 / 2}\|v\|_{m, \Omega} \quad / /
$$

Theorem 3: Let $v \in H^{m}(\Omega)$ for some $m$ satisfying $m_{0} \leqq m \leqq r+1$. Then if $v \in H^{m}(\partial \Omega)$,
(a)

$$
\left|v-\Pi_{r, h} v\right|_{q, \partial \Omega} \leqq C h^{m-q}\left(|v|_{m, \partial \Omega}+\|v\|_{m, \Omega}\right)
$$

when $0 \leqq q \leqq 2$. If $\partial v / \partial n \in H^{m-1}(\partial \Omega)$, then
(b) $\quad\left|\frac{\partial v}{\partial n}-\frac{\partial}{\partial n}\left(\Pi_{r, h} v\right)\right|_{q, \partial \Omega} \leqq C h^{m-q-1}\left(\left|\frac{\partial v}{\partial n}\right|_{m-1, \partial \Omega}+\|v\|_{m, \Omega}\right)$
when $0 \leqq q \leqq 1$, where $\partial / \partial n$ denotes differentiation normal to $\partial \Omega$.
(*) Here $m$ is assumed to be a non-negative integer for $\left\}_{m, k}\right.$ to be defined.

Proof: The proofs of part (a) and (b) are similar, so we just prove (a). Let $T \in \mathscr{J}_{h}$ be an element with an edge $e$ on $\partial \Omega$, let $\delta=\operatorname{diam} T$ and let $\tau \subset T$ be the simple triangle having $e$ as an edge. Note that since $\partial \Omega$ is smooth, the length of $e, l(e)$, may be bounded by $C \delta$. Let $\varphi$ be the polynomial in arclength on $e$ of degree $r$ defined using the degrees of freedom associated with values, tangential derivatives, or orthogonalities on $e$ in the definition of $\Pi_{r, h}$ which interpolates both $v \mid e$ and $\Pi_{r, h} v \mid e$. Then

$$
\begin{aligned}
\left|v-\Pi \Pi_{r, h} v\right|_{q, e} \leqq|v-\varphi|_{q, e}+\mid \varphi- & \left.\Pi_{r, h} v\right|_{q, e} \\
& \leqq C\left(\delta^{m-q}|v|_{m, e}+\delta^{r-q+1}\left|\frac{\partial^{r+1}}{\partial s^{r+1}}\left(\Pi_{r, h} v\right)\right|_{0, e}\right)
\end{aligned}
$$

for $0 \leqq q \leqq m$. Furthermore,

$$
\begin{aligned}
&\left|\frac{\partial^{r+1}}{\partial s^{r+1}}\left(\Pi_{r, h} v\right)\right|_{0, e} \leqq C \delta^{1 / 2} \sup _{e}\left|\frac{\partial^{r+1}}{\partial s^{r+1}}\left(\Pi_{r, h} v\right)\right| \\
& \leqq C \delta^{1 / 2}\left\{\Pi_{r, h} v\right\}_{r+1, \tau}=C \delta^{1 / 2}\left\{\Pi_{r, h} v\right\}_{r, \tau} \\
& \leqq C \delta^{-1 / 2}\left\|\Pi_{r, h} v\right\|_{r, \tau} \leqq C \delta^{m-r-(3 / 2)}\left\|\Pi_{r, h} v\right\|_{m-1, \tau}
\end{aligned}
$$

for $1 \leqq m \leqq r+1$, where we have used the fact that $\Pi_{r, h} v$ is a polynomial of degree $r$ on $\tau$ and inequalities (2) and (3). Thus

$$
\begin{aligned}
\left|v-\Pi_{r, h} v\right|_{q, e} \leqq C \delta^{m-q}\left(|v|_{m, e}+\delta^{-1 / 2}\right. & \left.\left\|\Pi_{r, h} v\right\|_{m-1, \tau}\right) \\
& \leqq C h^{m-q}\left(|v|_{m, e}+h^{-1 / 2}\left\|\Pi_{r, h} v\right\|_{m-1, \tau}\right)
\end{aligned}
$$

Part (a) of the theorem now follows from lemma 2. //
Up to this point we have made no essential use of the macroelements of degree $r \geqq 5$ introduced in section 1: the results proved so far hold for the standard $C^{1}$ element $\left(\tau, P_{r}(\tau), \bar{\Sigma}_{r}\right), r \geqq 5$, of remark 3 , section 1 . Now, however, we demonstrate that the use of the macroelements of section 1 along $\partial \Omega$ results in a reduced interpolation error for boundary values and normal derivatives when measured in negative Sobolev norms. We deal first with the case of nonnodal degrees of freedom (see remark 1 , section 1) on $\partial \Omega$ because the proof is easier and the result is better in the sense that less smoothness is required of the function being interpolated. For $\Gamma$ a smooth $\operatorname{arc}$ in $\mathbf{R}^{2}$, let $\widetilde{P}_{m}(\Gamma)$ denote the space of functions on $\Gamma$ which are polynomials in arclength of degree not greater than $m$. Adopt the convention that $\widetilde{P}_{-1}(\Gamma)$ is the set consisting of the null function.

Theorem 4: Suppose that each element in $\mathscr{I}_{h}$ with an edge $e$ in $\partial \Omega$ is such that for $v \in \mathrm{H}^{m_{0}}(\Omega)$

$$
\begin{equation*}
\left\langle v-\Pi_{r, h} v, \psi\right\rangle_{e}=0 \quad \text { when } \psi \in \tilde{P}_{r-4}(e) \tag{a}
\end{equation*}
$$

and/or

$$
\begin{equation*}
\left\langle\frac{\partial}{\partial n}\left(v-\Pi_{r, h} v\right), \psi\right\rangle \in=0 \quad \text { when } \psi \in \tilde{P}_{r-3}(e) \tag{b}
\end{equation*}
$$

Suppose that $m_{0} \leqq m \leqq r+1$. Then there exists a $C$ such that
(a) for $v \in \mathrm{H}^{m}(\Omega) \cap H^{m}(\partial \Omega)$,
and/or $\left|v-\Pi_{r, h} v\right|_{-p, \partial \Omega} \leqq C h^{m+p}\left(|v|_{m, \partial \Omega}+\|v\|_{m, \Omega}\right), \quad 0 \leqq p \leqq r-3$,
(b) for $v \in H^{m}(\Omega)$ such that $\partial v / \partial n \in H^{m-1}(\partial \Omega)$,

$$
\left|\frac{\partial}{\partial n}\left(v-\Pi_{r, h} v\right)\right|_{-p, \partial \Omega} \leqq C h^{m+p-1}\left(|v|_{m-1, \partial \Omega}+\|v\|_{m, \Omega}\right), \quad 0 \leqq p \leqq r-2
$$

Proof: Let $\varphi \in H^{p}(\partial \Omega)$ and let $\psi \in L^{2}(\partial \Omega)$ be such that $\left.\psi\right|_{e} \in \widetilde{P}_{\bar{p}}(e)$ for all boundary edges $e$ in $\mathscr{I}_{h}$ and (cf. lemma 1)

$$
|\varphi-\psi|_{0, \partial \Omega} \leqq C h^{q}|\varphi|_{q, \partial \Omega}, \quad 0 \leqq q \leqq p
$$

Thus for $m_{0} \leqq m \leqq r+1$ and $v \in H^{m}(\Omega) \cap H^{m}(\partial \Omega)$,

$$
\begin{aligned}
\left\langle v-\Pi_{r, h} v, \varphi\right. & \rangle_{\partial \Omega}
\end{aligned} \quad\left\langle v-\Pi_{r, h} v, \varphi-\psi\right\rangle_{\partial \Omega} .
$$

proving part (a). Part $(b)$ is similar. //
Note that in part (a) the upper limit on $p$ is just the number of value degrees of freedom associated with the interior of an edge lying on $\partial \Omega$ and in part (b) the upper limit on $p$ is the number of normal derivative degrees of freedom associated with the interior of an edge lying on $\partial \Omega$. It is clear that if the standard $C^{1}$ element $\left(\tau, P_{r}(\tau), \bar{\Sigma}_{r}\right)$ were used along $\partial \Omega$ (after modification as in remark 1), then a result similar to theorem 4 would hold except that the upper limit on $p$ would be reduced by two in both parts of the theorem.

When nodal degrees of freedom are used along $\partial \Omega$, the orthogonality conditions in theorem 4 can no longer be satisfied exactly. Nevertheless (see [30]), it is possible to place the nodes on $\partial \Omega$ in such a way that the integrals involved in those orthogonality conditions are sufficiently small for the rates of convergence of theorem 4 to be retained, as will now be shown.

Let $e$ be an edge of $\mathscr{I}_{h}$ on $\partial \Omega$, let $l(e)$ be the length of $e$ and suppose $e$ is parametrized by $s \in[0, l(e)]$. Let

$$
A_{r, e}(s)=\frac{d^{r-3}}{d s^{r-3}}\left(s^{r-1}(l(e)-s)^{r-1}\right)
$$

and let

$$
B_{r, e}(s)=\frac{d^{r-2}}{d s^{r-2}}\left(s^{r-1}(l(e)-s)^{r-1}\right)
$$

Integration by parts shows that

$$
\left\langle A_{r, e}, \psi\right\rangle_{e}=0, \quad \psi \in \widetilde{P}_{r-4}(e)
$$

and

$$
\left\langle B_{r, e}, \psi\right\rangle_{e}=0, \quad \psi \in \widetilde{P}_{r-3}(e)
$$

It follows that $A_{r, e}$, whose degree is $r+1$, has $r-3$ distinct zeroes in the interior of $[0, l(e)]$ in addition to the second order zeroes at the endpoints and that $B_{r, e}$, whose degree is $r$, has $r-2$ distinct zeroes in the interior of $[0, l(e)]$ in addition to the zeroes at the endpoints. The zeroes of $B_{r, e}$ are the Lobatto quadrature points for $e$. Quadrature rules based on both $A_{r, e}$ and $B_{r, e}$ were extensively used in [14].

Theorem 5: (a) Suppose that each element in $\mathscr{I}_{h}$ with an edge $e$ on $\partial \Omega$ is a macroelement with nodes on e placed so that $v-\Pi_{r, h} v$ has the same zeroes on $e$ (each with the same order) as $A_{r, e}$. If $m_{0} \leqq m \leqq r+1,0 \leqq p \leqq r-3$, and $v \in H^{m}(\Omega) \cap H^{m+p}(\partial \Omega)$, then

$$
\left|v-\Pi_{r, h} v\right|_{-p, \partial \Omega} \leqq C h^{m+p}\left(|v|_{m+p, \partial \Omega}+\|v\|_{m, \Omega}\right) .
$$

(b) Suppose that each element in $\mathscr{I}_{h}$ with an edge e on $\partial \Omega$ is a macroelement with nodes on e placed so that $(\partial / \partial n)\left(v-\Pi_{r, h} v\right)$ has the same zeroes on e as $B_{r, e}$. If $m_{0} \leqq m \leqq r+1,0 \leqq p \leqq r-2, v \in H^{m}(\Omega)$, and $\partial v / \partial n \in H^{m+p-1}(\partial \Omega)$, then

$$
\left|\frac{\partial v}{\partial n}-\frac{\partial}{\partial n}\left(\Pi_{r, h} v\right)\right|_{-p, \partial \Omega} \leqq C h^{m+p-1}\left(\left|\frac{\partial v}{\partial n}\right|_{m+p-1, \partial \Omega}+\|v\|_{m, \Omega}\right)
$$

Proof: We give the proof of part (a); part (b) is similar.
Let $T, \tau, e, \delta$ and $l(e)$ be as in the proof of theorem 3. Let $\varphi \in H^{p}(\partial \Omega)$. Choose $\psi$ as in theorem 4 so that

$$
h^{-q}|\varphi-\psi|_{0, e}+|\psi|_{q, e} \leqq C|\varphi|_{q, e} \quad 0 \leqq q \leqq p
$$

Note that $\rho \delta \leqq l(e)$, where $\rho$ is the constant associated with the non-degeneracy of the family of triangulations. Since

$$
\left\langle v-\Pi_{r, h} v, \varphi\right\rangle_{\partial \Omega}=\left\langle v-\Pi_{r, h} v, \varphi-\psi\right\rangle_{\partial \Omega}+\left\langle v-\Pi_{r, h} v, \psi\right\rangle_{\partial \Omega}
$$

and

$$
\left\langle v-\Pi_{r, h} v, \varphi-\psi\right\rangle_{\partial \Omega} \leqq C h^{m+p}\left(|v|_{m . \partial \Omega}+\|v\|_{m, \Omega}\right)|\varphi|_{p, \partial \Omega}
$$

as in the proof of theorem 4, we just need to estimate $\left\langle v-\Pi_{r, h} v, \psi\right\rangle_{\partial \Omega}$.

Consider the linear functional $E$ on $C^{1}([0, l(e)])$ given by

$$
E(f)=\int_{0}^{l(e)}\left(f(s)-f_{I}(s)\right) d s
$$

where $f_{1}$ is the unique polynomial of degree $r$ such that the set of zeroes $f-f_{I}$ include the zeroes (counting multiplicity) of $A_{r, e}$. Note that $E(f)=0$ if $f^{I}$ is a polynomial of degree $2 r-3$ since $f-f_{I}$ factors into $A_{r, e}$ times a polynomial of degree $r-4$. Thus, by the Peano Kernel theorem, there exists $C$ such that for $f \in C^{k}([0, l(e)])$

$$
|E(f)| \leqq C l(e)^{k}\left\|\left(\frac{d}{d s}\right)^{k} f\right\|_{L^{1}((0, l(e)])}, \quad 2 \leqq k \leqq 2 r-2 .
$$

Let $\psi \in \widetilde{P}_{\bar{p}}(e)$. Then by Schwarz's inequality

$$
|E(f \psi)| \leqq C l(e)^{k}|f|_{k, e}|\psi|_{p, e}
$$

for $k$ an integer in the range $2 \leqq k \leqq 2 r-2$, with $C$ independent of $\psi$. Viewing $f \rightarrow E(f \psi)$ as a linear functional, it follows by interpolation [22] that the above holds for all real $k$ in that range.

Let $\theta \in \widetilde{P}_{r}(e)$ be such that $\theta-v$ (as a function of arc length) vanishes at the zeroes of $A_{r, e}$ (counting multiplicity). Then

$$
\begin{aligned}
& \left\langle v-\Pi_{r, h} v, \psi\right\rangle_{e}=\langle v-\theta, \psi\rangle_{e}+\left\langle\theta-\Pi_{r, h} v, \psi\right\rangle_{e} \\
& \quad \leqq C\left(\delta^{m+p}|v-\theta|_{m+p, e}+\delta^{2 r-2}\left|\theta-\Pi_{r, h} v\right|_{2 r-2, e}\right)|\psi|_{p, e}
\end{aligned}
$$

Using an argument similar to that in the proof of theorem 3, we see that

$$
\left|\Pi_{r, h} v\right|_{2 r-2, e} \leqq C \delta^{m-r-3 / 2}\left\|\Pi_{r, h} v\right\|_{m-1, \tau}
$$

Since $\theta$ is a polynomial, we see that

$$
\delta^{2 r-2}|\theta|_{2 r-2, e}=\delta^{2 r-2}|\theta|_{r, e} \leqq C \delta^{r+m-3}|\theta|_{m-1, e} \leqq C \delta^{r+m-3}|v|_{m+p, e}
$$

Thus from the above estimates and the fact that $p \leqq r-3$, we see that

$$
\left\langle v-\Pi_{r, h} v, \psi\right\rangle_{e} \leqq C\left(\delta^{m+p}|v|_{m+p, e}+\delta^{m+r-7 / 2}\left\|\Pi_{r, h} v\right\|_{m-1, r}\right)|\psi|_{p, e}
$$

Summing over $e$ and applying lemma 2, we get

$$
\left\langle v-\Pi_{r, h} v, \psi\right\rangle_{\partial \Omega} \leqq C h^{m+p}\left(|v|_{m+p, \partial \Omega}+\|v\|_{m, \Omega}\right)|\varphi|_{p, \partial \Omega,}
$$

which finishes the proof. //
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The last result in this section is that we can extract from the proofs of theorems 3, 4 and 5 the fact that $S_{r, h}$ contains subspaces consisting of functions with "nearly zero" boundary values and normal derivatives (see [24] and [30]).

Thforem 6: Let $S_{r, h}^{0}$ be any subspace of $S_{r . h}$ such that if $\chi \in S_{r, h}^{0}$ and $e$ is an edge of.$Y_{h}$ on $i \Omega$, then
(a) $\chi=0$ and $\partial \chi / \partial s=0$ at the endpoints of $e$, and either $\langle\chi, \psi\rangle_{e}=0$ for $\psi \in \widetilde{P}_{r-4}(e)$ or $\chi$ vanishes at the zeroes of $A_{r, e}$ in the interior of $e$, and/or
(b) $\partial \chi / \partial n=0$ at the endpoints of $e$ and either $\langle\partial \chi / \partial n, \psi\rangle_{e}=0$ for $\psi \in \widetilde{P}_{r-3}(e)$ or $\partial \chi / \partial n$ vanishes at the zeroes of $B_{r, e}$ in the interior of $e$,

Then, for $\chi \in S_{r, h}^{0}$ and $0 \leqq m \leqq 2$,
(a)

$$
|\chi|_{-p, \partial \Omega} \leqq C h^{m+p+(1 / 2)}\|\chi\|_{m, \Omega}
$$

when $-2 \leqq p \leqq r-3$ and $m+p+1 / 2 \geqq 0$, and/or
(b)

$$
\left|\frac{\partial \chi}{\partial n}\right|_{-p, \AA \Omega} \leqq C h^{m+p-(1 / 2)}\|\chi\|_{m, \Omega}
$$

when $-1 \leqq p \leqq r-2$ and $m+p-1 / 2 \geqq 0$
Proof: As usual, we only prove part (a).
Let $e$ be an edge of $\mathscr{I}_{h}$ on $\partial \Omega$, let $l(e)$ be the length of $e$, let $\tau$ be the simple triangle associated with $\mathscr{I}_{h}$ which has $e$ as an edge and let $\delta$ be the diameter of $\tau$. Since the family of triangulations is non-degenerate,

$$
C^{-1} \delta \leqq l(e) \leqq C \delta
$$

No matter which version of $(a)$ is satisfied on $e$, if $\chi \in S_{r, h}^{0}$, then $\chi \mid e$ has $r+1$ zeros, counting multiplicity, so by the Poincaré inequality,

$$
|\chi|_{q, e} \leqq C \delta^{r-q+1}|\chi|_{r+1, e}
$$

for $0 \leqq q \leqq r+1$. As in the proof of theorem 3, the inequalities (2) and (3) imply that, for $0 \leqq m \leqq r \leqq j$,

$$
\begin{equation*}
|\chi|_{j, e} \leqq C \delta^{m-r-(1 / 2)}\|\chi\|_{m, \tau} \tag{4}
\end{equation*}
$$

Thus, for $0 \leqq q \leqq r+1$ and $0 \leqq m \leqq r$,

$$
|\chi|_{q, e} \leqq C \delta^{m-q+(1 / 2)}\|\chi\|_{m, \tau} .
$$

Summing over all $e$ and using the fact that $\delta \leqq h$, we obtain

$$
|\chi|_{q, \partial \Omega} \leqq C h^{m-q+(1 / 2)}\|\chi\|_{m, \Omega},
$$

for $0 \leqq q \leqq 2,0 \leqq m \leqq 2$, provided that $m-q+(1 / 2) \geqq 0$. Letting $q=-p$, we see that part (a) of the theorem is now proved for $-2 \leqq p \leqq 0$.

To prove part ( $a$ ) in the range $0 \leqq p \leqq r-3$, let $\varphi$ and $\psi$ be as in the proof of theorem 5. Then

$$
\langle\chi, \varphi-\psi\rangle_{e} \leqq|\chi|_{0, e}|\varphi-\psi|_{0, e} \leqq C \delta^{m+p+(1 / 2)}\|\chi\|_{m, \tau}|\varphi|_{p, e}
$$

for $0 \leqq m \leqq r$. Furthermore, either $\langle\chi, \psi\rangle_{e}$ is automatically zero by the definition of $S_{r, h}^{0}$ in terms of orthogonalities or $\chi \psi$ vanishes at the zeros of $A_{r, e}$. Therefore, as in the proof of theorem 5,

$$
\langle\chi, \psi\rangle_{e} \leqq C \delta^{2 r-2}|\chi|_{2 r-2, e}|\psi|_{p, e}
$$

Using (4) and the definition of $\psi$,

$$
\langle\chi, \psi\rangle_{e} \leqq C \delta^{r+m-(5 / 2)}\|\chi\|_{m, \tau}|\psi|_{p, e}
$$

For $0 \leqq m \leqq r$ and $0 \leqq p \leqq r-3$, we therefore have

$$
\langle\chi, \varphi\rangle_{e}=\langle\chi, \varphi-\psi\rangle_{e}+\langle\chi, \psi\rangle_{e} \leqq C \delta^{m+p+(1 / 2)}\|\chi\|_{m, \tau}|\varphi|_{p, e}
$$

Thus, for $0 \leqq m \leqq 2,0 \leqq p \leqq r-3$, and $m+p+1 / 2 \geqq 0$,

$$
\langle\chi, \varphi\rangle_{\partial \Omega} \leqq C h^{m+p+(1 / 2)}\|\chi\|_{m, \Omega}|\varphi|_{p, \partial \Omega},
$$

which completes the proof. //
Remark 5: Since the techniques of proof in Theorems 2-6 are purely local, it is not necessary to segregate the methods of orthogonality and interpolation. Indeed, when the boundary data is singular, it would be wise to use orthogonality to impose the boundary conditions near the singularity, while interpolation could be used away from the singularity (cf. the different smoothness requirements on the data in theorem 4 versus theorem 5). Also, it is not necessary to require that $\partial \Omega$ be smooth globally; if $\partial \Omega$ is piecewise smooth and if a boundary vertex is placed at every point of $\partial \Omega$ where it is not $C^{\infty}$, the same results follow. //

Remark 6: It may be desirable to impose the orthogonalities in theorem 4 by evaluating the integral over $e$ using a numerical quadrature rule:

$$
\int_{e} f \sim \sum_{i=1}^{k} w_{i} f\left(z_{i}\right)
$$

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In the case that the weights $w_{i}$ are all nonzero and $\left\{z_{i}\right\}$ corresponds to the set of zeroes of $A_{r, e}\left(\right.$ resp. $\left.B_{r, e}\right)$ in the interior of $[0, l(e)]$ then orthogonality of $v-\Pi_{r, h} v$ to $\Psi \in \widetilde{P}_{r-4}(e)$ [resp. $(\partial / \partial n)\left(v-\Pi_{r, h} v\right)$ to $\left.\Psi \in \widetilde{P}_{r-3}(e)\right]$ with respect to the quadrature on $e$ is equivalent to vanishing at the quadrature points, i.e., the situation covered by theorem 5. The error induced by using other quadrature rules may be studied using the techniques in the proof of theorem 5. /

## 3. APPLICATION TO THE PLATE BENDING PROBLEM

Consider the bilinear form $a(.,$.$) on H^{2}(\Omega)$ defined by

$$
a(u, v)=\frac{\mathrm{D}}{2} \int_{\Omega}\left\{\Delta u \Delta v-(1-v)\left(u_{x x} v_{y y}+u_{y y} v_{x x}-2 u_{x y} v_{x y}\right)\right\} d x d y
$$

where D and $v$ are constants such that $D>0$ and $0 \leqq v \leqq 1 / 2$.
Let $V$ be the subspace of $H^{2}(\Omega)$ consisting of functions that vanish on $\partial \Omega$. Since $a(v, v) \geqq(1 / 2) D(1-v)[v]_{2, \Omega}^{2}$, it follows from Rellich's lemma [1, chapt. 10] that for some $\gamma<\infty$.

$$
\begin{equation*}
a(v, v) \geqq \frac{1}{\gamma}\|v\|_{2, \Omega}^{2} \quad \text { for all } v \in V \tag{5}
\end{equation*}
$$

Given $F \in H^{-2}(\Omega)$ and $g \in H^{2}(\Omega)$, there is a unique $u \in H^{2}(\Omega)$ such that

$$
u-g \in V(\text { i. e., } u=g \text { on } \partial \Omega)
$$

and

$$
a(u, v)=F(v) \quad \text { for all } v \in V
$$

(the Lax-Milgram theorem). Suppose that the inner product with $F$ is given by

$$
F(v)=\int_{\Omega} f v d x+\int_{\partial \Omega} M \frac{\partial v}{\partial n} d s
$$

where $f \in L^{2}(\Omega)$ and $M \in L^{2}(\partial \Omega)$. Then $u$ is the solution to the simply supported plate bending problem corresponding to a loading $f$, an edge displacement $g$, and a moment $M$ applied to the edge. The constant $v$ in the definition of $a(.,$. is Poisson's ratio and $D$ is the flexural rigidity [3]. When $f \in H^{-4}(\Omega)$, $g \in H^{s-(1 / 2)}(\partial \Omega)$, and $M \in H^{s-(5 / 2)}(\partial \Omega)$, then $u \in H^{s}(\Omega)$, with the obvious norm inequality ( $s \geqq 2$ ), since $u$ is related to ( $f, g, M$ ) by a properly elliptic boundary value problem [22], [1]. We now consider a Galerkin approximation to $u$.

Let $\Pi_{r, h}, 0<h \leqq h_{0}$, be a family of interpolation operators as studied in the previous section, and let $S_{r, h}^{0}$ be the image of $V \cap H^{m_{0}}(\Omega)$ via the mapping $\Pi_{r, h}$.

Suppose that

$$
\begin{equation*}
a(\chi, \chi) \geqq \frac{1}{2 \gamma}\|\chi\|_{2, \Omega}^{2} \text { for all } \chi \in S_{r, h}^{0} . \tag{6}
\end{equation*}
$$

(As can be seen from [30], (6) follows from (5) for $h_{0}$ sufficiently small.) Then there is a uniquely determined $u^{h} \in S_{r, h}$ such that

$$
u^{h}-\Pi_{r, h} g \in S_{r, h}^{0}\left(\text { i. e., } u^{h} \text { interpolates } g \text { on } \partial \Omega\right)
$$

and

$$
a\left(u^{h}, \chi\right)=F(\chi) \quad \text { for all } \chi \in S_{r, h}^{0} .
$$

Note that $u$ and $u^{h}$ depend only on the values of $g$ on $\partial \Omega$.
The Galerkin methods above are direct generalizations to a fourth order elliptic boundary value problem of the methods studied by Blair [4] and in the papers [2] and [30].

Theorem 7: Let $\Pi_{r, h}$ be as in either theorem 4 or 5 . Suppose that (6) holds for $0<h \leqq h_{0}$ and that $u \in H^{m}(\Omega)$ for $m$ in the range $7 / 2<m_{0} \leqq m \leqq r+1$. Then

$$
\left\|u-u^{h}\right\|_{2, \Omega} \leqq c h^{m-2}\|u\|_{m, \Omega} .
$$

Suppose further that $g \in H^{m}(\partial \Omega)$ and that $\Pi_{r, h}$ is as in theorem 4. Then

$$
\left\|u-u^{h}\right\|_{s, \Omega} \leqq c h^{m-s}\left(\|u\|_{m, \Omega}+|g|_{m, \partial \Omega}\right)
$$

for $3-r \leqq s \leqq 2$. When $\Pi_{r, h}$ is as in theorem 5 , suppose that $g \in H^{q+m}(\partial \Omega)$, $0 \leqq q \leqq r-3$. Then

$$
\left\|u-u^{h}\right\|_{s, \Omega} \leqq c h^{m-s}\left(\|u\|_{m, \Omega}+|g|_{m+q, \partial \Omega}\right)
$$

for $-q \leqq s \leqq 2$.
Proof: As in [30, section A. 1] (6), Green's theorem and the trace theorem [22] imply that, for $m>7 / 2$,

$$
\left\|u-u^{h}\right\|_{2 . \Omega} \leqq C_{1}\left\|u-\Pi_{r, h} u\right\|_{2, \Omega}+c_{2}\|u\|_{m, \Omega} \sup _{\chi \in S_{r, h}^{0}} \frac{|\chi|_{7 / 2-m, \partial \Omega}}{\|\chi\|_{2, \Omega}} .
$$

Thus theorems 2 and 6 yield the first conclusion. Now let $\varphi \in H^{k}(\Omega)$, $0 \leqq k \leqq r-3$, and let $\Phi \in V$ solve $\left.\left({ }^{*}\right), a(\Phi), v\right)=(\varphi, v)$ for all $v \in V$. By elliptic regularity theory (see above), $\|\Phi\|_{k+4, \Omega} \leqq c\|\varphi\|_{k, \Omega}$. Integration by parts (Green's theorem) and the trace theorem yield

$$
\begin{aligned}
\left(u-u^{h}, \varphi\right) \leqq \mid a\left(u-u^{h},\right. & \left.\Phi-\Pi_{r, h} \Phi\right) \mid \\
& +c\|\Phi\|_{k+4, \Omega}\left|g-u^{h}\right|_{-(1 / 2)-k, \partial \Omega}+c\|u\|_{m, \Omega}\left|\Pi_{r, h} \Phi\right|_{7 / 2-m, \partial \Omega} .
\end{aligned}
$$

(*) Hence forth (.. .) and $\langle., \text {.〉 will denote (., .) })_{\Omega}$ and $\langle., .\rangle_{\partial \Omega}$.
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Using the first part of this theorem, theorem 2, and either theorem 4 or 5,

$$
\left(u-u^{h}, \varphi\right) \leqq c\left(h^{m+k}\|u\|_{m, \Omega}+\left|g-u^{h}\right|_{-(1 / 2)-k, \partial \Omega}\right)\|\Phi\|_{k+4, \Omega}
$$

(recall that $\Phi \in V$, i.e., $\Phi \equiv 0$ on $\partial \Omega$, so that the restrictions of theorems 4 and 5 are the same). Thus regardless of the type of interpolation,

$$
\left\|u-u^{h}\right\|_{-k, \Omega} \leqq c\left(h^{m+k}\|u\|_{m, \Omega}+\left|g-u^{h}\right|_{-(1 / 2)-k, \partial \Omega}\right)
$$

Now consider $u^{h}$. Theorem 6 implies that

$$
\begin{aligned}
&\left|g-u^{h}\right|_{-(1 / 2)-k, \partial \Omega} \leqq\left|u^{h}-\Pi_{r, h} u\right|_{-(1 / 2)-k, \partial \Omega}+\left|g-\Pi_{r, h} u\right|_{-(1 / 2)-k, \partial \Omega} \\
& \leqq c h^{k+2}\left\|u^{h}-\Pi_{r, h} u\right\|_{2, \Omega}+\left|u-\Pi_{r, h} u\right|_{-(1 / 2)-k, \partial \Omega} .
\end{aligned}
$$

Inserting $(u-u)$ in the first term and using the triangle inequality, the first part of this theorem and theorem 2 imply that

$$
\left\|u-u^{h}\right\|_{-k, \Omega} \leqq c\left(h^{m+k}\|u\|_{m, \Omega}+\left|u-\Pi_{r, h} u\right|_{-(1 / 2)-k, \partial \Omega}\right)
$$

Apply either theorem 4 or 5 to estimate the second term on the right hand side, and the result follows. //

The above theorem can clearly be extended to allow more general coefficients in $a(.,$.$) . The clamped plate problem may be treated by simply changing V$ in the above to be the subspace of $H^{2}(\Omega)$ consisting of functions that vanish to second order on $\partial \Omega$ and using part ( $b$ ) of theorems 4-6.

## 4. THE $H^{1}$-GALERKIN METHOD FOR SECOND ORDER ELLIPTIC PROBLEMS

Consider the Dirichlet problem

$$
\begin{gather*}
L u=\nabla \cdot(a \nabla u)=f \text { in } \Omega,  \tag{7}\\
u=g \quad \text { on } \partial \Omega,
\end{gather*}
$$

where $a=a(x)$ is a smooth, positive function on $\bar{\Omega}$. Again let $V=\left\{v \in H^{2}(\Omega): v=0\right.$ on $\left.\partial \Omega\right\}$, and let $S_{r, h}^{0}$ be the image of $V \cap H^{m_{0}}(\Omega)$ under the mapping $\Pi_{r, h}$, where $\Pi_{r, h}$ is determined by either the interpolation procedure associated with theorem 5 or the orthogonality conditions associated with theorem 4. Assume that $g$ has a (theoretical) extension to $H^{m_{0}}(\Omega)$, so that $\Pi_{r, h} g$ is definable; practically, this involves only the values of $g$ on $\partial \Omega$. Then the $H^{1}$-Galerkin method for approximating the solution of (7) consists of finding $u^{h} \in S_{r, h}$ such that

$$
\begin{gather*}
\left(L u^{h}, \Delta v\right)=(f, \Delta v), \quad v \in S_{r, h}^{0}  \tag{8}\\
u^{h}-\Pi_{r, h} g \in S_{r, h}^{0}
\end{gather*}
$$

$H^{1}$-Galerkin methods have been proposed earlier in [16] for both elliptic and parabolic problems for the special case $L=\Delta$, and the nonlinear Dirichlet problem based on (7) with $a=a(x, u)$, plus a linear problem with lower order terms included, have been treated in [13]. In both these papers the Dirichlet boundary condition was imposed weakly through penalty-like terms on the boundary. The method (8) corresponding to interpolation was mentioned briefly in [13], but no analysis was given in that case. The orthogonality method (8) presents an $H^{1}$ analogue of the method of Blair [4].

Note that the algebraic equations arising from (8) do not, in general, generate a symmetric matrix, and care must be taken to show that a solution of (8) exists. There are significant practical advantages of (8) over least squares methods [6,7], particularly for nonlinear problems and in applications to transient problems, since the algebraic equations become simpler. The analysis of ( 8 ) below is similar to that given in [13], but the details are noticeably different. Both rely on ideas discussed by Schatz [28] earlier.

Lemma 3 (Gårding inequality): There exist constants $h_{0}>0, \rho>0$, and $C$ such that

$$
(L v, \Delta v) \geqq \rho\|v\|_{2, \Omega}^{2}-C\|v\|_{0, \Omega}^{2}
$$

for $0<h \leqq h_{0}$ and $v \in S_{r, h}^{0}$.
Proof: Since $a$ is bounded below positively and $\nabla a$ is bounded, it is trivial to see that

$$
(L v, \Delta v) \geqq \rho_{1}\|\Delta v\|_{0, \Omega}^{2}-C_{1}\|v\|_{1, \Omega}^{2}, \quad v \in H^{2}(\Omega)
$$

For $v \in S_{r, h}^{0}$, theorem 6 implies that $|v|_{3 / 2, \tilde{\Omega} \Omega} \leqq C h\|v\|_{2, \Omega}$. Since $\|v\|_{2, \Omega}$ and $\|\Delta v\|_{0, \Omega}+|v|_{3 / 2, \partial \Omega}$ are equivalent [22], a simple version [1] of interpolation of Sobolev norms implies that

$$
(L v, \Delta v) \geqq \rho_{2}\|v\|_{2, \Omega}^{2}-C_{1}\left(\varepsilon\|v\|_{2, \Omega}^{2}+\varepsilon^{-1}\|v\|_{0, \Omega}^{2}\right) \geqq \rho\|v\|_{2, \Omega}^{2}-C\|v\|_{0, \Omega}^{2}
$$

for $v \in S_{r, h}^{0}$ and $h$ sufficiently small.
Lemma 4: If $\zeta=u-u^{h}$ and $s \in[-2, r-3]$, then

$$
\|\zeta\|_{-s, \Omega} \leqq C\left\{h^{s+2}\|\zeta\|_{2, \Omega}+|\zeta|_{-s-(1 / 2), \partial \Omega}\right\} .
$$

Proof: Let $w \in H^{s}(\Omega)$ and determine $\chi \in H^{s+2}(\Omega)$ and $\psi \in H^{s+(1 / 2)}(\partial \Omega)$ such that

$$
(v, w)=(L v, \chi)+\langle v, \psi\rangle
$$

for all $v \in H^{2}(\Omega)$. The existence of $\chi$ and $\psi$ follows from lemma 6.1 of DouglasDupont [15]; also the following bound holds:

$$
\|\chi\|_{s+2, \Omega}+|\psi|_{s+(1 / 2), \Delta \Omega} \leqq C\|w\|_{s, \Omega}
$$

$\left(\psi=\left.a(\partial \chi / \partial n)\right|_{\partial \Omega}\right.$ whenever the latter makes sense). Let $\varphi$ solve Poisson's equations:

$$
\begin{aligned}
& \Delta \varphi=\chi \quad \text { in } \Omega \\
& \varphi=0 \quad \text { on } \partial \Omega .
\end{aligned}
$$

Then $\varphi \in H^{s+4}(\Omega)$ and

$$
(\zeta, w)=(L \zeta, \Delta \varphi)+\langle\zeta, \psi\rangle
$$

Since $(L \zeta, \Delta v)=0$ for $v \in S_{r, h}^{0}$, an appropriate choice of $v \in S_{r, h}^{0}$ yields $(\zeta, w)=(L \zeta, \Delta(\varphi-v))+\langle\zeta, \psi\rangle$

$$
\begin{aligned}
\leqq C\|\zeta\|_{2, \Omega} h^{s+2}\|\varphi\|_{s+4, \Omega} & +|\zeta|_{-s-(1 / 2), \Delta \Omega}|\psi|_{s+(1 / 2), \partial \Omega} \\
& \leqq C\left(\|\zeta\|_{2, \Omega} h^{s+2}+|\zeta|_{-s-(1 / 2), \Delta \Omega}\right)\|w\|_{s, \Omega}
\end{aligned}
$$

and the proof is finished. //
Lemma 5: Let $\alpha=0$ if orthogonality determines $\Pi_{r, h}$, and let $\alpha=s+(1 / 2)$ if interpolation determines $\Pi_{r, h}$. If $-1 / 2 \leqq s \leqq r-(7 / 2)$, then

$$
|\zeta|_{-s-(1 / 2), \partial \Omega} \leqq C h^{s+3}\|\zeta\|_{2, \Omega}+C h^{m+s+(1 / 2)}\left(\|u\|_{m, \Omega}+|g|_{m+\alpha, \partial \Omega}\right)
$$

for $m_{0} \leqq m \leqq r+1$.
Proof: Let $-1 / 2 \leqq s \leqq r-(7 / 2)$. Then

$$
|\zeta|_{-s-(1 / 2), \partial \Omega} \leqq\left|u-\Pi_{r, h} u\right|_{-s-(1 / 2), \partial \Omega}+\left|\Pi_{r, h} u-u^{h}\right|_{-s-(1 / 2), \partial \Omega}
$$

Apply either theorem 4 or 5 to the first term on the right hand side and theorem 6 to the second. After a trivial simplification, the desired inequality results. //

Logically it would perhaps have been better to show existence and uniqueness of a solution of (8) before treating lemmas 4 and 5 , but it would have induced an unnecessary duplication of argument.

Lemma 6: For h sufficiently small there exists a unique solution $u^{h} \in S_{r, h}$ of (8).
Proof: It is clear that uniqueness implies existence and that the difference $z$ of two solutions of (8) is an element of $S_{r, h}^{0}$ satisfying $(L z, \Delta v)=0$ for $v \in S_{r, h}^{0}$. Thus, $z$ corresponds to $\zeta$ in the case $u \equiv 0$, and lemma 3,4 , and 5 with $s=0$ imply that

$$
\|z\|_{2, \Omega} \leqq C h^{2}\|z\|_{2, \Omega} .
$$

Hence $z=0$ for small $h . / /$

Theorem 8: Let $\zeta=u-u^{h}$, where $u$ is the solution of (7) and $u^{h}$ is the unique solution of (8) for small $h$. Let $m_{0} \leqq k \leqq r+1$ and

$$
m=m(p)=\max \left\{m_{0}, k-p-(1 / 2)\right\}
$$

where $m_{0}$ appears in condition (c) of section 2. Then, if $r \geqq 4$,

$$
\|\zeta\|_{p, \Omega} \leqq C\left(h^{k-p}\|u\|_{k, \Omega}+h^{m+(1 / 2)}|g|_{m+\alpha, \partial \Omega}\right), \quad 0 \leqq p \leqq 2
$$

where $\alpha=0$ for $\Pi_{r, h}$ determined by orthogonality and $\alpha=1 / 2$ for $\Pi_{r, h}$ determined by interpolation. If $r=3,|g|_{m+\alpha, a \Omega}$ should be replaced by $|g|_{m+x+(1 / 2) \text { i } \Omega}$ in the inequality. If $0 \leqq s \leqq r-(7 / 2)$, then

$$
\|\zeta\|_{-s, \Omega} \leqq C\left(h^{k+s}\|u\|_{k, \Omega}+h^{n+s+(1 / 2)}|g|_{n+\beta, \Omega \Omega}\right)
$$

where $n=\max \left\{m_{0}, k-(1 / 2)\right\}$ and $\beta=0$ or $s+(1 / 2)$ if $\Pi_{r, h}$ is determined by orthogonality or interpolation, respectively. If $r-(7 / 2) \leqq s \leqq r-3$, then

$$
\|\zeta\|_{-s, \Omega} \leqq C\left(h^{k+s}\|u\|_{k, \Omega}+h^{t+r-3}|g|_{t+\gamma, \partial \Omega}\right)
$$

where $t=\max \left(m_{0}, k+s-r+3\right)$ and $\gamma=0$ or $r-3$ for $\Pi_{r, h}$ as above.
Proof: We have seen that $\|\zeta\|_{0, \Omega} \leqq C h^{2}\|\zeta\|_{2, \Omega}+|\zeta|_{-1 / 2, \partial \Omega}$. Let

$$
\xi=\Pi_{r, h} u-u^{h} \in S_{r, h}^{0},
$$

and apply the Gårding inequality to $\xi$. Then

$$
\begin{aligned}
\rho\|\xi\|_{2, \Omega}^{2} \leqq(L \xi, \Delta \xi)+C & \|\xi\|_{0, \Omega}^{2} \\
& =\left(L\left(\Pi_{r, h} u-u\right), \Delta \xi\right)+C\|\xi\|_{0, \Omega}^{2} \\
& \leqq C\left\|\Pi_{r, h} u-u\right\|_{2, \Omega}\|\xi\|_{2, \Omega}+C\|\xi\|_{0, \Omega}^{2}
\end{aligned}
$$

and

$$
\|\xi\|_{2, \Omega} \leqq C\left(\|u\|_{k, \Omega} h^{k-2}+\|\xi\|_{0, \Omega}\right), \quad m_{0} \leqq k \leqq r+1
$$

Thus,

$$
\|\zeta\|_{2, \Omega} \leqq C\left(\|u\|_{k, \Omega} h^{k-2}+\|\zeta\|_{0, \Omega}\right) \leqq C\left(\|u\|_{k, \Omega} h^{k-2}+\|\zeta\|_{2, \Omega} h^{2}+|\zeta|_{-1 / 2, \partial \Omega}\right)
$$

For small $h$ and by lemma 5 , if $r \geqq 4$,

$$
\begin{aligned}
&\|\zeta\|_{2, \Omega} \leqq C\left(\|u\|_{k, \Omega} h^{k-2}+|\zeta|_{-1 / 2, \partial \Omega}\right) \\
& \leqq C\left\{\|u\|_{k, \Omega} h^{k-2}+h^{m+(1 / 2)}\left(\|u\|_{m, \Omega}+|g|_{m+\alpha, \partial \Omega}\right)\right\}
\end{aligned}
$$

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Choose $m=\max \left\{m_{0}, k-5 / 2\right\}$ as called for in the statement of theorem 8 to obtain the desired bound for $p=2$ and $r \geqq 4$. The remainder of the proof consists of a careful application of lemmas 4 and 5 . //

The two $H^{1}$-Galerkin methods for the Dirichlet problem can be used to motivate $H^{1}$-treatments of parabolic problems. See [16] for a simple case (however, with penalty-set boundary values) and [15] for a somewhat analogous development. See also [33] for another related concept.

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