Nonconforming Finite Elements for the Stokes Problem*

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Dedicated to Professor Eugene Isaacson on the occasion of his 70th birthday

Abstract. A new stability result is obtained for the approximation of the stationary Stokes problem by nonconforming piecewise cubic approximations to the velocities and a discontinuous piecewise quadratic approximation to the pressure. The basic result is that for most reasonable meshes, these elements form a stable pair without the addition of quartic bubble functions (which had previously been added to insure stability).

1. Introduction. In the finite element approximation of the velocity-pressure formulation of the stationary Stokes equations using triangular finite elements, several approaches appear in the literature. One might classify these into conforming schemes which use a continuous piecewise polynomial approximation to the velocity and a discontinuous piecewise polynomial approximation to the pressure, conforming schemes which use continuous piecewise polynomial approximations for both velocity and pressure, macroelement schemes in which the pressure elements (usually discontinuous) are defined on a coarser mesh than the velocity elements, and nonconforming schemes in which the velocities are only continuous at appropriate Gauss points on the triangle edges, and discontinuous pressures are used.

In an early paper on the finite element approximation of the Stokes problem by the first author and P.-A. Raviart [4], several combinations of conforming and nonconforming velocity elements and discontinuous pressure elements were analyzed. Specifically, in the case of conforming velocity elements, it was shown that piecewise constant pressures could be paired with piecewise quadratic velocities (an idea suggested by Fortin [5]) to give a (suboptimal) O(h) energy norm convergence rate. When discontinuous piecewise linear pressures are used, a corresponding velocity space of continuous piecewise quadratics augmented by cubic bubble functions improves the convergence rate to $O(h^2)$. In the case of discontinuous quadratic pressures, a corresponding velocity space of continuous piecewise cubics augmented by two quartic bubbles for each component of velocity gives an energy norm convergence rate of $O(h^3)$. In later work of Bernardi and Raugel [1], following another idea of Fortin [6], it was shown that pairing piecewise constant pressures with a velocity space consisting of continuous piecewise linear functions augmented by one (vector) quadratic per edge also yield an O(h) energy norm convergence rate. In the case of higher-order elements, Scott and Vogelius [10] proved that except for some exceptional meshes, the use of discontinuous piecewise polynomials of degree

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 $n\ (n\geq 3)$ for pressure and continuous piecewise polynomials of degree n+1 for velocity (without adding any bubble functions) give methods with optimal-order convergence rates.

Since the pairing of constant pressures with continuous piecewise linear velocities is not convergent (for most meshes, the set of such velocities with zero divergence is {0}) and the substitution of quadratic velocities leads to a suboptimal convergence rate, Crouzeix and Raviart also studied the use of nonconforming velocity spaces. They showed that constant pressures paired with nonconforming piecewise linear velocities give an optimal O(h) energy norm convergence rate. The case of nonconforming quadratic velocities was not considered in their paper, but was considered in a later paper of Fortin [7], where it was proved that nonconforming quadratic velocities may be paired with discontinuous linear pressures to produce an $O(h^2)$ method. Crouzeix and Raviart also considered the case of nonconforming cubic velocities and discontinuous quadratic pressures, but obtained an $O(h^3)$ convergence rate only by augmenting each component of the velocity space by two quartic bubble functions (as in the conforming case). The purpose of this paper is to improve this last result by showing that for most commonly used meshes, the quartic bubble functions are not needed to give a convergent method of optimal order. In light of the work of Scott and Vogelius (mentioned above) that conforming velocity elements of degree ≥ 4 already give optimal-order methods, the three lowest-order cases of nonconforming elements appear to be the only ones of interest.

An outline of this paper is as follows. In Section 2 we define the notation to be used and recall some of the theory of the finite element approximation of the stationary Stokes equations. In the now standard approach to this problem, the essential feature of the analysis is the verification of an appropriate form of the infsup condition. The verification of such a condition for the choice of nonconforming cubic velocities and discontinuous quadratic pressures is done in the remaining sections for various types of meshes.

2. Notation and Preliminaries. For Ω a polygonal domain in \mathbb{R}^2 , we thus consider the approximation of the stationary Stokes problem: Find $\mathbf{u} = (u_1, u_2)$ and p satisfying

$$-\nu \Delta \mathbf{u} + \nabla p = \mathbf{f} \quad \text{in } \Omega,$$
$$\operatorname{div} \mathbf{u} = 0 \quad \text{in } \Omega,$$
$$\mathbf{u} = 0 \quad \text{on } \partial \Omega,$$

where **u** is the fluid velocity, p is the pressure, **f** are the body forces per unit mass, and ν is the viscosity. The variational formulation of this problem is:

Find
$$\mathbf{u} = (u_1, u_2) \in (H_0^1(\Omega))^2$$
 and $p \in L^2(\Omega)/\mathbb{R}$ satisfying

$$a(\mathbf{u}, \mathbf{v}) - (p, \operatorname{div} \mathbf{v}) = (\mathbf{f}, \mathbf{v}) \quad \text{for all } \mathbf{v} \in (H_0^1(\Omega))^2,$$

 $(\operatorname{div} \mathbf{u}, q) = 0 \quad \text{for all } q \in L^2(\Omega)/\mathbb{R},$

where

$$a(\mathbf{u}, \mathbf{v}) = \nu \sum_{i,j=1}^{2} \int_{\Omega} \frac{\partial u_{i}}{\partial x_{j}} \frac{\partial v_{i}}{\partial x_{j}} dx$$

and (\cdot, \cdot) denotes the $L^2(\Omega)$ inner product.

The finite element methods we consider may be described abstractly as follows. We let τ_h denote a triangulation of Ω by triangles T of diameter $\leq h$. We then denote by \mathbf{W}_h a finite-dimensional approximation of $(H_0^1(\Omega))^2$. Since we are considering nonconforming methods, $\mathbf{W}_h \not\subseteq (H_0^1(\Omega))^2$, but $\mathbf{v}_h|_T \in H^1(T)$ for all $\mathbf{v}_h \in \mathbf{W}_h$ and $T \in \tau_h$. We assume that

$$\|\mathbf{v}_h\|_{1,h} = \left(\sum_{T \in \tau_h} |\mathbf{v}_h|_{1,T}^2\right)^{1/2}$$

is a norm on \mathbf{W}_h . Let Q_h denote a finite-dimensional subspace of $L^2(\Omega)/\mathbb{R}$. The approximation scheme is then: Find $\mathbf{u}_h \in \mathbf{W}_h$, $p_h \in Q_h$ satisfying

$$a_h(\mathbf{u}_h, \mathbf{v}_h) - (p_h, \operatorname{div}_h \mathbf{v}_h) = (\mathbf{f}, \mathbf{v}_h) \quad \text{for all } \mathbf{v}_h \in \mathbf{W}_h,$$

 $(\operatorname{div}_h \mathbf{u}_h, q) = 0 \quad \text{for all } q \in Q_h,$

where

$$a_h(\mathbf{u}, \mathbf{v}) = \nu \sum_{T \in \tau_h} \sum_{i, i=1}^2 \int_T \frac{\partial u_i}{\partial x_j} \frac{\partial v_i}{\partial x_j} dx$$

and $\operatorname{div}_h \mathbf{v}$ is the $L^2(\Omega)$ function whose restriction to each triangle $T \in \tau_h$ is given by $(\operatorname{div} \mathbf{v})|_T$.

The analysis of this type of method is well understood. If the space \mathbf{W}_h were conforming, the general theory of saddle-point problems developed by Babuška and Brezzi could be applied directly. In that case the only difficulty in the analysis is the verification of the inf-sup condition

(2.1)
$$\inf_{0 \neq q_h \in Q_h} \sup_{0 \neq \mathbf{v}_h \in \mathbf{W}_h} \frac{(\operatorname{div} \mathbf{v}_h, q_h)}{\|\mathbf{v}_h\|_1 \|q_h\|_0} \geq \gamma.$$

When (2.1) holds, one obtains the quasi-optimal error estimate

$$\|\mathbf{u} - \mathbf{u}_h\|_1 + \|p - p_h\|_0 \le C \inf(\|\mathbf{u} - \mathbf{v}_h\|_1 + \|p - q_h\|_0),$$

where the inf is taken over all $\mathbf{v}_h \in \mathbf{W}_h$ and $q_h \in Q_h$. For nonconforming methods, a straightforward modification of this result leads to the conclusion that if

(2.2)
$$\inf_{0 \neq q_h \in Q_h} \sup_{0 \neq \mathbf{v}_h \in \mathbf{W}_h} \frac{(\operatorname{div}_h \mathbf{v}_h, q_h)}{\|\mathbf{v}_h\|_{1,h} \|q_h\|_0} \geq \gamma,$$

then one obtains the error estimate

(2.3)
$$\|\mathbf{u} - \mathbf{u}_h\|_{1,h} + \|p - p_h\|_0$$

$$\leq C \left[\inf(\|\mathbf{u} - \mathbf{v}_h\|_{1,h} + \|p - q_h\|_0) + \sup \frac{\sum_T \int_{\partial T} \left(\nu \frac{\partial \mathbf{u}}{\partial n} - p\mathbf{n}\right) \cdot \mathbf{w}_h}{\|\mathbf{w}_h\|_{1,h}} \right],$$

where the inf is taken over all $\mathbf{v}_h \in \mathbf{W}_h$ and $q_h \in Q_h$ and the sup is taken over all $\mathbf{w}_h \in \mathbf{W}_h$.

Since [4] predated the work of Brezzi [3], the analysis in the former paper does not proceed by giving a direct verification of condition (2.2). However, using the interpolants constructed in [4], condition (2.2) can be easily verified on each triangle for the elements considered, with the global result following immediately. This local verification of condition (2.2) depends on the fact that appropriate bubble

functions have been added to the basic nonconforming spaces. In order to verify (2.2) without adding these bubble functions, we instead seek to verify (2.2) first on a patch of elements. This idea was previously used by Boland and Nicolaides [2] and Stenberg [11]. The basic technique is described in Girault and Raviart [8, pp. 129–132] for the case of conforming elements. In our case, we first define a collection of subdomains $\Omega_r \subset \Omega$, $r = 1, 2, \ldots, R$, such that $\overline{\Omega}_r$ is a union of triangles $T_{0,r}, T_{1,r}, \ldots, T_{M_r,r} \in \tau_h$ and

$$\overline{\Omega} = \bigcup_{r=1}^R \overline{\Omega}_r.$$

We assume that there exists a constant L independent of h such that for all $T \in \tau_h$ the number of r such that $\Omega_r \cap T \neq \emptyset$ is bounded by L. We next define

$$\mathbf{V}_h(\Omega_r) = \{\mathbf{u}_h : \mathbf{u}_h|_{T_{k,r}} \in (P_3)^2, k = 0, 1, \dots, M_r, \mathbf{u}_h = 0 \text{ in } \Omega/\Omega_r,$$

$$\mathbf{u}_h \text{ continuous at the Gauss points on all triangle sides}\},$$

$$H_h(\Omega_r) = \left\{ q \in L^2(\Omega_r) : q|_{T_{k,r}} \in P_2 \text{ and } \int_{T_{k,r}} q \, dx = 0, \right.$$

$$\left. k = 0, 1, \dots, M_r, \ q = 0 \text{ in } \Omega/\Omega_r \right\}.$$

The weak local verification of condition (2.2) on the patch of elements Ω_r will consist of establishing the validity of the following hypothesis:

Hypothesis H0. There exists a positive constant γ^* , independent of h and r, such that

(2.4)
$$\inf_{0 \neq q_h \in H_h} \sup_{0 \neq \mathbf{v}_h \in \mathbf{V}_h} \frac{(\operatorname{div}_h \mathbf{v}_h, q_h)}{\|\mathbf{v}_h\|_{1,h,\Omega_r} \|q_h\|_{0,\Omega_r}} \geq \gamma^*.$$

We shall do this by showing the validity of

Hypothesis H1. For all $q_h \in H_h(\Omega_r)$, there exists $\mathbf{u}_h \in \mathbf{V}_h(\Omega_r)$ such that

(2.5)
$$\operatorname{div}_{h}\mathbf{u}_{h} = q_{h} \quad \text{and} \quad \|\mathbf{u}_{h}\|_{1,h,\Omega_{r}} \leq \|q_{h}\|_{0,\Omega_{r}}/\gamma^{*}.$$

As in the papers mentioned above (e.g. see Theorem 1.12 of [8]), we can establish the following result. We include a proof for the sake of completeness and also since we do not assume, as in [2] and [8], that the domains Ω_r are disjoint.

THEOREM 2.1. Let

 $\mathbf{W}_h = \{\mathbf{u}_h : \mathbf{u}_h|_T \in (P_3)^2, \ T \in \tau_h, \mathbf{u}_h \ continuous \ at \ the \ Gauss \ points \\ on \ all \ triangle \ sides, \ \mathbf{u}_h = 0 \ at \ the \ Gauss \ points \ on \ \partial \Omega\},$

$$Q_h = \left\{ q \in L^2(\Omega), q|_T \in P_2, \ T \in \tau_h, \int_{\Omega} q \, dx = 0 \right\}.$$

Then, if hypothesis H1 is satisfied, condition (2.2) holds with a constant γ independent of h.

Proof. For $q_h \in Q_h$, one can construct by standard techniques a function $\mathbf{w}_h \in \mathbf{W}_h$ satisfying for all $T \in \tau_h$ the conditions

$$\int_T \operatorname{div}_h \mathbf{w}_h \, dx = \int_T q_h \, dx, \qquad \|\mathbf{w}_h\|_{1,h,\Omega} \le \|q_h\|_{0,\Omega}/\gamma_1,$$

with γ_1 independent of h. Setting

$$(2.6) p_h = q_h - \operatorname{div}_h \mathbf{w}_h,$$

we may write

$$p_h = \sum_{r=1}^R p_{h,r},$$

where

$$p_{h,r} \in H_h(\Omega_r)$$
 and $p_{h,r} \cdot p_{h,r'} = 0$ for $r \neq r'$.

Hence we have

$$||p_h||_{0,\Omega}^2 = \sum_{r=1}^R ||p_{h,r}||_{0,\Omega_r}^2.$$

Using (2.5), we may write

$$p_{h,r} = \operatorname{div}_h \mathbf{u}_{h,r} \quad \text{with } \|\mathbf{u}_{h,r}\|_{1,h,\Omega_r} \le \|p_{h,r}\|_{0,\Omega_r}/\gamma^*.$$

Then we set

(2.7)
$$\mathbf{u}_h = \sum_{r=1}^R \mathbf{u}_{h,r} \quad \text{and} \quad \mathbf{v}_h = \mathbf{u}_h + \mathbf{w}_h,$$

so that we have

$$\operatorname{div}_{h}\mathbf{v}_{h} = q_{h}$$

and

$$\|\mathbf{u}_{h}\|_{1,h,\Omega}^{2} = \sum_{T \in \tau_{h}} \|\mathbf{u}_{h}\|_{1,T}^{2} \leq L \sum_{T \in \tau_{h}} \sum_{r=1}^{R} \|\mathbf{u}_{h,r}\|_{1,T}^{2}$$

$$\leq L \sum_{r=1}^{R} \|\mathbf{u}_{h,r}\|_{1,h,\Omega_{r}}^{2} \leq L \|p_{h}\|_{0,\Omega}^{2} / \gamma^{*2}.$$

From (2.6) it follows that

$$||p_h||_{0,\Omega} \le ||q_h||_{0,\Omega} + ||\operatorname{div}_h \mathbf{w}_h||_{0,\Omega} \le ||q_h||_{0,\Omega} + \sqrt{2} ||\mathbf{w}_h||_{1,h,\Omega} \le (1 + \sqrt{2}/\gamma_1) ||q_h||_{0,\Omega}.$$

Using (2.6) and (2.7), we then obtain

 $\|\mathbf{v}_h\|_{1,h,\Omega} \le \|\mathbf{u}_h\|_{1,h,\Omega} + \|\mathbf{w}_h\|_{1,h,\Omega} \le \sqrt{L}(\gamma_1 + \sqrt{2})\|q_h\|_{0,\Omega}/(\gamma^*\gamma_1) + \|q_h\|_{0,\Omega}/\gamma_1.$ Finally,

$$\|\mathbf{v}_h\|_{1,h,\Omega} \leq \|q_h\|_{0,\Omega}/\gamma,$$

where

$$\gamma = \gamma^* \gamma_1 / [\sqrt{L}(\gamma_1 + \sqrt{2}) + \gamma^*].$$

Together with (2.8), this implies (2.2).

It then follows directly from (2.3) and standard estimates for nonconforming methods (see [4] and [9]) that if $\mathbf{u} \in (H^4(\Omega))^2$ and $p \in H^3(\Omega)$, then

(2.9)
$$\|\mathbf{u} - \mathbf{u}_h\|_{1,h} + \|p - p_h\|_0 \le Ch^3(\|\mathbf{u}\|_4 + \|p\|_3).$$

Thus, to establish Theorem 2.1 and the error estimate (2.9), we need only show that Hypothesis H1 is valid. Before doing so in the next sections, we first state some basic facts that will facilitate the forthcoming analysis. We recall that the barycentric coordinates $\lambda_i = \lambda_i(\mathbf{x})$, $1 \le i \le 3$, of any point $\mathbf{x} \in \mathbb{R}^2$ with respect to the points A_1 , A_2 , A_3 are the unique solutions of the linear system

(2.10)
$$\sum_{j=1}^{3} A_j \lambda_j = \mathbf{x},$$

(2.11)
$$\sum_{i=1}^{3} \lambda_{i} = 1.$$

Denoting by $\overrightarrow{A_1 A_2}$ the vector from A_1 to A_2 , we have

(2.12)
$$\mathbf{x} = A_1 + \lambda_2 \overrightarrow{A_1 A_2} + \lambda_3 \overrightarrow{A_1 A_3}.$$

Hence,

(2.13)
$$I = \nabla \mathbf{x}^t = \nabla \lambda_2 (\overrightarrow{A_1 A_2})^t + \nabla \lambda_3 (\overrightarrow{A_1 A_3})^t,$$

and for any vector u,

(2.14)
$$\mathbf{u} = (\mathbf{u} \cdot \nabla \lambda_2) \overrightarrow{A_1 A_2} + (\mathbf{u} \cdot \nabla \lambda_3) \overrightarrow{A_1 A_3}.$$

By applying (2.13) to the vector $\nabla \lambda_2$, it easily follows that

(2.15)
$$\nabla \lambda_2 \cdot \overrightarrow{A_1 A_2} = 1 \quad \text{and} \quad \nabla \lambda_2 \cdot \overrightarrow{A_1 A_3} = 0.$$

3. Verification of Hypothesis H1—Mesh I. The verification of condition (2.5) will be done for three choices of the subdomain Ω_r . In the first of these, $\Omega_r = T_0 \cup T_1 \cup T_2 \cup T_3$, where the four triangles T_0 , T_1 , T_2 , T_3 are aligned in the configuration depicted in Figure 3.1.

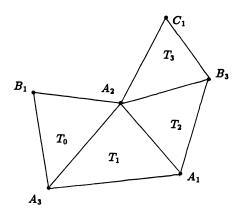


FIGURE 3.1

We denote by

 (μ_1, μ_2, μ_3) the barycentric coordinates of B_3 with respect to A_1, A_2, A_3 ,

 (ν_1, ν_2, ν_3) the barycentric coordinates of B_1 with respect to A_1, A_2, A_3 ,

 (τ_1, τ_2, τ_3) the barycentric coordinates of C_1 with respect to A_1, A_2, B_3 .

The verification of Hypothesis H1 will depend on a geometrical condition on the mesh patch Ω_r given in terms of the barycentric coordinates defined above. Following the verification, we describe some common situations which are and are not covered. To simplify the proof of our result, it is convenient to first define the following subspace of $\mathbf{V}_h(\Omega_r)$. Let

$$\hat{\mathbf{V}}_h(\Omega_r) = \left\{ \mathbf{u}_h \in \mathbf{V}_h(\Omega_r) : \int_s \mathbf{u}_h \, d\sigma = 0 \text{ for all sides } s \text{ of } \Omega_r \right\}.$$

The verification of condition (2.5) is contained in the following proposition.

PROPOSITION 3.1. If $D \equiv -\nu_2 \tau_2 \mu_1^2 - (\nu_1 - 1)\tau_2 \mu_2 \mu_3 + (\tau_1 - 1)\nu_2 \mu_2 \neq 0$, then the operator div_h from $\hat{\mathbf{V}}_h$ into H_h is one-to-one and for all $\mathbf{u}_h \in \hat{\mathbf{V}}_h$

$$\|\mathbf{u}_h\|_{1,h,\Omega_r} \leq \frac{C(\theta)}{|D|} \|\operatorname{div}_h \mathbf{u}_h\|_{0,\Omega_r},$$

where C is a constant which depends only on the smallest angle θ of the triangles T_0 , T_1 , T_2 , and T_3 .

Proof. First we remark that $\mathbf{u}_h \in \hat{\mathbf{V}}_h$ implies $\int_{T_k} \operatorname{div}_h \mathbf{u}_h \, dx = 0$, so that $\operatorname{div}_h(\hat{\mathbf{V}}_h) \subseteq H_h$. Since $\operatorname{dim} \hat{\mathbf{V}}_h = \operatorname{dim} H_h = 20$, div_h is one-to-one if and only if $\operatorname{Ker}(\operatorname{div}_h) = \{0\}$.

In the triangle T_k , we may write

$$\mathbf{u}_{h} = \sum_{i=1}^{3} \mathbf{u}_{i}^{(k)} \lambda_{i}^{3} + \sum_{i \neq j} \mathbf{u}_{ij}^{(k)} \lambda_{i}^{2} \lambda_{j} + \mathbf{u}_{123}^{(k)} \lambda_{1} \lambda_{2} \lambda_{3},$$

$$q = \sum_{i=1}^{3} q_{i}^{(k)} \lambda_{i}^{2} + \sum_{i < j} q_{ij}^{(k)} \lambda_{i} \lambda_{j},$$

where $\lambda_i = \lambda_i^{(k)}$ are the barycentric coordinates in the triangle T_k . We shall omit the superscript (k) when there is no ambiguity.

We now turn to the lengthy process of expressing the continuity of \mathbf{u}_h at the Gauss points and the satisfaction of the equations $\operatorname{div}_h \mathbf{u}_h = q$ on each triangle T and $\int_s \mathbf{u}_h d\sigma = 0$ for each triangle side s, in terms of the degrees of freedom of \mathbf{u}_h and q. On the side $\overrightarrow{A_2A_3}$ or $\overrightarrow{A_2B_3}$, the condition $\int_s \mathbf{u}_h d\sigma = 0$ becomes

(3.1)
$$3(\mathbf{u}_2 + \mathbf{u}_3) + (\mathbf{u}_{23} + \mathbf{u}_{32}) = 0.$$

Similarly, we have on $\overrightarrow{A_2B_1}$, $\overrightarrow{A_2A_1}$, $\overrightarrow{A_2C_1}$

(3.2)
$$3(\mathbf{u}_2 + \mathbf{u}_1) + (\mathbf{u}_{21} + \mathbf{u}_{12}) = 0.$$

The continuity at the Gauss points on $\overrightarrow{A_2A_3}$ (barycentric coordinates $(0, 1/2, 1/2), (0, \theta, 1 - \theta), (0, 1 - \theta, \theta)$ with $\theta(1 - \theta) = 1/10$) gives

$$(3.3) \qquad (\mathbf{u}_2 + \mathbf{u}_3 + \mathbf{u}_{23} + \mathbf{u}_{32})^{(0)} = (\mathbf{u}_2 + \mathbf{u}_3 + \mathbf{u}_{23} + \mathbf{u}_{32})^{(1)},$$

(3.4)
$$(\theta^{3}\mathbf{u}_{2} + (1-\theta)^{3}\mathbf{u}_{3} + \theta(1-\theta)[\theta\mathbf{u}_{23} + (1-\theta)\mathbf{u}_{32}])^{(0)}$$

$$= (\theta^{3}\mathbf{u}_{2} + (1-\theta)^{3}\mathbf{u}_{3} + \theta(1-\theta)[\theta\mathbf{u}_{23} + (1-\theta)\mathbf{u}_{32}])^{(1)},$$

(3.5)
$$([1-\theta]^3 \mathbf{u}_2 + \theta^3 \mathbf{u}_3 + \theta(1-\theta)[(1-\theta)\mathbf{u}_{23} + \theta\mathbf{u}_{32}])^{(0)}$$

$$= ([1-\theta]^3 \mathbf{u}_2 + \theta^3 \mathbf{u}_3 + \theta(1-\theta)[(1-\theta)\mathbf{u}_{23} + \theta\mathbf{u}_{32}])^{(1)}.$$

Since $\theta(1-\theta)=1/10$ implies $\theta^2=\theta-1/10$ and $\theta^3=(9\theta-1)/10$, we have

$$\theta^{3}\mathbf{u}_{2} + (1-\theta)^{3}\mathbf{u}_{3} + \theta(1-\theta)[\theta\mathbf{u}_{23} + (1-\theta)\mathbf{u}_{32}]$$

$$= \{(9\theta - 1)\mathbf{u}_{2} + (8-9\theta)\mathbf{u}_{3} + [\theta\mathbf{u}_{23} + (1-\theta)\mathbf{u}_{32}]\}/10$$

$$= \{(1-\theta)[3(\mathbf{u}_{2} + \mathbf{u}_{3}) + \mathbf{u}_{23} + \mathbf{u}_{32}] + (2\theta - 1)(9\mathbf{u}_{2} + \mathbf{u}_{23})$$

$$+ (5-6\theta)(\mathbf{u}_{2} + \mathbf{u}_{3})\}/10.$$

Using an analogous rewriting of Eq. (3.5) and taking into account Eq. (3.1), we have the continuity at the Gauss points on $\overrightarrow{A_2A_3}$ and $\overrightarrow{A_2B_3}$ if and only if

(3.6)
$$(\mathbf{u}_2 + \mathbf{u}_3)^{(k)} = (\mathbf{u}_2 + \mathbf{u}_3)^{(k+1)}, \qquad k = 0, 2,$$

(3.7)
$$(9\mathbf{u}_2 + \mathbf{u}_{23})^{(k)} = (9\mathbf{u}_2 + \mathbf{u}_{23})^{(k+1)}, \qquad k = 0, 2.$$

Similarly, continuity at the Gauss points on $\overrightarrow{A_2A_1}$ becomes

$$(3.8) (\mathbf{u}_1 + \mathbf{u}_2)^{(1)} = (\mathbf{u}_1 + \mathbf{u}_2)^{(2)},$$

(3.9)
$$(9\mathbf{u}_1 + \mathbf{u}_{12})^{(1)} = (9\mathbf{u}_1 + \mathbf{u}_{12})^{(2)}.$$

The continuity with 0 at the Gauss points of $\overrightarrow{A_1A_3}$, $\overrightarrow{A_1B_3}$, $\overrightarrow{B_1A_3}$, $\overrightarrow{C_1B_3}$ (which implies $\int_s \mathbf{u}_h d\sigma = 0$) may be written

$$(3.10) \mathbf{u}_1 + \mathbf{u}_3 = 0,$$

$$(3.11) 9\mathbf{u}_1 + \mathbf{u}_{13} = 0,$$

$$(3.12) 9\mathbf{u}_3 + \mathbf{u}_{31} = 0.$$

Noticing that

$$\int_T q \, dx = 0 \Leftrightarrow 2(q_1 + q_2 + q_3) + q_{12} + q_{23} + q_{13} = 0,$$

the equation $\operatorname{div}_h \mathbf{u}_h = q$ on T is equivalent to

$$(3.13) 3\mathbf{u}_1 \cdot \nabla \lambda_1 + \mathbf{u}_{12} \cdot \nabla \lambda_2 + \mathbf{u}_{13} \cdot \nabla \lambda_3 = q_1,$$

$$(3.14) 3\mathbf{u}_2 \cdot \nabla \lambda_2 + \mathbf{u}_{23} \cdot \nabla \lambda_3 + \mathbf{u}_{21} \cdot \nabla \lambda_1 = q_2,$$

$$(3.15) 3\mathbf{u}_3 \cdot \nabla \lambda_3 + \mathbf{u}_{31} \cdot \nabla \lambda_1 + \mathbf{u}_{32} \cdot \nabla \lambda_2 = q_3,$$

(3.16)
$$2\mathbf{u}_{23} \cdot \nabla \lambda_2 + 2\mathbf{u}_{32} \cdot \nabla \lambda_3 + \mathbf{u}_{123} \cdot \nabla \lambda_1 = q_{23},$$

$$(3.17) 2\mathbf{u}_{31} \cdot \nabla \lambda_3 + 2\mathbf{u}_{13} \cdot \nabla \lambda_1 + \mathbf{u}_{123} \cdot \nabla \lambda_2 = q_{13}.$$

Taking into account (3.2), continuity with 0 at the Gauss points of $\overrightarrow{A_2B_1}$ and $\overrightarrow{A_2C_1}$ is equivalent to

(3.18)
$$(\mathbf{u_1} + \mathbf{u_2})^{(k)} = 0, \quad k = 0, 3,$$

(3.19)
$$(9\mathbf{u}_1 + \mathbf{u}_{12})^{(k)} = 0, \quad k = 0, 3.$$

Using all of these equations, we now study $Ker(div_h)$. To do so, we set q = 0. In triangle T_0 we obtain from (3.10), (3.11), (3.12), (3.2), (3.18), and (3.19) that

$$\mathbf{u_2} = \mathbf{u_3} = -\mathbf{u_1}, \qquad \mathbf{u_{31}} = \mathbf{u_{21}} = -\mathbf{u_{13}} = -\mathbf{u_{12}} = 9\mathbf{u_1}.$$

We deduce from (3.13) that

$$(3.20) 12\mathbf{u}_1 \cdot \nabla \lambda_1 = 0,$$

from (3.14) that

$$\mathbf{u_{23}} \cdot \nabla \lambda_3 = 3\mathbf{u_1} \cdot \nabla \lambda_2$$

and from (3.15) that

$$\mathbf{u_{32}} \cdot \nabla \lambda_2 = 3\mathbf{u}_1 \cdot \nabla \lambda_3.$$

Therefore,

$$(9\mathbf{u_2} + \mathbf{u_{23}}) \cdot \nabla \lambda_3 = -9\mathbf{u_1} \cdot \nabla \lambda_3 + 3\mathbf{u_1} \cdot \nabla \lambda_2 = 12\mathbf{u_1} \cdot \nabla \lambda_2,$$

and using (3.1),

$$(9\mathbf{u_2} + \mathbf{u_{23}}) \cdot \nabla \lambda_2 = (6\mathbf{u_2} - 3\mathbf{u_3} - \mathbf{u_{32}}) \cdot \nabla \lambda_2 = 0.$$

Let

(3.21)
$$\alpha = (\mathbf{u_1} \cdot \nabla \lambda_2)^{(0)}.$$

Then

$$(\mathbf{u_2} + \mathbf{u_3})^{(0)} = 2\alpha \overrightarrow{A_2} \overrightarrow{A_3},$$

 $(9\mathbf{u_2} + \mathbf{u_{23}})^{(0)} = 12\alpha \overrightarrow{B_1} \overrightarrow{A_3}.$

In triangle T_1 we use (3.6) and (3.7) with k = 0 to obtain

(3.22)
$$(\mathbf{u_2} + \mathbf{u_3})^{(1)} = 2\alpha \overrightarrow{A_2} \overrightarrow{A_3},$$

(3.23)
$$(9\mathbf{u_2} + \mathbf{u_{23}})^{(1)} = 12\alpha \overrightarrow{B_1 A_3},$$

and using (3.1),

$$(9\mathbf{u}_3 + \mathbf{u}_{32})^{(1)} = -12\alpha \overrightarrow{B_1 A_2}.$$

Next, we use (3.15) and (3.12) to get

$$12(\mathbf{u}_3 \cdot \nabla \lambda_3 - \alpha \overrightarrow{B_1 A_2} \cdot \nabla \lambda_2) = 0,$$

i.e.,

(3.24)
$$\mathbf{u}_3 \cdot \nabla \lambda_3 = \alpha (1 - \nu_2).$$

Let

(3.25)
$$\beta = (\mathbf{u}_3 \cdot \nabla \lambda_1)^{(1)}.$$

Then

$$\mathbf{u}_3 = \beta \overrightarrow{A_2} \overrightarrow{A_1} + \alpha (1 - \nu_2) \overrightarrow{A_2} \overrightarrow{A_3},$$

and from (3.10),

$$\mathbf{u}_1 + \mathbf{u}_2 = \mathbf{u}_2 - \mathbf{u}_3 = 2\alpha \overrightarrow{A_2 A_3} - 2\mathbf{u}_3 = -2(\beta \overrightarrow{A_2 A_1} - \alpha \overrightarrow{\nu_2 A_2 A_3}),$$

i.e.,

(3.26)
$$(\mathbf{u}_1 + \mathbf{u}_2)^{(1)} = -2(\beta \overrightarrow{A_2} \overrightarrow{A_1} - \alpha \nu_2 \overrightarrow{A_2} \overrightarrow{A_3}).$$

From (3.13) and (3.11) we obtain

(3.27)
$$(9\mathbf{u}_1 + \mathbf{u}_{12}) \cdot \nabla \lambda_2 = 9\mathbf{u}_1 \cdot \nabla \lambda_2 - 3\mathbf{u}_1 \cdot \nabla \lambda_1 - \mathbf{u}_{13} \cdot \nabla \lambda_3$$
$$= -12\mathbf{u}_1 \cdot \nabla \lambda_1 = 12\mathbf{u}_3 \cdot \nabla \lambda_1 = 12\beta.$$

From (3.2), (3.10), (3.14), and (3.22),

$$\begin{aligned} (9\mathbf{u}_1 + \mathbf{u}_{12}) \cdot \nabla \lambda_1 &= (6\mathbf{u}_1 - 3\mathbf{u}_2 - \mathbf{u}_{21}) \cdot \nabla \lambda_1 \\ &= (6\mathbf{u}_1 - 3\mathbf{u}_2) \cdot \nabla \lambda_1 + 3\mathbf{u}_2 \cdot \nabla \lambda_2 + \mathbf{u}_{23} \cdot \nabla \lambda_3 \\ &= (6\mathbf{u}_1 - 3\mathbf{u}_2) \cdot \nabla \lambda_1 + 3\mathbf{u}_2 \cdot \nabla \lambda_2 - 9\mathbf{u}_2 \cdot \nabla \lambda_3 + 12\alpha \overrightarrow{B_1 A_3} \cdot \nabla \lambda_3 \\ &= 6\mathbf{u}_1 \cdot \nabla \lambda_1 + 6\mathbf{u}_2 \cdot (\nabla \lambda_2 - \nabla \lambda_3) + 12\alpha(1 - \nu_3) \\ &= -6\mathbf{u}_3 \cdot \nabla \lambda_1 + 12\alpha \overrightarrow{A_2 A_3} \cdot (\nabla \lambda_2 - \nabla \lambda_3) - 6\mathbf{u}_3 \cdot (\nabla \lambda_2 - \nabla \lambda_3) \\ &+ 12\alpha(1 - \nu_3) \\ &= 12\mathbf{u}_3 \cdot \nabla \lambda_3 - 12\alpha(1 + \nu_3) = 12\alpha(\nu_1 - 1). \end{aligned}$$

With (3.27) that gives

(3.28)
$$(9\mathbf{u}_1 + \mathbf{u}_{12})^{(1)} = 12[\alpha(\nu_1 - 1)\overrightarrow{A_3}\overrightarrow{A_1} + \beta\overrightarrow{A_3}\overrightarrow{A_2}].$$

Similarly, we obtain

(3.29)
$$(\mathbf{u}_1 + \mathbf{u}_2)^{(2)} = -2(\delta \overrightarrow{A_2} \overrightarrow{A_1} - \gamma \tau_2 \overrightarrow{A_2} \overrightarrow{B_3}),$$

(3.30)
$$(9\mathbf{u}_1 + \mathbf{u}_{12})^{(2)} = 12[\gamma(\tau_1 - 1)\overrightarrow{B_3}\overrightarrow{A_1} + \delta\overrightarrow{B_3}\overrightarrow{A_2}],$$

where

$$(3.31) \gamma = (\mathbf{u}_1 \cdot \nabla \lambda_2)^{(3)}$$

and

(3.32)
$$\delta = (\mathbf{u}_3 \cdot \nabla \lambda_1)^{(2)}.$$

We remark that

$$\begin{split} \overrightarrow{A_2B_3} &= \mu_1\overrightarrow{A_2A_1} + \mu_3\overrightarrow{A_2A_3}, \\ \overrightarrow{B_3A_1} &= (1-\mu_1)\overrightarrow{A_3A_1} - \mu_2\overrightarrow{A_3A_2}, \\ \overrightarrow{B_3A_2} &= -\mu_1\overrightarrow{A_3A_1} + (1-\mu_2)\overrightarrow{A_3A_2}. \end{split}$$

Thus Eq. (3.8) becomes (from (3.26) and (3.29))

$$\beta = \delta - \tau_2 \mu_1 \gamma,$$

$$\nu_2 \alpha = \tau_2 \mu_3 \gamma,$$

and Eq. (3.9) becomes

$$(\nu_1 - 1)\alpha = (1 - \mu_1)(\tau_1 - 1)\gamma - \mu_1\delta,$$

$$\beta = -(\tau_1 - 1)\mu_2\gamma + (1 - \mu_2)\delta.$$

We thus obtain the linear system

$$\begin{pmatrix} 0 & 1 & \mu_1 \tau_2 & -1 \\ \nu_2 & 0 & -\mu_3 \tau_2 & 0 \\ \nu_1 - 1 & 0 & -(\tau_1 - 1)(1 - \mu_1) & \mu_1 \\ 0 & 1 & \mu_2(\tau_1 - 1) & -(1 - \mu_2) \end{pmatrix} \begin{pmatrix} \alpha \\ \beta \\ \gamma \\ \delta \end{pmatrix} = 0.$$

The determinant of this linear system is

$$D = -(\nu_1 - 1)\tau_2\mu_3\mu_2 + \nu_2[(\tau_1 - 1)\mu_2 - \mu_1^2\tau_2].$$

Therefore, if $D \neq 0$, we have $\alpha = \beta = \gamma = \delta = 0$. We then deduce from (3.20) and (3.21) that $\mathbf{u}_1^{(0)} = 0$. Consequently,

$$\mathbf{u}_{2}^{(0)} = \mathbf{u}_{3}^{(0)} = \mathbf{u}_{12}^{(0)} = \mathbf{u}_{21}^{(0)} = \mathbf{u}_{31}^{(0)} = \mathbf{u}_{13}^{(0)} = \mathbf{u}_{23}^{(0)} = 0.$$

It then follows from (3.1) that $\mathbf{u}_{32}^{(0)} = 0$ and from (3.16) and (3.17) that $\mathbf{u}_{123}^{(0)} = 0$. Hence $\mathbf{u}_h^{(0)} = 0$. In a similar way, we obtain $\mathbf{u}_h^{(1)} = \mathbf{u}_h^{(2)} = \mathbf{u}_h^{(3)} = 0$, so that $\mathbf{u}_h = 0$. Since this shows that $\operatorname{Ker}(\operatorname{div}_h) = \{0\}$, we get that div_h is one-to-one.

Finally, we consider the inequality stated in Proposition 3.1. To prove this result, we consider the previous equations with $q \neq 0$. First observe that Eqs. (3.1)–(3.15) give the values of

$$(\mathbf{u}_i \cdot \nabla \lambda_l)^{(k)}, (\mathbf{u}_{ij} \cdot \nabla \lambda_l)^{(k)}, \qquad k = 0, 1, 2, 3,$$

as linear functions of $(q_1, q_2, q_3)^{(k')}$, k' = 0, 1, 2, 3, and thus we obtain

$$|(\mathbf{u}_{i} \cdot \nabla \lambda_{l})^{(k)}| \leq \frac{C_{1}}{|D|} \max_{i,k'} |q_{i}^{(k')}|,$$

$$|(\mathbf{u}_{ij} \cdot \nabla \lambda_{l})^{(k)}| \leq \frac{C_{2}}{|D|} \max_{i,k'} |q_{i}^{(k')}|.$$

We next use (3.16) and (3.17) to obtain

$$\begin{split} |(\mathbf{u}_{123} \cdot \nabla \lambda_1)^{(k)}| &\leq |q_{23}^{(k)}| + \frac{4C_2}{|D|} \max_{i,k'} |q_i^{(k')}|, \\ |(\mathbf{u}_{123} \cdot \nabla \lambda_2)^{(k)}| &\leq |q_{13}^{(k)}| + \frac{4C_2}{|D|} \max_{i,k'} |q_i^{(k')}|. \end{split}$$

By standard scaling arguments, we have that

$$\begin{aligned} |q_i^{(k)}| &\leq C ||q||_{L^2(T_k)} / A_k^{1/2}, \\ |q_{ij}^{(k)}| &\leq C ||q||_{L^2(T_k)} / A_k^{1/2}, \\ ||\operatorname{Grad} \mathbf{u}||_{L^2(T_k)} &\leq C(\theta) \left(\sum |\mathbf{u}_i| + \sum |\mathbf{u}_{ij}| + |\mathbf{u}_{123}| \right), \end{aligned}$$

where A_k denotes the area of T_k . The inequality stated in Proposition 3.1 follows immediately.

In order to understand the condition on the determinant D under which Proposition 3.1 is valid, we now examine this condition in more detail. For the configuration of triangles depicted in Figure 3.1, we have that $\nu_1 < 0$, $\mu_3 < 0$, $\tau_1 < 0$. If we assume that

$$\nu_2 \ge 0, \ \tau_2 \ge 0, \ \mu_2 \ge 0, \ \nu_2 + \mu_2 \ne 0, \ \nu_2 + \tau_2 \ne 0, \ \tau_2 + \mu_2 \ne 0,$$

then it is easy to see that $D \neq 0$. Some meshes not covered by Proposition 1 are:

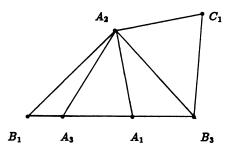


FIGURE 3.2. $\nu_2=\mu_2=0$

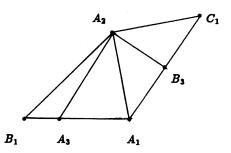


FIGURE 3.3. $\nu_2 = \tau_2 = 0$

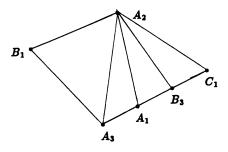


FIGURE 3.4. $au_2 = \mu_2 = 0$

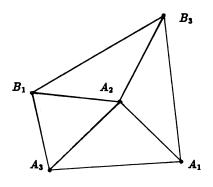


FIGURE 3.5

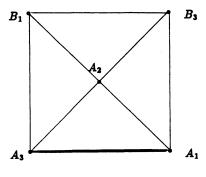


FIGURE 3.6

Remark. The limiting case $C_1 = B_1$ is allowed. In this case, $\mathbf{u}_h = 0$ at the Gauss points on $\overrightarrow{B_1 A_2}$. Therefore, the mesh depicted in Figure 3.5 and, in particular, the special case depicted in Figure 3.6 are covered by Proposition 3.1.

4. Verification of Hypothesis H1—Mesh II. We next consider the mesh in Figure 4.1.

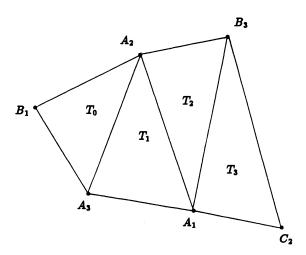


FIGURE 4.1

We use the same notation with the addition that (ξ_1, ξ_2, ξ_3) are the barycentric coordinates of C_2 with respect to A_1 , A_2 , B_3 . The verification of condition (2.5) is contained in the following proposition.

Proposition 4.1. If

$$\Delta = (\nu_1 - 1)\xi_1\mu_2\mu_3 - \nu_2(\xi_2 - 1)\mu_1 + \nu_2\xi_1[(1 - \mu_1)^2 - \mu_3] \neq 0,$$

then the operator div_h from $\hat{\mathbf{V}}_h$ into H_h is one-to-one and for all $\mathbf{u}_h \in \hat{\mathbf{V}}_h$,

$$\|\mathbf{u}_h\|_{1,h,\Omega_r} \leq \frac{C(\theta)}{|\Delta|} \|\mathrm{div}_h \mathbf{u}_h\|_{0,\Omega_r}.$$

Proof. We follow the proof of Proposition 3.1 up to Eq. (3.27). In triangle T_1 , it follows from (3.2) that

$$9\mathbf{u}_2 + \mathbf{u}_{21} = 6(\mathbf{u}_2 + \mathbf{u}_1) - 9\mathbf{u}_1 - \mathbf{u}_{12}$$

and from (3.26) and (3.27) that

$$(4.1) (9\mathbf{u}_2 + \mathbf{u}_{21})^{(1)} = 12\{\alpha\nu_2\overrightarrow{A_2}A_3 + [\alpha(\nu_1 - 1) + \beta]\overrightarrow{A_1}A_3\}.$$

By symmetry, we obtain

(4.2)
$$(\mathbf{u}_1 + \mathbf{u}_2)^{(2)} = -2[\delta \overrightarrow{A_1 A_2} - \gamma \xi_1 \overrightarrow{A_1 B_3}],$$

$$(4.3) (9\mathbf{u}_1 + \mathbf{u}_{12})^{(2)} = 12\{\gamma \xi_1 \overrightarrow{A_1 B_3} + [\gamma(\xi_2 - 1) + \delta]\overrightarrow{A_2 B_3}\}.$$

Then (3.8) and (3.9) become

$$\begin{split} \beta &= -\delta - \gamma \xi_1(\mu_1 - 1), \\ \nu_2 \alpha &= \xi_1 \mu_3 \gamma, \\ (\nu_1 - 1) \alpha &= (\mu_1 - 1) \xi_1 \gamma + \mu_1 [(\xi_2 - 1) \gamma + \delta], \\ \beta &= \mu_2 \xi_1 \gamma + (\mu_2 - 1) [(\xi_2 - 1) \gamma + \delta]. \end{split}$$

We thus obtain the linear system

$$\begin{pmatrix} 0 & 1 & (\mu_1-1)\xi_1 & 1 \\ \nu_2 & 0 & -\mu_3\xi_1 & 0 \\ \nu_1-1 & 0 & \xi_1-\mu_1(\xi_1+\xi_2-1) & -\mu_1 \\ 0 & 1 & -\xi_1\mu_2-(\mu_2-1)(\xi_2-1) & (1-\mu_2) \end{pmatrix} \begin{pmatrix} \alpha \\ \beta \\ \gamma \\ \delta \end{pmatrix} = 0.$$

The remainder of the proof of Proposition 4.1 follows in an analogous fashion to that of Proposition 3.1.

We now examine the determinant Δ in more detail in order to understand the condition under which Proposition 4.1 is valid. For the configuration of triangles depicted in Figure 4.1, we have that $\nu_1 < 0$, $\mu_3 < 0$, $\xi_2 < 0$. If we assume that

$$\xi_1 \ge 0, \ \mu_2 \ge 0, \ \mu_1 \ge 0, \ \nu_2 \ge 0, \ \nu_2 + \xi_1 \ne 0, \ \nu_2 + \mu_2 \ne 0, \ \xi_1 + \mu_1 \ne 0,$$

then it is easy to see that $\Delta \neq 0$. Some meshes not covered by Proposition 4.1 are

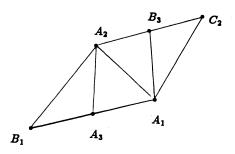


FIGURE 4.2. $\nu_2 = \xi_1 = 0$

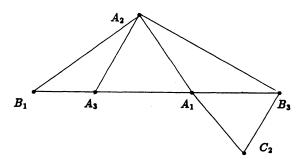


FIGURE 4.3. $\nu_2 = \mu_2 = 0$

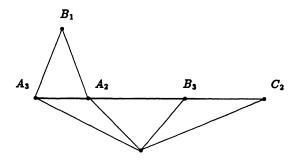


FIGURE 4.4. $\xi_1 = \mu_1 = 0$

However, the standard regular mesh, depicted in Figure 4.5, is covered by Proposition 4.1.

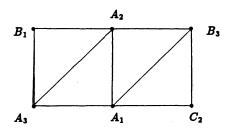


FIGURE 4.5

5. Verification of Hypothesis H1—Mesh III. Next, we turn to the case of three triangles aligned as in Figure 5.1.

We set $\Omega_r = T = T_1 \cup T_2 \cup T_3$. Let (ν_1, ν_2, ν_3) denote the barycentric coordinates of B with respect to the points A_1 , A_2 , A_3 and let (μ_1, μ_2, μ_3) denote the barycentric coordinates of a point $x \in T$ with respect to A_1 , A_2 , A_3 . To define the barycentric coordinates in each T_i , we also denote the vertex B in the triangle T_i by B_i . We then denote by $\lambda_j^{(i)}$ the barycentric coordinate of a point $x \in T_i$ with respect to the vertex of T_i with subscript j. Define the subspace X_h of V_h by

$$\mathbf{X}_h = \{\mathbf{u}_h \in (C^0(\mathbb{R}^2))^2 \colon \mathbf{u}_h|_{T_i} \in (P_3)^2, \mathbf{u}_h(B) = 0, \mathbf{u}_h = 0 \text{ in } \Omega/T\}.$$

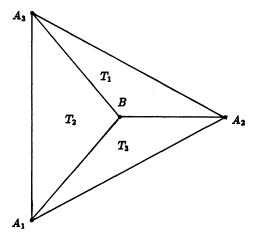


FIGURE 5.1

Note that $\mathbf{u}_h \in \mathbf{X}_h$ may be represented in the triangle T_i by

$$\mathbf{u}_{h} = \sum_{\substack{j=1\\j\neq i}}^{3} \mathbf{u}_{ij}^{(i)} \lambda_{i}^{2} \lambda_{j} + \sum_{\substack{j=1\\j\neq i}}^{3} \mathbf{u}_{ji}^{(i)} \lambda_{j}^{2} \lambda_{i} + \mathbf{u}_{123}^{(i)} \lambda_{1} \lambda_{2} \lambda_{3},$$

where we have omitted the superscript (i) on the barycentric coordinates. We will continue to omit superscripts (on \mathbf{u}_{ij} as well) when there is no ambiguity about the triangle. Using the continuity of \mathbf{u}_h , it follows by comparing the expressions for \mathbf{u}_h along common interior sides, that

$$\begin{aligned} \mathbf{u}_{13}^{(1)} &= \mathbf{u}_{23}^{(2)}, & \mathbf{u}_{31}^{(1)} &= \mathbf{u}_{32}^{(2)}, \\ \mathbf{u}_{12}^{(2)} &= \mathbf{u}_{13}^{(3)}, & \mathbf{u}_{21}^{(2)} &= \mathbf{u}_{31}^{(3)}, \\ \mathbf{u}_{23}^{(3)} &= \mathbf{u}_{21}^{(1)}, & \mathbf{u}_{32}^{(3)} &= \mathbf{u}_{12}^{(1)}. \end{aligned}$$

Using the above, it is easy to check that dim $X_h = 18$, while H_h , defined in Section 2, has dimension 15. To get a result analogous to Propositions 3.1 and 4.1, we define

$$\hat{\mathbf{V}}_h = \left\{ \mathbf{u}_h \in \mathbf{X}_h : \int_{s} \mathbf{u}_h \cdot \boldsymbol{\nu} \, d\sigma = 0, \text{ for all sides } s \text{ of } \Omega_r \right\},$$

where ν is the normal to s. Note that $\int_s \mathbf{u}_h \cdot \boldsymbol{\nu} \, d\sigma = 0$ is equivalent to the conditions

(5.2)
$$(\mathbf{u}_{31}^{(1)} + \mathbf{u}_{13}^{(1)}) \cdot \nabla \lambda_2^{(1)} = (\mathbf{u}_{21}^{(1)} + \mathbf{u}_{12}^{(1)}) \cdot \nabla \lambda_3^{(1)} = 0,$$

(5.3)
$$(\mathbf{u}_{12}^{(2)} + \mathbf{u}_{21}^{(2)}) \cdot \nabla \lambda_3^{(2)} = (\mathbf{u}_{32}^{(2)} + \mathbf{u}_{23}^{(2)}) \cdot \nabla \lambda_1^{(2)} = 0,$$

(5.4)
$$(\mathbf{u}_{32}^{(3)} + \mathbf{u}_{23}^{(3)}) \cdot \nabla \lambda_1^{(3)} = (\mathbf{u}_{13}^{(3)} + \mathbf{u}_{31}^{(3)}) \cdot \nabla \lambda_2^{(3)} = 0.$$

We will then prove:

PROPOSITION 5.1. The operator div from $\hat{\mathbf{V}}_h$ into H_h is one-to-one and for all $\mathbf{u}_h \in \hat{\mathbf{V}}_h$,

(5.5)
$$\|\mathbf{u}_h\|_{1,h,\Omega_r} \leq C(\theta) \|\operatorname{div}_h \mathbf{u}_h\|_{0,\Omega_r}.$$

Note that for this mesh, we are actually proving the inf-sup condition for *conforming* piecewise cubics.

Before proving Proposition 5.1, we gather some facts about the relationships among the various barycentric coordinates defined above. By writing B as a convex combination of the A_i , we easily obtain

$$\begin{split} &\mu_1 = \nu_1 \lambda_1^{(1)} = \nu_1 \lambda_2^{(2)} + \lambda_1^{(2)} = \nu_1 \lambda_3^{(3)} + \lambda_1^{(3)}, \\ &\mu_2 = \nu_2 \lambda_1^{(1)} + \lambda_2^{(1)} = \nu_2 \lambda_2^{(2)} = \nu_2 \lambda_3^{(3)} + \lambda_2^{(3)}, \\ &\mu_3 = \nu_3 \lambda_1^{(1)} + \lambda_3^{(1)} = \nu_3 \lambda_2^{(2)} + \lambda_3^{(2)} = \nu_3 \lambda_3^{(3)}. \end{split}$$

From these relations it follows that

$$\nu_3 \nabla \lambda_3^{(1)} = -\nu_1 \nabla \lambda_1^{(3)}, \quad \nu_1 \nabla \lambda_2^{(1)} = -\nu_2 \nabla \lambda_1^{(2)}, \quad \nu_2 \nabla \lambda_3^{(2)} = -\nu_3 \nabla \lambda_2^{(3)}.$$

By first expressing $\lambda_i^{(j)}$ as a linear function of μ_1 , μ_2 , μ_3 and then using the fact that $\sum_{i=1}^3 \nabla \lambda_i^{(j)} = 0$, we obtain the further relations

(5.6)
$$\nabla \lambda_2^{(1)} = (\nabla \mu_2 + \nu_2 \nabla \lambda_3^{(1)})/(1 - \nu_2),$$

(5.7)
$$\nabla \lambda_3^{(1)} = (\nabla \mu_3 + \nu_3 \nabla \lambda_2^{(1)})/(1 - \nu_3),$$

(5.8)
$$\nabla \lambda_3^{(2)} = (\nabla \mu_3 + \nu_3 \nabla \lambda_1^{(2)})/(1 - \nu_3),$$

(5.9)
$$\nabla \lambda_1^{(2)} = (\nabla \mu_1 + \nu_1 \nabla \lambda_3^{(2)}) / (1 - \nu_1),$$

(5.10)
$$\nabla \lambda_1^{(3)} = (\nabla \mu_1 + \nu_1 \nabla \lambda_2^{(3)}) / (1 - \nu_1),$$

(5.11)
$$\nabla \lambda_2^{(3)} = (\nabla \mu_2 + \nu_2 \nabla \lambda_1^{(3)}) / (1 - \nu_2).$$

Using these facts, we are now ready to prove Proposition 5.1.

Proof of Proposition 5.1. From the equation $\operatorname{div} \mathbf{u}_h = q$, we get in triangle T_1 the equations

$$\mathbf{u}_{12} \cdot \nabla \lambda_2 + \mathbf{u}_{13} \cdot \nabla \lambda_3 = q_1,$$

$$\mathbf{u}_{21} \cdot \nabla \lambda_1 = q_2,$$

$$\mathbf{u_{31}} \cdot \nabla \lambda_1 = q_3,$$

$$(5.15) 2\mathbf{u}_{12} \cdot \nabla \lambda_1 + 2\mathbf{u}_{21} \cdot \nabla \lambda_2 + \mathbf{u}_{123} \cdot \nabla \lambda_3 = q_{12},$$

$$\mathbf{u}_{123} \cdot \nabla \lambda_1 = q_{23},$$

in triangle T_2 the equations

$$\mathbf{u}_{12} \cdot \nabla \lambda_2 = q_1,$$

$$\mathbf{u}_{21} \cdot \nabla \lambda_1 + \mathbf{u}_{23} \cdot \nabla \lambda_3 = q_2,$$

$$\mathbf{u}_{32} \cdot \nabla \lambda_2 = q_3,$$

$$\mathbf{u}_{123} \cdot \nabla \lambda_2 = q_{13},$$

$$(5.21) 2\mathbf{u}_{23} \cdot \nabla \lambda_2 + 2\mathbf{u}_{32} \cdot \nabla \lambda_3 + \mathbf{u}_{123} \cdot \nabla \lambda_1 = q_{23},$$

and in triangle T_3 the equations

$$\mathbf{u}_{13} \cdot \nabla \lambda_3 = q_1,$$

$$\mathbf{u}_{23} \cdot \nabla \lambda_3 = q_2,$$

$$\mathbf{u}_{31} \cdot \nabla \lambda_1 + \mathbf{u}_{32} \cdot \nabla \lambda_2 = q_3,$$

$$\mathbf{u}_{123} \cdot \nabla \lambda_3 = q_{12},$$

$$(5.26) 2\mathbf{u}_{31} \cdot \nabla \lambda_3 + 2\mathbf{u}_{13} \cdot \nabla \lambda_1 + \mathbf{u}_{123} \cdot \nabla \lambda_1 = q_{13}.$$

Now observe that $\mathbf{u}_{21}^{(1)} = \mathbf{u}_{23}^{(3)}$ is determined from Eqs. (5.13) and (5.23), $\mathbf{u}_{31}^{(1)} = \mathbf{u}_{32}^{(2)}$ is determined form Eqs. (5.14) and (5.19), and $\mathbf{u}_{12}^{(2)} = \mathbf{u}_{13}^{(3)}$ is determined from Eqs. (5.17) and (5.22). More precisely, using the relationships among the barycentric coordinates, the functions $\mathbf{u}_{ji}^{(i)} \cdot \nabla \lambda_k^{(i)}$, $i = 1, 3, j = 1, 3, k = 1, 3, j \neq i$, may be written as linear combinations of $q_j^{(i)}$, i = 1, 3, j = 1, 3, and these $\mathbf{u}_{ji}^{(i)} = 0$ if the above $q_j^{(i)} = 0$. Then $\mathbf{u}_{ij}^{(i)} \cdot \nabla \lambda_k^{(i)}$, $i \neq j \neq k$, is determined by Eqs. (5.2), (5.3), and (5.4). We deduce from (5.6), (5.7), and (5.12) that in triangle T_1

(5.27)
$$\frac{1}{1-\nu_2}\mathbf{u}_{12} \cdot \nabla \mu_2 + \frac{1}{1-\nu_3}\mathbf{u}_{13} \cdot \nabla \mu_3 \\
= q_1 - \frac{\nu_2}{1-\nu_2}\mathbf{u}_{12} \cdot \nabla \lambda_3^{(1)} - \frac{\nu_3}{1-\nu_3}\mathbf{u}_{13} \cdot \nabla \lambda_2^{(1)}.$$

Similarly, in T_2

(5.28)
$$\frac{1}{1-\nu_3}\mathbf{u}_{23} \cdot \nabla \mu_3 + \frac{1}{1-\nu_1}\mathbf{u}_{21} \cdot \nabla \mu_1 \\ = q_2 - \frac{\nu_3}{1-\nu_3}\mathbf{u}_{23} \cdot \nabla \lambda_1^{(2)} - \frac{\nu_1}{1-\nu_1}\mathbf{u}_{21} \cdot \nabla \lambda_3^{(2)}$$

and in T_3

(5.29)
$$\frac{1}{1-\nu_1}\mathbf{u}_{31} \cdot \nabla \mu_1 + \frac{1}{1-\nu_2}\mathbf{u}_{32} \cdot \nabla \mu_2 \\ = q_3 - \frac{\nu_1}{1-\nu_1}\mathbf{u}_{31} \cdot \nabla \lambda_2^{(3)} - \frac{\nu_2}{1-\nu_2}\mathbf{u}_{32} \cdot \nabla \lambda_1^{(3)}.$$

Using the equivalences in (5.1), the above system of equations may be reduced to a linear system in the three unknowns $\mathbf{u}_{12}^{(1)} \cdot \nabla \mu_2$, $\mathbf{u}_{23}^{(2)} \cdot \nabla \mu_3$, $\mathbf{u}_{31}^{(3)} \cdot \nabla \mu_1$, with matrix

$$\begin{pmatrix} (1-\nu_2)^{-1} & (1-\nu_3)^{-1} & 0\\ 0 & (1-\nu_3)^{-1} & (1-\nu_1)^{-1}\\ (1-\nu_2)^{-1} & 0 & (1-\nu_1)^{-1} \end{pmatrix}.$$

The determinant of this matrix is $2[(1-\nu_1)(1-\nu_2)(1-\nu_3)]^{-1} \neq 0$, which allows us to determine $\mathbf{u}_{ij}^{(i)}$. Finally, we use Eqs. (5.15), (5.16), (5.20), (5.21), (5.25), and (5.26) to determine $\mathbf{u}_{123}^{(i)}$ as linear functions of the q_{ij} and q_i . Thus we see that the operator div from $\hat{\mathbf{V}}_h$ into H_h is one-to-one. Inequality (5.5) follows by the same scaling argument used in the proof of Proposition 3.1.

6. Additional Results. In checking the validity of Hypothesis H1 for a given mesh, it may be the case that there are a few triangles which do not easily fit into one of the mesh types covered in the previous sections. In that case, the following lemmas may be useful.

LEMMA 6.1. Let T_0 and T_1 be two triangles with a common side s and $\Omega_0 = T_0 \cup T_1$. Then, for all $q_h \in H_h(T_0)$, there exists $\mathbf{u}_h \in \mathbf{V}_h(\Omega_0)$ such that

$$\operatorname{div}_h \mathbf{u}_h = q_h \quad \text{in } T_0, \quad \int_s \mathbf{u}_h \, d\sigma = 0, \quad \text{and} \quad \|\mathbf{u}_h\|_{1,h,\Omega_0} \leq C(\theta) \|q_h\|_{0,T_0},$$

where C depends only on the smallest angle θ of the triangles T_0 and T_1 .

Proof.

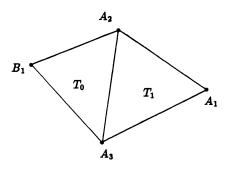


FIGURE 6.1

Using the notation of Proposition 3.1 and Figure 6.1, we see that the requirement that $\mathbf{u}^{(0)} = 0$ at the Gauss points of $\overrightarrow{B_1 A_2}$ and $\overrightarrow{B_1 A_3}$ leads to the conditions

$$\mathbf{u}_2 = \mathbf{u}_3 = -\mathbf{u}_1, \quad \mathbf{u}_{13} = -\mathbf{u}_{31} = -9\mathbf{u}_1, \quad \mathbf{u}_{12} = -\mathbf{u}_{21} = -9\mathbf{u}_1,$$

and the requirement that $\int_s \mathbf{u}_h d\sigma = 0$ on $\overrightarrow{A_3 A_2}$ can be implemented by setting

$$\mathbf{u}_{23} = 3\mathbf{u}_1 + \kappa, \quad \mathbf{u}_{32} = 3\mathbf{u}_1 - \kappa,$$

with κ to be determined. Setting $\alpha = \mathbf{u}_1 \cdot \nabla \lambda_2$, the equation $\operatorname{div}_h \mathbf{u}_h = q_h$ in T_0 leads to the equations

$$\begin{aligned} \mathbf{u}_{1} \cdot \nabla \lambda_{1} &= q_{1}/12, \\ \kappa \cdot \nabla \lambda_{3} &= q_{2} - q_{1}/2 + 6\alpha, \\ \kappa \cdot \nabla \lambda_{2} &= -q_{3} + q_{1} + 6\alpha, \\ \mathbf{u}_{123} \cdot \nabla \lambda_{2} &= q_{13} + 3q_{1} + 18\alpha, \\ \mathbf{u}_{123} \cdot \nabla \lambda_{3} &= q_{12} + 3q_{1}/2 - 18\alpha. \end{aligned}$$

It is easy to check that choosing $\alpha = -q_1/24$ and defining **u** in triangle T_1 by

$$\mathbf{u}_{i}^{(1)} = \mathbf{u}_{i}^{(0)}, \quad \mathbf{u}_{ij}^{(1)} = \mathbf{u}_{ij}^{(0)}, \quad \mathbf{u}_{123}^{(1)} = \mathbf{u}_{123}^{(0)}$$

gives the desired result.

LEMMA 6.2. If $\overline{\Omega}_s = \overline{\Omega}_r \cup T$, where T has a common side with Ω_r , and if Hypothesis H1 holds for Ω_r , then Hypothesis H1 holds for Ω_s .

Proof. Let $q_h \in H_h(\Omega_s)$. Lemma 6.1 then implies that there exists $\mathbf{u}_0 \in \mathbf{V}_h(\Omega_s)$ such that

$$\operatorname{div}_h \mathbf{u}_0 = q_h \quad \text{in } T \quad \text{and} \quad \|\mathbf{u}_0\|_{1,h,\Omega_s} \le C_1 \|q_h\|_{0,\Omega_s}.$$

Setting $q_1 = q_h - \operatorname{div}_h \mathbf{u}_0$, we observe that $\int_s \mathbf{u}_h d\sigma = 0$ implies that $q_1 \in H_h(\Omega_r)$.

Hence, there exists $\mathbf{u}_1 \in \mathbf{V}_h(\Omega_r)$ with $\operatorname{div}_h \mathbf{u}_1 = q_1$ and

$$\|\mathbf{u}_1\|_{1,h,\Omega_s} \le \|q_1\|_{0,\Omega_r}/\gamma_r \le (1+C_1\sqrt{2})\|q_h\|_{0,\Omega_s}/\gamma_r.$$

Taking $\mathbf{u}_h = \mathbf{u}_0 + \mathbf{u}_1$, we get that $\operatorname{div}_h \mathbf{u}_h = q_h$ and

$$\|\mathbf{u}_h\|_{1,h,\Omega_s} \leq (C_1 + [1 + C_1\sqrt{2}]/\gamma_r)\|q_h\|_{0,\Omega_s} = \|q_h\|_{0,\Omega_s}/\gamma_s$$

with
$$\gamma_s = \gamma_r / (1 + C_1 \sqrt{2} + C_1 \gamma_r)$$
.

Remark. Finally, since we have not found a counterexample, we conjecture that the spaces \mathbf{W}_h and Q_h form a stable Stokes pair for any triangulation of a convex polygon satisfying the minimal angle condition and containing an interior vertex.

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