

**Lectures on  
Topics In Finite Element Solution of  
Elliptic Problems**

**By  
Bertrand Mercier**

**Tata Institute of Fundamental Research  
Bombay  
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Bertrand Mercier**

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# Preface

THESE NOTES SUMMARISE a course on the finite element solution of Elliptic problems, which took place in August 1978, in Bangalore.

I would like to thank Professor Ramanathan without whom this course would not have been possible, and Dr. K. Balagangadharan who welcomed me in Bangalore.

Mr. Vijayasundaram wrote these notes and gave them a much better form than what I would have been able to.

Finally, I am grateful to all the people I met in Bangalore since they helped me to discover the smile of India and the depth of Indian civilization.

**Bertrand Mercier**

Paris, June 7, 1979.



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# Chapter 1

## Sobolev Spaces

IN THIS CHAPTER the notion of Sobolev space  $H^1(\Omega)$  is introduced. **1**  
We state the Sobolev imbedding theorem, Rellich theorem, and Trace theorem for  $H^1(\Omega)$ , without proof. For the proof of the theorems the reader is referred to ADAMS [1].

### 1.1 Notations

Let  $\Omega \subset \mathbb{R}^n$  ( $n = 1, 2$  or  $3$ ) be an open set. Let  $\Gamma$  denote the boundary of  $\Omega$ , it is assumed to be bounded and smooth. Let

$$L^2(\Omega) = \left\{ f : \int_{\Omega} |f|^2 dx < \infty \right\} \quad \text{and}$$
$$(f, g) = \int_{\Omega} fg dx$$

Then  $L^2(\Omega)$  is a Hilbert space with  $(\cdot, \cdot)$  as the scalar product.

## 1.2 Distributions

Let  $\mathcal{D}(\Omega)$  denote the space of infinitely differentiable functions with compact support in  $\Omega$ .  $\mathcal{D}(\Omega)$  is a nonempty set. If

$$f(x) = \begin{cases} \exp\left(\frac{1}{|x|^2-1}\right) & \text{if } |x| < 1 \\ 0 & \text{if } |x| \geq 1 \end{cases}$$

then  $f(x) \in \mathcal{D}(\Omega)$ ,  $\Omega = \mathbb{R}$ .

The topology chosen for  $\mathcal{D}(\Omega)$  is such that a sequence of elements  $\phi_n$  in  $\mathcal{D}(\Omega)$  converges to an element  $\phi$  belonging to  $\mathcal{D}(\Omega)$  in  $\mathcal{D}(\Omega)$  if there exists a compact set  $K$  such that

$$\text{supp } \phi_n, \text{supp } \phi \subset K$$

- 2  $D^\alpha \phi_n \rightarrow D^\alpha \phi$  uniformly for each multi-index  $\alpha = (\alpha_1, \dots, \alpha_n)$  where  $D^\alpha \phi$  stands for

$$\frac{\partial^{\alpha_1 + \dots + \alpha_n} \phi}{\partial^{\alpha_1} x_1 \dots \partial^{\alpha_n} x_n}.$$

A continuous linear functional on  $\mathcal{D}(\Omega)$  is said to be a *distribution*. The space of distributions is denoted by  $\mathcal{D}'(\Omega)$ . We use  $\langle \cdot, \cdot \rangle$  for the duality bracket between  $\mathcal{D}'(\Omega)$  and  $\mathcal{D}(\Omega)$ .

**EXAMPLE 1.** (a) A square integrable function defines a distribution: If  $f \in L^2(\Omega)$  then

$$\langle f, \phi \rangle = \int_{\Omega} f \phi \, dx \quad \text{for all } \phi \in \mathcal{D}(\Omega)$$

can be seen to be a distribution. We identify  $L^2(\Omega)$  as a space of distribution, i.e.

$$L^2(\Omega) \subset \mathcal{D}'(\Omega).$$

- (b) The dirac mass  $\delta$ , concentrated at the origin, defined by

$$\langle \delta, \phi \rangle = \phi(0) \quad \text{for all } \phi \in \mathcal{D}(\Omega)$$

defines a distribution.

**DEFINITION. Derivation of a Distribution**

If  $f$  is a smooth function and  $\phi \in \mathcal{D}(\Omega)$  then using integration by parts we obtain

$$\int_{\Omega} \frac{\partial f}{\partial x_i} \phi \, dx = - \int_{\Omega} f \frac{\partial \phi}{\partial x_i} \, dx.$$

This gives a motivation for defining the derivative of a distribution. 3

If  $T \in \mathcal{D}'(\Omega)$  and  $\alpha$  is a multi index then  $D^\alpha T \in \mathcal{D}'(\Omega)$  is defined by

$$\langle D^\alpha T, \phi \rangle = (-1)^{|\alpha|} \langle T, D^\alpha \phi \rangle \quad \forall \phi \in \mathcal{D}(\Omega).$$

If  $T_n, T \in \mathcal{D}'(\Omega)$  then we say  $T_n \rightarrow T$  in  $\mathcal{D}'(\Omega)$  if

$$\langle T_n, \phi \rangle \rightarrow \langle T, \phi \rangle \quad \text{for all } \phi \in \mathcal{D}(\Omega).$$

The derivative mapping  $D^\alpha : \mathcal{D}' \rightarrow \mathcal{D}'$  is continuous since if  $T_n \rightarrow T$  in  $\mathcal{D}'$  then

$$\begin{aligned} \langle D^\alpha T_n, \phi \rangle &= (-1)^{|\alpha|} \langle T_n, D^\alpha \phi \rangle \\ &\rightarrow (-1)^{|\alpha|} \langle T, D^\alpha \phi \rangle \\ &= \langle D^\alpha T, \phi \rangle \quad \text{for all } \phi \in \mathcal{D}(\Omega). \end{aligned}$$

**1.3 Sobolev Space**

The Sobolev space  $H^1(\Omega)$  is defined by

$$H^1(\Omega) = \left\{ v \in L^2(\Omega) : \frac{\partial v}{\partial x_i} \in L^2(\Omega), 1 \leq i \leq n \right\}$$

where the derivatives are taken in the sense of distribution.

$$f \in L^2(\Omega) \text{ need not imply } \frac{\partial f}{\partial x_i} \in L^2(\Omega).$$

**EXAMPLE 2.** Let  $\Omega = [-l, l]$

$$f(x) = \begin{cases} -1 & \text{if } x < 0 \\ 0 & \text{if } x \geq 0. \end{cases}$$

Then  $f \in L^2[-l, l]$ ; but  $df/dx = \delta$  is not given by a locally integrable function and hence not by an  $L^2$  function.

We define an inner product  $(\cdot, \cdot)_1$  in  $H^1(\Omega)$  as follows:

4

$$(u, v)_1 = (u, v) + \sum_{i=1}^n \left( \frac{\partial u}{\partial x_i}, \frac{\partial v}{\partial x_i} \right) \quad \text{for all } u, v \in H^1(\Omega).$$

Let  $\|\cdot\|_1$  be the norm associated with this inner product. Then

**LEMMA 1.**  $H^1(\Omega)$  with  $\|\cdot\|_1$  is a Hilbert spaces.

*Proof.* Let  $u_j$  be a Cauchy sequence in  $H^1(\Omega)$ . This imply

$$\{u_j\}, \left\{ \frac{\partial u_j}{\partial x_i} \right\} \quad i = 1, 2, \dots, n$$

are Cauchy in  $L^2$ . Hence there exists  $v, v_i \in L^2(\Omega)$   $1 \leq i \leq n$  such that

$$\begin{aligned} u_j &\rightarrow v \text{ in } L^2(\Omega), \\ \frac{\partial u_j}{\partial x_i} &\rightarrow v_i \text{ in } L^2(\Omega), \quad 1 \leq i \leq n, \end{aligned}$$

For any  $\phi \in \mathcal{D}(\Omega)$ ,

$$\left\langle \frac{\partial u_j}{\partial x_i}, \phi \right\rangle = - \left\langle u_j, \frac{\partial \phi}{\partial x_i} \right\rangle \rightarrow - \left\langle u, \frac{\partial \phi}{\partial x_i} \right\rangle = \left\langle \frac{\partial u}{\partial x_i}, \phi \right\rangle.$$

But

$$\left\langle \frac{\partial u_j}{\partial x_i}, \phi \right\rangle \rightarrow \langle v_i, \phi \rangle.$$

Hence

$$v_i = \frac{\partial u}{\partial x_i}.$$

Thus

$$\begin{aligned} u_j &\rightarrow u \text{ in } L^2(\Omega) \\ \frac{\partial u_j}{\partial x_i} &\rightarrow \frac{\partial u}{\partial x_i} \text{ in } L^2(\Omega). \end{aligned}$$

5 This proves  $u_j \rightarrow u$  in  $H^1(\Omega)$ . □

## 1.4 Negative Properties of $H^1(\Omega)$

- (a) The functions in  $H^1(\Omega)$  need not be continuous except in the case  $n = 1$ .

**EXAMPLE 3.** Let

$$\begin{aligned}\Omega &= \{(x, y) \in \mathbb{R}^2 : x^2 + y^2 < r_0^2\}, r_0 < 1. \\ f(r) &= (\log 1/r)^k, k < 1/2 \quad \text{where} \\ r &= (x^2 + y^2)^{1/2}.\end{aligned}$$

Then  $f \in H^1(\Omega)$  but  $f$  is not continuous at the origin.

In the case  $n = 1$ , if  $u \in H^1(\Omega)$ ,  $\Omega \subset \mathbb{R}^1$  then  $u$  can be shown to be continuous using the formula

$$u(y) - u(x) = \int_x^y \frac{du}{dx}(s) ds$$

where  $du/ds$  denotes the distributional derivative of  $u$ .

- (b)  $\mathcal{D}(\Omega)$  is not dense in  $H^1(\Omega)$ . To see this let  $u \in (\mathcal{D}(\Omega))^\perp$  in  $H^1$  and  $\phi \in \mathcal{D}(\Omega)$ . We have

$$(u, \phi)_1 = (u, \phi) + \sum_{i=1}^n \left\langle \frac{\partial u}{\partial x_i}, \frac{\partial \phi}{\partial x_i} \right\rangle = 0$$

$$\text{i.e.} \quad \langle u, \phi \rangle + \sum_{i=1}^n \left\langle -\frac{\partial^2 u}{\partial x_i^2}, \phi \right\rangle = 0$$

Thus

$$\langle -\Delta u + u, \phi \rangle = 0 \quad \text{for all } \phi \in \mathcal{D}(\Omega).$$

Hence

$$-\Delta u + u = 0 \quad \text{in } \mathcal{D}'(\Omega).$$

Let  $\Omega = \{x \in \mathbb{R}^n : |x| < 1\}$ ,

6

$$\begin{aligned} u(x) &= e^{r \cdot x} \quad \text{where } r \in \mathbb{R}^n, \\ \Delta u(x) &= |r|^2 e^{r \cdot x} = |r|^2 u. \\ &= u \quad \text{if } |r| = 1. \end{aligned}$$

Thus when  $n = 1$ ,  $u$  with  $r = \pm 1$  belongs to  $\mathcal{D}(\Omega)^\perp$ , when  $n > 1$  there are infinitely many  $r$ 's ( $r \in S^{n-1}$ ) such that  $u \in (\mathcal{D}(\Omega)^\perp)^\perp$ . Moreover these functions for different  $r$ 's are linearly independent. Therefore

$$\begin{aligned} \text{dimension } (\mathcal{D}(\Omega)^\perp) &\geq 2 \quad \text{if } n = 1 \\ \text{dimension } (\mathcal{D}(\Omega)^\perp) &= \infty \quad \text{if } n > 1. \end{aligned}$$

This proves the claim (b).

We shall define  $H_0^1(\Omega)$  as the closure of  $\mathcal{D}(\Omega)$  in  $H^1(\Omega)$ . We have the following inclusions

$$\mathcal{D}(\Omega) \underset{\text{dense}}{\subset} H_0^1(\Omega) \subset H^1(\Omega) \underset{\text{dense}}{\subset} L^2(\Omega).$$

## 1.5 Trace Theorem

Let  $\Omega$  be a bounded open subset of  $\mathbb{R}^n$  with a Lipschitz continuous boundary  $\Gamma$  : *i.e.* there exists finite number of local charts  $a_j$ ,  $1 \leq j \leq J$  from  $\{y' \in \mathbb{R}^{n-1} : |y'| < \alpha\}$  into  $\mathbb{R}^n$  and a number  $\beta > 0$  such that

$$\begin{aligned} \Gamma &= \bigcup_{j=1}^J \{(y', y_n) : y_n = a_j(y'), |y'| < \alpha\}, \\ \{(y', y_n) : a_j(y') < y_n < a_j(y') + \beta, |y'| < \alpha\} &\subset \Omega, 1 \leq j \leq J, \\ \{(y', y_n) : a_j(y') - \beta < y_n < a_j(y'), |y'| < \alpha\} &\subset C\bar{\Omega}, 1 \leq j \leq J. \end{aligned}$$

7 It can be proved that  $C^\infty(\bar{\Omega})$  is dense in  $H^1(\Omega)$ . If  $f \in C^\infty(\bar{\Omega})$  we define the trace of  $f$ , namely  $\gamma f$ , by

$$\gamma f = f|_\Gamma. \quad \text{Note } \gamma f \in L^2(\Gamma) \quad \text{if } f \in C^\infty(\bar{\Omega})$$



$\gamma : C^\infty(\overline{\Omega}) \rightarrow L^2(\Gamma)$  is continuous and linear with norm  $\|\gamma u\|_{L^2(\Gamma)} \leq C \|u\|_1$ . Hence this can be extended as continuous linear map from  $H^1(\Omega)$  to  $L^2(\Gamma)$ .

$H^1_0(\Omega)$  is characterised by

**THEOREM 2.**

$$H^1_0(\Omega) = \{v \in H^1(\Omega) : \gamma v = 0\}$$

## 1.6 Dual Spaces of $H^1(\Omega)$ and $H^1_0(\Omega)$

The mapping

$$I : H^1(\Omega) \rightarrow (L^2(\Omega))^{n+1} \text{ defined by}$$

$$I(v) = \left( v, \frac{\partial v}{\partial x_1}, \dots, \frac{\partial v}{\partial x_n} \right)$$

is easily seen to be an isometric isomorphism of  $H^1(\Omega)$  into subspace of  $(L^2(\Omega))^{n+1}$ . If  $f \in (H^1(\Omega))'$  then  $F : I(H^1(\Omega)) \rightarrow \mathbb{R}$  with  $F(Iu) = f(u)$  is a continuous linear functional on  $I(H^1(\Omega))$ . Hence by Hahn Banach theorem  $F$  can be extended to  $(L^2(\Omega))^{n+1}$ . Therefore, there exists  $(v, v_1, \dots, v_n) \in (L^2(\Omega))^{n+1}$  such that

$$f(u) = F(Iu) = (v, u) + \sum_{i=1}^n (v_i, \partial u / \partial x_i).$$

This representation is not unique since  $F$  cannot be extended uniquely to  $(L^2(\Omega))^{n+1}$ . For all  $\phi \in \mathcal{D}(\Omega)$  we have

$$f(\phi) = (v, \phi) - \sum_{i=1}^n \left\langle \frac{\partial v_i}{\partial x_i}, \phi \right\rangle,$$

Thus

$$f|_{\mathcal{D}(\Omega)} = v - \sum_{i=1}^n \frac{\partial v_i}{\partial x_i}.$$

Conversely if  $T \in \mathcal{D}'(\Omega)$  is given by

$$T = v - \sum_{i=1}^n \frac{\partial v_i}{\partial x_i},$$

where  $v, v_i \in L^2(\Omega)$ ,  $1 \leq i \leq n$  then  $T$  can be extended as a continuous linear functional on  $H^1(\Omega)$  by the prescription

$$T(u) = (v, u) + \sum_{i=1}^n \left( v_i, \frac{\partial u}{\partial x_i} \right) \quad \text{for all } u \in H^1(\Omega).$$

The extension of  $T$  to  $H^1(\Omega)$  need not be unique. But we will prove that the extension of  $T$  to  $H^1_0(\Omega)$  is unique. Let  $\tilde{T} \in (H^1_0(\Omega))'$  be such that  $\tilde{T}|_{\mathcal{D}(\Omega)} = T$ .

Let  $u \in H^1_0(\Omega)$ . Then there exists  $u_m \in \mathcal{D}(\Omega)$  such that  $u_m \rightarrow u$  in  $H^1$ . Now

$$\begin{aligned} \tilde{T}(u) &= \tilde{T}(\lim_{\text{in } H^1} u_m) = \lim_m \tilde{T}(u_m) \\ &= \lim_m T(u_m) \\ &= \lim_m \left[ (v, u_m) + \sum_{i=1}^n \left( v_i, \frac{\partial u_m}{\partial x_i} \right) \right] \\ &= (v, u) + \sum_{i=1}^n \left( v_i, \frac{\partial u}{\partial x_i} \right) \end{aligned}$$

Thus

$$\tilde{T} = T \text{ on } H^1_0(\Omega).$$

- 9 Hence we identify  $(H^1_0(\Omega))'$  with a space of distribution and we denote it by  $H^{-1}(\Omega)$ . That is

$$H^{-1}(\Omega) = \left\{ v - \sum_{i=1}^n \frac{\partial v_i}{\partial x_i} : (v, v_1, \dots, v_n) \in (L^2(\Omega))^{n+1} \right\} \subset \mathcal{D}'(\Omega).$$

**EXERCISE 1.** Show that  $\partial/\partial x_i : L^2(\Omega) \rightarrow H^{-1}(\Omega)$  is continuous.

## 1.7 Positive Properties of $H^1_0(\Omega)$ and $H^1(\Omega)$ .

**THEOREM 3. (Poincaré's Inequality).** Let  $\Omega$  be an open bounded subset of  $\mathbb{R}^n$ . Then there exists a constant  $C(\Omega)$  such that

$$\int_{\Omega} v^2 dx \leq C(\Omega) \int_{\Omega} |\nabla v|^2 dx \quad \text{for all } v \in H^1_0(\Omega).$$

*Proof.* We shall prove the inequality for the functions in  $\mathcal{D}(\Omega)$  and use the density of  $\mathcal{D}(\Omega)$  in  $H^1_0(\Omega)$ .

Since  $\Omega$  is bounded, we have

$$\Omega \subset [a_1, b_1] \times \dots \times [a_n, b_n].$$

For any  $u(x) \in \mathcal{D}(\Omega)$ , we have

$$u(x) = \int_{a_i}^{x_i} \frac{\partial u}{\partial x_i}(x_1, \dots, x_{i-1}, t, x_{i+1}, \dots, x_n) dt.$$

Thus

$$|u(x)| \leq (b_i - a_i)^{1/2} \left( \int_{a_i}^{b_i} \left| \frac{\partial u}{\partial x_i} \right|^2 dt \right)^{1/2}$$

Squaring both sides and integrating we obtain

$$\int_{a_1}^{b_1} \dots \int_{a_n}^{b_n} |u(x)|^2 dx \leq (b_i - a_i)^2 \int_{a_1}^{b_1} \dots \int_{a_n}^{b_n} \left| \frac{\partial u}{\partial x_i} \right|^2 dx.$$

Thus

$$\int_{\Omega} |u(x)|^2 dx \leq C(\Omega) \int_{\Omega} |\nabla u|^2 dx$$

10

where

$$C(\Omega) = (b_1 - a_1)^2 + \dots + (b_n - a_n)^2.$$

Let  $v \in H^1_0(\Omega)$ . Then there exists  $u_n \in \mathcal{D}(\Omega)$  such that  $u_n \rightarrow v$  in  $H^1$ , which implies  $u_n \rightarrow v$  in  $L^2$  and  $\frac{\partial u_n}{\partial x_i} \rightarrow \frac{\partial v}{\partial x_i}$  in  $L^2$ . Using this and the inequality for smooth functions we arrive at the result.  $\square$

**REMARK 1.** The theorem is not true for functions in  $H^1(\Omega)$ . For example a nonzero constant function belongs to  $H^1(\Omega)$  but does not satisfy the above inequality.

We state Sobolev imbedding theorem and Rellich's theorem, which have many important applications.

**SOBOLEV IMBEDDING THEOREM 4.** *If  $\Omega$  is an open bounded set having a Lipschitz continuous boundary then we have the imbedding*

$$H^1(\Omega) \hookrightarrow L^p(\Omega),$$

*$p < q$  or  $p \leq q$  according as  $n = 2$  or  $n > 2$  where*

$$1/q = 1/2 - 1/n.$$

**RELLICH'S THEOREM 5.** *The above imbedding is compact for  $p < q$ .*

## Chapter 2

# Abstract Variational Problems and Examples

IN SECTION 1 OF this chapter, we give a variational formulation of the Dirichlet problem. In section 2 we prove the existence and uniqueness results for the abstract variational problem. In the remaining sections, we deal with the Neumann problem, Elasticity problem, Stokes problem and Mixed problem and their variational formulations. 11

### 2.1 Dirichlet Problem

The Dirichlet problem is to find  $u$  such that

$$-\Delta u = f \quad \text{in } \Omega \quad (2.1)$$

$$u = 0 \quad \text{on } \Gamma \quad (2.2)$$

where  $\Omega \subset \mathbb{R}^n$  is a bounded open set with smooth boundary  $\Gamma$  and  $f$  is a given function.

Multiplying equation (2.1) by a smooth function  $v$  which vanishes on  $\Gamma$  and integrating, we obtain

$$\int_{\Omega} -\Delta u \cdot v \, dx = \int_{\Omega} f \cdot v \, dx \quad (2.3)$$

Formally, using integration by parts and the fact that  $v = 0$  on  $\Gamma$ , we see that

$$\int_{\Omega} \nabla v \cdot \nabla u \, dx = \int_{\Gamma} \frac{\partial u}{\partial n} v \, d\Gamma - \int_{\Omega} \Delta u \, v \, dx = \int_{\Omega} -\Delta u \cdot v \, dx \quad (2.4)$$

Equations (2.3) and (2.4) give

$$\int_{\Omega} \nabla u \cdot \nabla v \, dx = \int_{\Omega} f \, v \, dx.$$

Now, setting

$$a(u, v) = \int_{\Omega} \nabla u \cdot \nabla v \, dx$$

12 and

$$L(v) = \int_{\Omega} f \, v \, dx,$$

problem (2.1) can be formulated thus:

Find  $u \in V$  such that

$$a(u, v) = L(v) \quad \text{for all } v \in V, \quad (2.5)$$

where  $V$  has to be chosen suitably.

Since  $a(\cdot, \cdot)$  is symmetric, (2.5) is Euler's condition for the minimization problem

$$J(u) = \inf_{v \in V} J(v), \quad (2.6)$$

where

$$J(v) = 1/2 a(v, v) - L(v).$$

If  $u$  is the solution of the minimization problem then it can be shown that

$$(J'(u), v) = 0 \quad \text{for all } v \in V, \quad (2.7)$$

where  $(J'(u), v)$  is the Gateaux derivative of  $J$  in the direction  $v$ .

When  $a(\cdot, \cdot)$  is symmetric one can show that problems (2.5) and (2.6) are equivalent.

To have  $J(v)$  finite, we want our space  $V$  to be such that  $\nabla v \in L^2(\Omega)$ ,  $f \in L^2(\Omega)$  for all  $v \in V$ . The largest space satisfying the above conditions and (2.2) is  $H_0^1(\Omega)$  and hence we choose  $V$  to be  $H_0^1(\Omega)$ .

If  $u$  is a solution of (2.5) then

$$\int_{\Omega} \nabla u \cdot \nabla v \, dx = \int_{\Omega} f v \, dx \quad \text{for all } v \in \mathcal{D}(\Omega) \subset H_0^1(\Omega).$$

This implies that

13

$$\langle -\Delta u, v \rangle = \langle f, v \rangle \quad \text{for all } v \in \mathcal{D}(\Omega);$$

so

$$-\Delta u = f \quad \text{in } \mathcal{D}'. \quad (2.8)$$

Conversely, if  $u \in H_0^1(\Omega)$  satisfies (2.8), retracing the above steps we obtain

$$a(u, v) = f(v) \quad \text{for all } v \in \mathcal{D}(\Omega). \quad (2.9)$$

Since  $\mathcal{D}(\Omega)$  is dense in  $H_0^1(\Omega)$ , (2.9) holds for all  $v \in V = H_0^1(\Omega)$ . Thus  $u$  is the solution of (2.5).

## 2.2 Abstract Variational Problem.

We now prove the existence and uniqueness theorem for the abstract variational problem.

**THEOREM 1.** *Let  $V$  be a Hilbert space and  $a(\cdot, \cdot) : V \times V \rightarrow \mathbb{R}$  be continuous and bilinear. Further, assume that  $a(\cdot, \cdot)$  is **coercive**: there exists  $\alpha > 0$  such that  $a(v, v) \geq \alpha \|v\|_V^2$  for all  $v \in V$ . Let  $L$  be a continuous linear functional on  $V$ . Then the problem:*

*To find  $u \in V$  such that*

$$a(u, v) = L(v), \quad \text{for all } v \in V \quad (2.10)$$

*has a unique solution.*

*Proof.* (i) **Uniqueness.** Let  $u_1, u_2 \in V$  be two solutions of (2.10). Therefore

$$\begin{aligned} a(u_1, v) &= L(v), \\ a(u_2, v) &= L(v), \quad \text{for all } v \in V. \end{aligned}$$

14 Subtracting one from the other, taking  $v = u_2 - u_1$  and using  $V$ -coercivity of  $a(\cdot, \cdot)$ , we obtain

$$\alpha \|u_1 - u_2\|_V^2 \leq a(u_1 - u_2, u_1 - u_2) = 0,$$

Thus  $u_1 = u_2$ .

(ii) **Existence when  $a(\cdot, \cdot)$  is symmetric.** Since  $a(\cdot, \cdot)$  is symmetric, the bilinear form  $a(u, v)$  is a scalar product on  $V$  and the associated norm  $a(v, v)^{1/2}$  is equivalent to the norm in  $V$ . Hence, by the Riesz representation theorem there exists  $\sigma \in V$  such that

$$a(\sigma, v) = L(v) \quad \text{for all } v \in V.$$

Hence the theorem is true in the symmetric case.

(iii) **Existence in the general case.** Let  $w \in V$ . The function  $L_w : V \rightarrow \mathbb{R}$  defined by

$$L_w(v) = (w, v) - \rho(a(w, v) - L(v))$$

is linear and continuous. Hence by the Riesz representation theorem there exists a  $u \in V$  such that

$$L_w(v) = (u, v).$$

Let  $T : V \rightarrow V$  be defined by

$$Tw = u$$

where  $u$  is the solution of the equation

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



15 we will show that  $T$  is a contraction mapping. Hence  $T$  has a unique fixed point which will be the solution of (2.10).

Let

$$u_1 = Tw_1, u_2 = Tw_2.$$

Thus

$$(u_1 - u_2, v) = (w_1 - w_2, v) - \rho a(w_1 - w_2, v) \quad \forall v \in V \quad (2.11)$$

Let  $A : V \rightarrow V$ , where  $Au$  is the unique solution of

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

which exists by the Riesz representation theorem.

$$\|Au\| = \sup_{v \in V} \frac{|(Au, v)|}{\|v\|} = \sup_{v \in V} \frac{|a(u, v)|}{\|v\|} \leq M \|u\|,$$

where  $|a(u, v)| \leq M \|u\| \|v\|$ . So  $A$  is continuous. Equation (2.11) can be written as

$$(u_1 - u_2, v) = (w_1 - w_2, v) - \rho(A(w_1 - w_2), v) \quad \text{for all } v \in V,$$

which implies that

$$u_1 - u_2 = (w_1 - w_2) - \rho A(w_1 - w_2).$$

So

$$\begin{aligned} \|u_1 - u_2\|^2 &= \|w_1 - w_2\|^2 - 2\rho(A(w_1 - w_2), w_1 - w_2) \\ &\quad + \rho^2 \|A(w_1 - w_2)\|^2 \\ &\leq \|w_1 - w_2\|^2 - 2\rho a(w_1 - w_2, w_1 - w_2) \\ &\quad + \rho^2 M^2 \|w_1 - w_2\|^2, \quad (\text{using the continuity of } A) \\ &\leq \|w_1 - w_2\|^2 - 2\rho\alpha \|w_1 - w_2\|^2 + \rho^2 M^2 \|w_1 - w_2\|^2, \end{aligned}$$

since

$$a(w_1 - w_2, w_1 - w_2) \geq \alpha \|w_1 - w_2\|^2.$$

So

$$\|u_1 - u_2\|^2 \leq (1 - 2\rho\alpha + \rho^2 M^2) \|w_1 - w_2\|^2.$$

That is,

$$\|Tw_1 - Tw_2\| \leq \sqrt{(1 - 2\rho\alpha + \rho^2 M^2)} \|w_1 - w_2\|.$$

Choosing  $\rho$  in  $]0, \alpha/2M[$ , we obtain that  $T$  is a contraction.  $\square$

This proves the theorem.

**REMARK 1.** *This theorem also gives an algorithm to find the solution of equation (2.10). Let  $u^0 \in V$  be given. Let  $u^{n+1} = Tu^n$ . Then  $u^n \rightarrow w_0$ , which is the fixed point of  $T$ , and also the solution of (2.10).*

### 2.3 Neumann's problem.

Neumann's problem is to find an  $u$  such that

$$-\Delta u + cu = f \quad \text{in } \Omega, \quad (2.12)$$

$$\frac{\partial u}{\partial n} = g \quad \text{on } \Gamma. \quad (2.13)$$

We now do the calculations formally to find out the bilinear form  $a(\cdot, \cdot)$ , the linear functional  $L(\cdot)$  and the space  $V$ .

For smooth  $v$ , (2.12) implies

$$\int_{\Omega} (-\Delta u + cu)v \, dx = \int_{\Omega} fv \, dx. \quad (2.14)$$

17 From Green's formula,

$$\int_{\Omega} \nabla u \cdot \nabla v \, dx = \int_{\Gamma} \frac{\partial u}{\partial n} v \, d\Gamma - \int_{\Omega} v \Delta u \, dx,$$

and by (2.14) we obtain

$$\int_{\Omega} (\nabla u \cdot \nabla v + cuv) \, dx = \int_{\Omega} fv \, dx + \int_{\Gamma} \frac{\partial u}{\partial n} v \, d\Gamma = \int_{\Omega} fv \, dx + \int_{\Gamma} gv \, d\Gamma,$$

since  $\frac{\partial u}{\partial n} = g$  on  $\Gamma$ , by (2.13). This suggests the definitions:

$$a(u, v) = \int_{\Omega} (\nabla u \cdot \nabla v + cuv) dx \quad (2.15)$$

$$L(v) = \int_{\Omega} fv dx + \int_{\Gamma} gv d\Gamma, \quad (2.16)$$

$$V = H^1(\Omega), \quad (2.17)$$

where  $f \in L^2(\Omega)$  and  $g \in L^2(\Gamma)$ .

Clearly  $a(u, v)$  is bilinear, continuous and symmetric.

$$\begin{aligned} a(v, v) &= \int_{\Omega} ((\nabla v)^2 + cv^2) dx \\ &\geq \min\{1, c\} \|v\|_1^2, \end{aligned}$$

which shows  $a(\cdot, \cdot)$  is  $H^1(\Omega)$ -coercive.

$L(\cdot)$  is a continuous linear functional on  $H^1(\Omega)$ . Hence by the theorem there exists a unique  $u \in V = H^1(\Omega)$  such that

$$\int_{\Omega} (\nabla u \cdot \nabla v + cuv) dx = \int_{\Omega} fv dx + \int_{\Gamma} gv d\Gamma \quad \text{for all } v \in H^1(\Omega) \quad (2.18)$$

From (2.18) we obtain that for all  $v \in \mathcal{D}(\Omega)$ ,

18

$$\langle -\Delta u + cu, v \rangle = \langle f, v \rangle.$$

Hence

$$-\Delta u + cu = f \quad \text{in } \mathcal{D}'(\Omega) \quad (2.19)$$

To find the boundary condition we use Green's formula:

$$\int_{\Omega} \nabla u \cdot \nabla v dx = \int_{\Omega} -\Delta u \cdot v dx + \int_{\Gamma} \frac{\partial u}{\partial n} v d\Gamma,$$

which holds for all  $u \in H^2(\Omega)$  and for all  $v \in H^1(\Omega)$ .

Assuming that our solution  $u \in H^2(\Omega)$ , from (2.19) we have

$$\int_{\Omega} (-\Delta u + cu)v = \int_{\Omega} fv \quad \text{for all } v \in H^1(\Omega).$$

Using Green's formula we obtain

$$\int_{\Omega} (\nabla u \nabla v + cuv) dx = \int_{\Gamma} \frac{\partial u}{\partial n} v dx + \int_{\Omega} fv dx.$$

This, together with (2.18), implies

$$\int_{\Gamma} \left( g - \frac{\partial u}{\partial n} \right) v d\Gamma = 0 \quad \text{for all } v \in H^1(\Omega).$$

Hence we get the desired boundary condition

$$\frac{\partial u}{\partial n} = g \quad \text{on } \Gamma.$$

If  $u \in H^2(\Omega)$ , these are still valid in "some sense" which is given in LIONS–MAGENES [29].

**REMARK 2.** *Even when  $g = 0$  we cannot take the space*

$$V_1 = \left\{ v \in H^1(\Omega) : \frac{\partial v}{\partial n} = 0 \quad \text{on } \Gamma \right\}$$

- 19** *to be the basic space  $V$ , since  $V_1$  is not closed. In the Neumann problem 2.3, we obtain the boundary condition from Green's formula. In the case of Dirichlet problem 2.1, we impose the boundary condition in the space itself.*

**REGULARITY THEOREM (FOR DIRICHLET PROBLEM) 2.** *If  $\Gamma$  is  $C^2$  or  $\Omega$  is a convex polygon and  $f \in L^2(\Omega)$ , then the solution  $u$  of the Dirichlet problem (2.1), (2.2) belongs to  $H^2(\Omega)$ .*

**REGULARITY THEOREM (FOR THE NEUMANN PROBLEM) 3.** *If  $\Gamma$  is  $C^2$  or  $\Omega$  is a convex polygon,  $f \in L^2(\Omega)$  and  $g$  belongs to a space finer than  $L^2(\Omega)$  (for example  $g \in H^1(\Gamma)$ ), then the solution  $u$  of the Neumann problem (2.12), (2.13) belongs to  $H^2(\Omega)$ .*

For a proof of these theorems the reader is referred to NECAS [33].

## 2.4 Mixed Problem.

In Sections 2.1 and 2.3 we found the variational formulation from the partial differential equation. In the general case it is difficult to formulate the variational problem from the p.d.e. In fact a general p.d.e. need not give rise to a variational problem. So in this section, we will take a general variational problem and find out the p.d.e. satisfied by its solution.

Let  $\Omega$  be a bounded open set with boundary  $\Gamma$ . Let  $\Gamma = \Gamma_0 \cup \Gamma_1$  where  $\Gamma_0$  and  $\Gamma_1$  are disjoint. Let

$$V = \{v \in H^1(\Omega) : v = 0 \text{ on } \Gamma_0\}. \quad (2.20)$$

It is easy to see that  $V$  is closed and hence a Hilbert space with  $\|\cdot\|_1$  norm

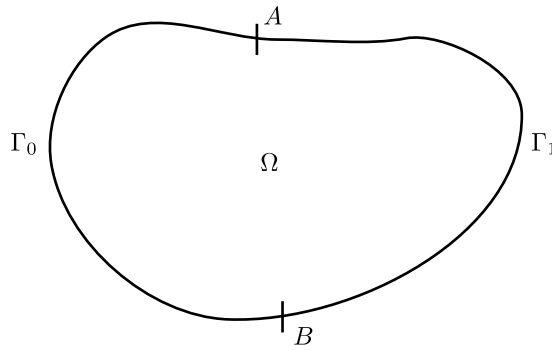


Figure 2.1:

We will use summation convention here afterwards. Let

$$a(u, v) = \int_{\Omega} \left( a_{ij}(x) \frac{\partial u}{\partial x_i} \frac{\partial v}{\partial x_j} + a_0 uv \right) dx, \quad (2.21)$$

$$L(v) = \int_{\Omega} f v dx + \int_{\Gamma} g v d\Gamma, \quad (2.22)$$

where  $a_0 > 0$ ,  $a_{ij}$  are smooth and there exists two constants  $\alpha_0$  and  $\alpha_1$  such that

$$\alpha_1 \xi_i \xi_i \geq a_{ij}(x) \xi_i \xi_j \geq \alpha_0 \xi_i \xi_i \quad \text{for all } x \in \Omega, \xi \in \mathbb{R}^n \quad (2.23)$$

i.e. the quadratic form  $a_{ij}(x) \xi_i \xi_j$  is uniformly continuous and uniformly positive definite.

Inequality (2.23) implies that the bilinear form  $a(\cdot, \cdot)$  is continuous and  $V$ -coercive. Formally we have

$$a(u, v) = \int_{\Omega} \left[ -\frac{\partial}{\partial x_j} \left( a_{ij} \frac{\partial u}{\partial x_i} \right) v + a_0 uv \right] dx + \int_{\Gamma} a_{ij} \frac{\partial u}{\partial x_i} n_j v d\Gamma \quad (2.24)$$

21 Let

$$\frac{\partial u}{\partial \nu_A} = a_{ij} \frac{\partial u}{\partial x_i} n_j,$$

and

$$Au = -\frac{\partial}{\partial x_j} \left( a_{ij} \frac{\partial u}{\partial x_i} \right).$$

If  $v \in \mathcal{D}(\Omega)$  then the equation

$$a(u, v) = L(v) \quad (2.25)$$

becomes

$$\langle Au, v \rangle = \langle f, v \rangle.$$

Therefore

$$Au = f \quad \text{in } \mathcal{D}'(\Omega). \quad (2.26)$$

Now for all  $v \in V$ , we have

$$\begin{aligned} a(u, v) &= \int_{\Omega} Au \cdot v + \int_{\Gamma} \frac{\partial u}{\partial \nu_A} v d\Gamma \\ &= \int_{\Omega} Au \cdot v + \int_{\Gamma_1} \frac{\partial u}{\partial \nu_A} v d\Gamma \\ L(v) &= \int_{\Omega} f v dx + \int_{\Gamma_1} g v d\Gamma. \end{aligned}$$

Equations (2.25) and (2.26) imply, for all  $v \in V$ ,

$$\int_{\Gamma_1} \frac{\partial u}{\partial \nu_A} v \, d\Gamma = \int_{\Gamma_1} g v \, d\Gamma.$$

From this we obtain formally

$$\frac{\partial u}{\partial \nu_A} = g \quad \text{on } \Gamma_1. \quad (2.27)$$

Thus the boundary value problem corresponding to the variational problem

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

with  $a(\cdot, \cdot)$ ,  $L(\cdot)$  and  $V$  given by the equations (2.20) - (2.22) is

$$\begin{aligned} Au &= f \quad \text{in } \Omega, \\ \frac{\partial u}{\partial \nu_A} &= g \quad \text{on } \Gamma_1, \\ u &= 0 \quad \text{on } \Gamma_0. \end{aligned} \quad (2.28)$$

**REMARK 3.** Even when  $f$  and  $g$  are smooth the solution  $u$  of the problem (2.28) may not be in  $H^2(\Omega)$ . In general, we will have a singularity at the transition points  $A, B$  on  $\Gamma$ . But if  $\Gamma_0$  and  $\Gamma_1$  make a corner then the solution  $u$  may be in  $H^2(\Omega)$  provided that the boundary functions  $f, g$  satisfy some compatibility conditions. For regularity theorems the reader is referred to an article by PIERRE GIRSVARD [22].

**EXERCISE 1. Transmission Problem.** Let  $\Omega, \Omega_1, \Omega_2$  be open sets such that  $\Omega = \Omega_1 \cup \Omega_2 \cup S$  where  $\Omega_1$  and  $\Omega_2$  are disjoint subsets of  $\Omega$  and  $S$  is the interface between them. Let

$$\begin{aligned} a(u, v) &= \sum_{i=1}^2 \int_{\Omega_i} a_i \nabla u \cdot \nabla v \, dx, \\ L(v) &= \int_{\Omega} f v \, dx, \end{aligned}$$

where  $a_i > 0, i = 1, 2$ , and  $f \in L^2(\Omega)$ . If  $u$  is the solution of the problem 23

$$a(u, v) = L(v) \quad \text{for all } v \in H_0^1(\Omega),$$

and

$$u_i = u|_{\Omega_i}, f_i = f|_{\Omega_i}$$

then show that

$$-a_i \Delta u_i = f_i \quad \text{on } \Omega_i, i = 1, 2;$$

$$u_1 = u_2 \quad \text{on } S,$$

$$a_1 \frac{\partial u_1}{\partial n} = a_2 \frac{\partial u_2}{\partial n} \quad \text{on } S.$$

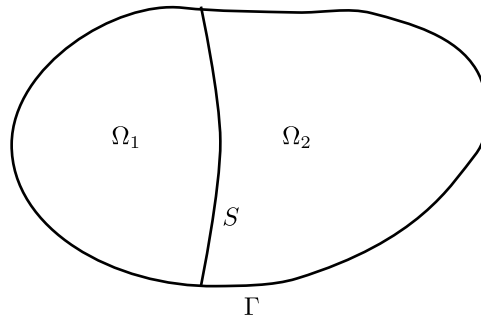


Figure 2.2:

**EXERCISE 2. Fourier Condition.** Let

$$V = H^1(\Omega),$$

$$a(u, v) = \int_{\Omega} \nabla u \cdot \nabla v \, dx + \int_{\Gamma} uv \, d\Gamma,$$

$$L(v) = \int_{\Omega} fv \, dx + \int_{\Gamma} gv \, d\Gamma,$$

What is the boundary value problem associated with this? Interpret the problem.



## 2.5 Elasticity Problem.

- (a) **3-DIMENSIONAL CASE.** Let  $\Omega \subset \mathbb{R}^3$  be a bounded, connected open set. Let  $\Gamma$  be the boundary of  $\Omega$  and let  $\Gamma$  be split into two parts  $\Gamma_0$  and  $\Gamma_1$ . Let  $\Omega$  be occupied by an elastic medium, which we assume to be continuous. Let the elastic material be fixed along  $\Gamma_0$ . Let  $(f_i)$  be the body force acting in  $\Omega$  and  $(g_i)$  be the pressure load acting along  $\Gamma_1$ . Let  $(u_i(x))$  denote the displacement at  $x$ . 24

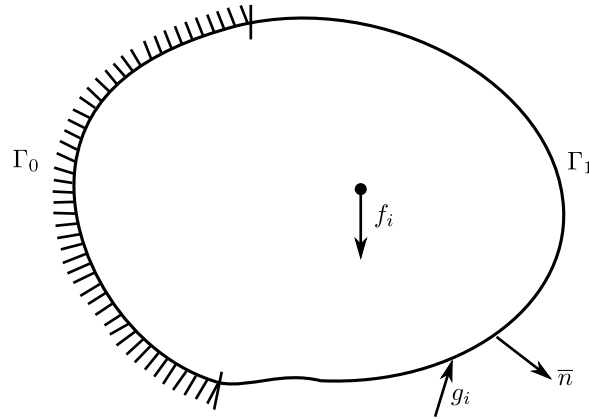


Figure 2.3:

In linear elasticity the stress-strain relation is

$$\sigma_{ij}(u) = \lambda(\operatorname{div} u)\delta_{ij} + 2\mu\varepsilon_{ij}(u), \quad \varepsilon_{ij}(u) = \frac{1}{2} \left( \frac{\partial u_i}{\partial x_j} + \frac{\partial u_j}{\partial x_i} \right), \quad (2.29)$$

where  $\sigma_{ij}$  and  $\varepsilon_{ij}$  denote the components of the stress and strain tensors respectively.

The problem is to find  $\sigma_{ij}$  and  $u_i$ , given  $(f_i)$  in  $\Omega$ ,  $(g_i)$  on  $\Gamma_1$  and  $(u_i) = 0$  on  $\Gamma_0$ .

The equations of equilibrium are

$$\frac{\partial}{\partial x_j} \sigma_{ij} + f_i = 0 \quad \text{in } \Omega, \quad (2.30)$$

$$\sigma_{ij}n_j = g_i \quad \text{on } \Gamma_1, \quad i = 1, 2, 3 \quad (2.30b)$$

$$u_i = 0 \quad \text{on } \Gamma_0. \quad (2.30c)$$

25 We have used the summation convention in the above equations. We choose

$$v = \{v \in (H^1(\Omega))^3 : v = 0 \quad \text{on } \Gamma_0\}, \quad (2.31)$$

$$a(u, v) = \int_{\Omega} \sigma_{ij}(u) \varepsilon_{ij}(v) dx, \quad (2.32)$$

$$L(v) = \int_{\Gamma_1} g_i v_i d\Gamma + \int_{\Omega} f_i v_i dx. \quad (2.33)$$

Using (2.29),  $a(u, v)$  can be written as

$$a(u, v) = \int_{\Omega} (\lambda \operatorname{div} u \cdot \operatorname{div} v + 2\mu \varepsilon_{ij}(u) \varepsilon_{ij}(v)) dx,$$

from which it is clear that  $a(\cdot, \cdot)$  is symmetric. That  $a(\cdot, \cdot)$  is  $V$ -elliptic is a nontrivial statement and the reader can refer to CIARLET [9]. Formal application of Green's formula will show that the boundary value problem corresponding to the variational problem (2.31) - (2.33) is (2.30).

$a(\cdot, \cdot)$  can be interpreted as the internal work and  $L(\cdot)$  as the work of the external loads. Thus, the equation

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

is a reformulation of the theorem of virtual work.

26 (b) **PLATE PROBLEM.** Let  $2\eta$  be the thickness of the plate. By allowing  $\eta \rightarrow 0$  in (a) we obtain the equations for the plate problem. It will be a two dimensional problem.

We have to find the bending moments  $M_{ij}$  and displacement ( $u_i$ ). These two satisfy the equations:

$$M_{ij} = \alpha \Delta u \delta_{ij} + \beta \frac{\partial^2 u}{\partial x_i \partial x_j}, \quad (2.34)$$

$$\frac{\partial^2 M_{ij}}{\partial x_i \partial x_j} = f \quad \text{in } \Omega, \quad (2.35)$$

$$u = 0 \quad \text{on } \Gamma, \quad (2.36)$$

and

$$\frac{\partial u}{\partial n} = 0 \quad \text{if the plate is clamped,} \quad (2.37)$$

$$M_{ij}n_i n_j = 0 \quad \text{if the plate is simply supported} \quad (2.37a)$$

We take

$$V = \begin{cases} H^2_0(\Omega) = \{v \in H^2 : V = \frac{\partial v}{\partial n} = 0 \text{ on } \Gamma\}, & \text{if the plate is clamped;} \\ H^2(\Omega) \cap H^1_0(\Omega), & \text{if the plate is simply supported} \end{cases} \quad (2.38)$$

Formally, using Green's formula we obtain

$$\begin{aligned} \int_{\Omega} \frac{\partial^2 M_{ij}}{\partial x_i \partial x_j} v \, dx &= - \int_{\Omega} \frac{\partial M_{ij}}{\partial x_j} \frac{\partial v}{\partial x_i} \, dx + \int_{\Gamma} \frac{\partial M_{ij}}{\partial x_j} v n_i \, d\Gamma \\ &= \int_{\Omega} M_{ij} \frac{\partial^2 v}{\partial x_i \partial x_j} \, dx - \int_{\Gamma} M_{ij} n_j \frac{\partial v}{\partial x_i} \, d\Gamma \\ &\quad + \int_{\Gamma} \frac{\partial M_{ij}}{\partial x_j} v n_i \, d\Gamma \quad \text{for all } v \in V. \end{aligned} \quad (2.39)$$

But

27

$$\int_{\Gamma} \frac{\partial M_{ij}}{\partial x_j} v n_i \, d\Gamma = 0 \quad \text{for all } v \in V, \quad \text{since } v = 0 \quad \text{on } \Gamma,$$

and

$$\int_{\Gamma} M_{ij} n_j \frac{\partial v}{\partial x_i} \, d\Gamma = \int_{\Gamma} M_{ij} n_j \left( n_i \frac{\partial v}{\partial n} + s_i \frac{\partial v}{\partial s} \right) \, d\Gamma,$$

where  $\partial v/\partial n$  denotes the normal derivative of  $v$  and  $\partial v/\partial s$  denotes the tangential derivative. By (2.37) and (2.37a) we have

$$\int_{\Gamma} M_{ij} n_j \frac{\partial v}{\partial x_i} d\Gamma = 0.$$

Hence

$$\int_{\Omega} f v dx = \int_{\Omega} \frac{\partial^2 M_{ij}}{\partial x_i \partial x_j} v dx = \int_{\Omega} M_{ij} \frac{\partial^2 v}{\partial x_i \partial x_j} dx \quad \text{for all } v \in V.$$

We therefore choose

$$a(u, v) = \int_{\Omega} M_{ij} \frac{\partial^2 v}{\partial x_i \partial x_j} dx = \int_{\Omega} \left( \alpha \Delta u \cdot \Delta v + \beta \frac{\partial^2 u}{\partial x_i \partial x_j} \frac{\partial^2 v}{\partial x_i \partial x_j} \right) dx \quad (2.39)$$

and

$$L(v) = \int_{\Omega} f v dx. \quad (2.40)$$

$a(\cdot, \cdot)$  can be proved to be  $V$ -coercive if  $\beta \geq 0$  and  $\alpha \geq 0$ .

**REGULARITY THEOREM 4.** *When  $\Omega$  is smooth and  $f \in L_2(\Omega)$ , then the solution  $u$  of the problem*

$$\begin{aligned} -\Delta u &= f & \text{in } \Omega, \\ u &= 0 & \text{on } \Gamma, \end{aligned}$$

**28** *belongs to  $H^2(\Omega)$ . Moreover, we have*

$$\|u\|_2 \leq C \|f\|_0 = C \|\Delta u\|_0,$$

*where  $C$  is a constant.*

This proves the coerciveness of  $a(u, v)$  above for  $\beta = 0$  and  $\alpha > 0$ .

## 2.6 Stokes Problem.

The motion of an incompressible, viscous fluid in a region  $\Omega$  is governed by the equations

$$-\Delta u + \nabla p = f \quad \text{in } \Omega, \quad (2.41)$$

$$\operatorname{div} u = 0 \quad \text{in } \Omega, \quad (2.42)$$

$$u = 0 \quad \text{in } \Gamma; \quad (2.43)$$

where  $u = (u_i)_{i=1,\dots,n}$  denotes the velocity of the fluid and  $p$  denotes the pressure. We have to solve for  $u$  and  $p$ , given  $f$ .

We impose the condition (2.42) in the space  $V$  itself. That is, we define

$$V = \{v \in (H_0^1(\Omega))^n : \operatorname{div} v = 0\} \quad (2.44)$$

Taking the scalar product on both sides of equation (2.41) with  $v \in V$  and integrating, we obtain

$$\int_{\Omega} f \cdot v \, dx = \int_{\Omega} \frac{\partial u_i}{\partial x_j} \frac{\partial v_i}{\partial x_j},$$

since

$$\begin{aligned} - \int_{\Omega} v \cdot \Delta u &= - \int_{\Omega} v_j \frac{\partial^2 u_j}{\partial x_i \partial x_i} \\ &= \int_{\Omega} \frac{\partial v_j}{\partial x_i} \frac{\partial u_j}{\partial x_i} - \int_{\Gamma} v_j \frac{\partial u_j}{\partial x_i} n_i, \end{aligned}$$

and

$$\int_{\Omega} \nabla p \cdot v = \int_{\Omega} \frac{\partial p}{\partial x_i} v_i = - \int_{\Omega} p \frac{\partial v_i}{\partial x_i} + \int_{\Gamma} p v_i n_i = 0$$

29

as  $v \in V$ . Therefore we define

$$a(u, v) = \int_{\Omega} \frac{\partial u_i}{\partial x_j} \frac{\partial v_i}{\partial x_j} \, dx \quad (2.45)$$

$$L(v) = \int_{\Omega} f \cdot v \, dx. \quad (2.46)$$

We now have the technical lemma.

**LEMMA 5.** *The space*

$$\vartheta = \{v \in (\mathcal{D}(\Omega))^n : \operatorname{div} v = 0\}$$

*is dense in  $V$ .*

The proof of this Lemma can be found in LADYZHENSKAYA [27].

The equation  $a(u, v) = L(v)$  for all  $v \in V$  with  $a(\cdot, \cdot), L(\cdot), v$  defined by (2.44) - (2.46) is then equivalent to

$$\langle \Delta u + f, \phi \rangle = 0 \quad \text{for all } \phi \in \vartheta, \quad (2.47)$$

where  $\langle \cdot, \cdot \rangle$  denotes the duality bracket between  $(\mathcal{D}'(\Omega))^n$  and  $(\mathcal{D}(\Omega))^n$ . Notice that (2.46) is not valid for all  $\phi \in (\mathcal{D}(\Omega))^n$  since  $(\mathcal{D}(\Omega))^n$  is not contained in  $\vartheta$ . To prove conversely that the solution of (2.46) satisfies (2.41), we need

**THEOREM 6.** *The annihilator  $\vartheta^\perp$  of  $\vartheta$  in  $(\mathcal{D}'(\Omega))^n$  is given by  $\vartheta^\perp = \{v : \text{there exists a } p \in \mathcal{D}'(\Omega) \text{ such that } v = \nabla p\}$ .*

30 Theorem 2.6 and Equation (2.46) imply that there exists a  $p \in \mathcal{D}'(\Omega)$  such that

$$\Delta u + f = \Delta p.$$

Since

$$u \in (H^1(\Omega))^n \quad \text{and} \quad f \in (L^2(\Omega))^n, \Delta u + f \in (H^{-1}(\Omega))^n.$$

Therefore

$$\nabla p \in (H^{-1}(\Omega))^n.$$

We now state

**THEOREM 7.** *If  $p \in \mathcal{D}'(\Omega)$  and  $\nabla p \in (H^{-1}(\Omega))^n$ , then  $p \in L^2(\Omega)$  and*

$$\|p\|_{L^2(\Omega)/\mathbb{R}} \leq C \|\nabla p\|_{(H^{-1}(\Omega))^n}$$

where  $C$  is a constant.

From this Theorem we obtain that  $p \in L^2(\Omega)$ . Thus, if  $f \in L^2(\Omega)$  and  $\Omega$  is smooth, we have proved that the problem (2.44) - (2.46) has a solution  $u \in V$  and  $p \in L^2(\Omega)$ .





## Chapter 3

# Conforming Finite Element Methods

IN CHAPTER 2 WE dealt with the abstract variational problems and some examples. In all our examples the function space  $V$  is infinite dimensional. Our aim is to approximate  $V$  by means of finite dimensional subspaces  $V_h$  and study the problem in  $V_h$ . Solving the variational problem in  $V_h$  will correspond to solving some system of linear equations. In this Chapter we will study an error estimate, the construction of  $V_h$  and examples of finite elements. 31

### 3.1 Approximate Problem.

The abstract variational problem is:

$$\text{find } u \in V \text{ such that } a(u, v) = L(v) \quad \text{for all } v \in V, \quad (3.1)$$

where  $a(\cdot, \cdot)$ ,  $L(\cdot)$ ,  $V$  are as in Chapter 2.

Let  $V_h$  be a finite dimensional subspace of  $V$ . Then the approximate problem corresponding to (3.1) is:

$$\text{find } u_h \in V_h \text{ such that } a(u_h, v) = L(v) \quad \text{for all } v \in V_h. \quad (3.2)$$

By the Lax-Milgram Lemma (Chapter 2, Theorem 2.1), (3.2) has a unique solution.

Let dimension  $(V_h) = N(h)$  and let  $(w_i)_{i=1, \dots, N(h)}$  be a basis of  $V_h$ . Let

$$u_h = \sum_{i=1}^{N(h)} u_i w_i, v_h = \sum_{j=1}^{N(h)} v_j w_j,$$

32 where  $u_i, v_j \in \mathbb{R}$ ,  $1 \leq i, j \leq N(h)$ . Substituting these in (3.2), we obtain

$$\sum_{i,j=1}^{N(h)} u_i v_j a(w_i, w_j) = \sum_{j=1}^{N(h)} L(w_j) v_j \quad (3.3)$$

Let

$$A^T = (a(w_i, w_j))_{i,j}, U = (u_i)_i, V = (v_i)_i, b = (L(w_i))_i$$

Then (3.3) can be written as

$$V^T A U = V^T b.$$

This is true for all  $V \in \mathbb{R}^{N(h)}$ . Hence

$$A U = b. \quad (3.4)$$

If the linear system (3.4) is solved, then we know the solution  $u_h$  of (3.2). This approximation method is called the Rayleigh-Galerkin method.

$A$  is positive definite since

$$\begin{aligned} V^T A V &= \sum_{i,j=1}^{N(h)} a(w_i, w_j) v_i v_j = a \left( \sum_{i=1}^{N(h)} w_i v_i, \sum_{j=1}^{N(h)} w_j v_j \right) \\ &\geq \alpha \left\| \sum_{i=1}^{N(h)} v_i w_i \right\|^2 \quad \text{for all } V \in \mathbb{R}^{N(h)}. \end{aligned}$$

$A$  is symmetric if the bilinear form  $a(\cdot, \cdot)$  is symmetric.

33 From the computational point of view it is desirable to have  $A$  as a sparse matrix, i.e.  $A$  has many zero elements. Usually  $a(\cdot, \cdot)$  will be given by an integral and the matrix  $A$  will be sparse if the support of the basis functions is "small". For example, if

$$a(u, v) = \int_{\Omega} \nabla u \cdot \nabla v \, dx,$$

then  $a(w_i, w_j) = 0$  if  $\text{Supp } w_i \cap \text{Supp } w_j = \emptyset$ .

Now we will prove a theorem regarding the error committed when the approximate solution  $u_h$  is taken instead of the exact solution  $u$ .

**THEOREM 1.** *If  $u$  and  $u_h$  denote the solutions of (3.1) and (3.2) respectively, then we have*

$$\|u - u_h\|_V \leq C \inf_{v_h \in V_h} \|u - v_h\|_V$$

*Proof.* We have

$$\begin{aligned} a(u, v) &= L(v) \quad \text{for all } v \in V, \\ a(u_h, v) &= L(v) \quad \text{for all } v \in V_h; \end{aligned}$$

so

$$a(u - u_h, v) = 0 \quad \text{for all } v \in V_h \quad (3.5)$$

By the  $V$ -coerciveness of  $a(\cdot, \cdot)$  we obtain

$$\begin{aligned} \|u - u_h\|^2 &\leq 1/\alpha a(u - u_h, u - u_h) \\ &= 1/\alpha a(u - u_h, u - v + v - u_h), \quad \text{for all } v \in V_h \\ &= 1/\alpha a(u - u_h, u - v), \quad \text{by (3.5)} \\ &\leq M/\alpha \|u - u_h\| \|u - v\| \end{aligned}$$

□

This proves the theorem with  $C = M/\alpha$ .

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### 3.2 Internal Approximation of $H^1(\Omega)$ .

Let  $\Omega \subset \mathbb{R}^2$  be a polygonal domain. Let  $T_h$  be a triangulation of  $\Omega$ : that is  $T_h$  a finite collection of triangles such that

$$\overline{\Omega} = \bigcup_{K \in T_h} \overline{K} \quad \text{and} \quad K \cap K' = \emptyset \quad \text{for} \quad K, K' \in T_h, K \neq K'.$$

Let  $P(K)$  be a function space defined on  $K$  such that  $P(K) \subset H^1(K)$ . Usually we take  $P(K)$  to be the space of polynomials of some degree. We have

**THEOREM 2.** *If*

$$V_h = \{v_h \in C^0(\Omega) : v_h|_K \in P(K), K \in T_h\}$$

where  $P(K) \subset H^1(K)$ , then  $V_h \subset H^1(\Omega)$ .

*Proof.* Let  $u \in V_h$  and  $v_i$  be a function defined on  $\Omega$  such that  $v_i|_K = \frac{\partial}{\partial x_i}(u|_K)$ . This makes sense since  $u|_K \in H^1(K)$ . Moreover  $v_i \in L^2(\Omega)$ , since

$$v_i|_K = \frac{\partial}{\partial x_i}(u|_K) \in L^2(K). \text{ We will show that } v_i = \frac{\partial u}{\partial x_i} \text{ in } \mathcal{D}'(\Omega).$$

For any  $\phi \in \mathcal{D}(\Omega)$ , we have

$$\begin{aligned} \langle v_i, \phi \rangle &= \int_{\Omega} v_i \phi \, dx = \sum_{K \in T_h} \int_K v_i \phi \, dx = \sum_{K \in T_h} \int_K \frac{\partial}{\partial x_i}(u|_K) \phi \, dx \\ &= \sum_{K \in T_h} - \int_K (u|_K) \frac{\partial \phi}{\partial x_i} \, dx + \int_K (u|_K) \phi \, n_i^K \, d\Gamma, \end{aligned}$$

35 where  $n_i^K$  is the  $i^{\text{th}}$  component of the outward drawn normal to  $\partial K$ . So

$$\langle v_i, \phi \rangle = - \int_{\Omega} u \frac{\partial \phi}{\partial x_i} \, dx + \sum_{K \in T_h} \int_{\partial K} (u|_K) \phi \, n_i^K \, d\Gamma \quad (3.6)$$

The second term on the right hand side of (3.6) is zero since  $u$  is continuous in  $\Omega$  and if  $K_1$  and  $K_2$  are two adjacent triangles then  $n_i^{K_1} = -n_i^{K_2}$ . Therefore

$$\langle v_i, \phi \rangle = - \int_{\Omega} u \frac{\partial \phi}{\partial x_i} \, dx = \left\langle \frac{\partial u}{\partial x_i}, \phi \right\rangle$$

which implies

$$v_i = \frac{\partial u}{\partial x_i} \text{ in } \mathcal{D}'(\Omega).$$

Hence  $u \in H^1(\Omega)$ . Thus  $V_h \subset H^1(\Omega)$ .

We assume that the triangulation  $T_h$  is such that if  $K_1, K_2 \in T_h$  are distinct, then either  $\overline{K_1} \cap \overline{K_2}$  is empty or equal to the common edge of

the triangles  $K_1$  and  $K_2$ . By this assumption we eliminate the possibility of a triangulation as shown in figure.

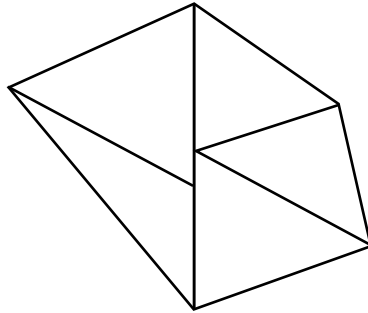


Figure 3.1:

**Construction of  $V_h$ .**

Let  $\Omega$  be a polygonal domain and  $T_h$  be a triangulation of  $\Omega$ , where

$$h = \max_{K \in T_h} (\text{diameter of } K).$$

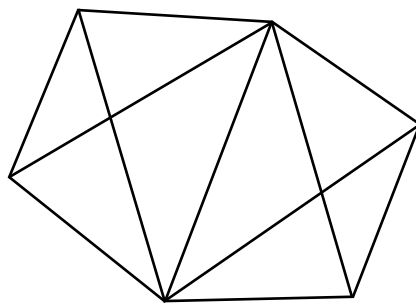


Figure 3.2:

Let

$$N(h) = \# \text{ nodes of the triangulation,} \quad (3.8)$$

$$P(K) = P_1(K) = \begin{array}{l} \text{polynomial of degree less than or equal to} \\ 1 \text{ in } x \text{ and } y \end{array} \quad (3.9)$$

Let

$$V_h = \{v_h : v_h|_K \in P_1(K), K \in T_h\} \quad (3.10)$$

We know that a polynomial of degree 1 in  $x$  and  $y$  is uniquely determined if its values on three non-collinear points are given. Using this we construct a basis for  $V_h$ . A function in  $V_h$  is uniquely determined if its value at all the nodes of the triangulation is given. Let the nodes of the triangulation be numbered  $\{1, 2, \dots, N(h)\}$ . Let  $w_i \in V_h$  be

$$w_i = \begin{cases} 1 & \text{at the } i^{\text{th}} \text{ node,} \\ 0 & \text{at other nodes.} \end{cases} \quad (3.11)$$

37 It is easy to see that  $w_i$  are linearly independent. If  $v \in V_h$ , then

$$v = \sum_{i=1}^{N(h)} v_i w_i \quad (3.12)$$

where  $v_i$  the value of  $v$  at the  $i^{\text{th}}$  node. This proves that  $\{w_i\}_{1, \dots, N(h)}$  is a basis of  $V_h$  and dimension of  $V_h = N(h)$ .

Moreover,  $\text{Supp } w_i \subset \cup \bar{K}$ , where the union is taken over all the triangles whose one of the vertices is the  $i^{\text{th}}$  node.

Hence if  $i^{\text{th}}$  node and  $j^{\text{th}}$  node are not the vertices of a triangle  $K$ , for any  $K \in T_h$ , then

$$\text{Supp } w_i \cap \text{Supp } w_j = \phi.$$

38 We will show that  $V_h$  given by (3.10) is contained in  $C^0(\bar{\Omega})$ . Let  $v \in V_h$  and let  $K_1, K_2 \in T_h$  be adjacent triangles. Let ' $\ell$ ' be the side common to both  $K_1$  and  $K_2$ .  $v|_{K_1}$  and  $v|_{K_2}$  are polynomials of degree less than or equal to one in  $x$  and  $y$ . Let  $\tilde{v}_1$  and  $\tilde{v}_2$  be the extensions of  $v|_{K_1}$  and  $v|_{K_2}$  to  $\bar{K}_1$  and  $\bar{K}_2$  respectively.  $\tilde{v}_1|_\ell$  and  $\tilde{v}_2|_\ell$  can be thought of as a polynomial of degree less than or equal to one in a **single variable** and hence can be determined uniquely if their values at two distinct points are known. But, by the definition of  $V_h$  in (3.10),  $\tilde{v}_1|_\ell$  and  $\tilde{v}_2|_\ell$  agree at the common vertices of  $K_1$  and  $K_2$ . Hence  $\tilde{v}_1|_\ell = \tilde{v}_2|_\ell$ . This proves that  $v$  is continuous across  $K_1$  and  $K_2$ . Thus  $v \in C^0(\bar{\Omega})$ . Hence  $V_h \subset C^0(\bar{\Omega})$ .

Using the theorem 3.2 we conclude that  $V_h \subset H^1(\Omega)$ . When we impose certain restrictions on  $T_h$ , it is possible to prove that  $d(u, V_h) \rightarrow 0$  as  $h \rightarrow 0$  where  $d(u, V_h)$  is the distance between the solution  $u$  of (3.1) and the finite dimensional space  $V_h$ . The reader can refer to CIARLET [9]. Thus  $V_h$  “approximates”  $H^1(\Omega)$ .

The finite element method and the finite difference scheme are the “same” when the triangulation is uniform. For elliptic problems the finite element method gives better results than the finite difference scheme.  $\square$

### 3.3 Finite Elements of Higher Degree.

**DEFINITION.** Let  $K$  be a triangle with vertices  $(a_i, i = 1, 2, 3)$ . Let the coordinates of  $a_i$  be  $a_{ij}$ ,  $j = 1, 2$ . For any  $x \in \mathbb{R}^2$ , the barycentric coordinates  $\lambda_i(x), i = 1, 2, 3$ , of  $x$  are defined to be the unique solution of the linear system

$$\begin{aligned} \sum_{i=1}^3 \lambda_i a_{ij} &= x_j, \quad j = 1, 2; \\ \sum_{i=1}^3 \lambda_i &= 1 \end{aligned} \tag{3.13}$$

Notice that the determinant of the coefficient matrix of the system (3.13) is twice the area of the triangle  $K$ . It is easy to see that the barycentric coordinates of  $a_1, a_2, a_3$  are  $(1, 0, 0)$ ,  $(0, 1, 0)$  and  $(0, 0, 1)$  respectively. The barycentric coordinate of the centroid  $G$  of  $K$  is  $(1/3, 1/3, 1/3)$ .

Using Cramer rule we find from (3.13) that

$$\lambda_1 = \frac{\begin{vmatrix} x_1 & a_{21} & a_{31} \\ x_2 & a_{22} & a_{32} \\ 1 & 1 & 1 \end{vmatrix}}{\begin{vmatrix} a_{11} & a_{21} & a_{31} \\ a_{12} & a_{22} & a_{32} \\ 1 & 1 & 1 \end{vmatrix}}$$

$$\begin{aligned} \text{i.e.} \quad \lambda_1 &= \frac{\text{area of the triangle } xa_2a_3}{\text{area of the triangle } a_1a_2a_3} \\ \text{Similarly, } \lambda_2 &= \frac{\text{area of the triangle } a_1xa_3}{\text{area of the triangle } a_1a_2a_3} \\ \lambda_3 &= \frac{\text{area of the triangle } a_1a_2x}{\text{area of the triangle } a_1a_2a_3} \end{aligned}$$

This geometric interpretation of the barycentric coordinates will be helpful in specifying the barycentric coordinates of a point. For example, the equation of the side  $a_2a_3$  in barycentric coordinates is  $\lambda_1 = 0$ .

- 40 DEFINITION.** A **finite element** is a triple  $(K, P_K, \Sigma_K)$ , where  $K$  is a polyhedron,  $P_K$  is polynomial space whose dimension is  $m$  and  $\Sigma_K$  is a set of distributions, whose cardinality is  $m$ . Further  $\Sigma_K = \{L_i \in \mathcal{D}' ; i = 1, 2, \dots, m\}$  is such that for given  $d_i \in \mathbb{R}$ ,  $1 \leq i \leq m$ , the equations

$$L_i(p) = d_i, \quad 1 \leq i \leq m$$

have a unique solution  $p \in P_K$ .

The elements  $L_i$  are called degree of freedom of  $P$ .

**EXAMPLE 1.** (Finite Element of Degree 1). Let  $K = a$  triangle,

$$\begin{aligned} P_K &= P_1(K) = \text{Polynomials of degree } \leq 1. \\ &= \text{Span } \{1, x, y\} \end{aligned}$$

$$\dim P_K = 3,$$

$$\Sigma_K = \{\delta_{a_i} : a_i \text{ vertices, } i = 1, 2, 3\},$$

where  $\delta_{a_i}$  denotes the dirac mass at the point  $a_i$ . Then  $(K, P_K, \Sigma_K)$  is a finite element.

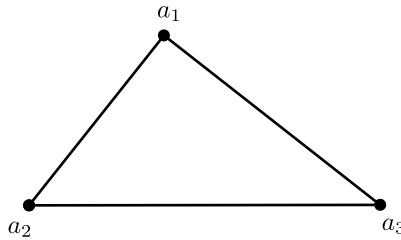


Figure 3.3:



This follows from the fact that  $p \in P_1(K)$  is uniquely determined if its values at three non collinear points are given.

$$\delta_{a_i}(\lambda_j) = \lambda_j(a_i) = \delta_{ij}.$$

Hence  $\lambda_j, j = 1, 2, 3$ , form a basis for  $P_1(K)$  and if  $p \in P_1(K)$  then 41

$$p = \sum_{i=1}^3 p(a_i)\lambda_i$$

**REMARK 1.** In the definition, dimension of  $P_K$  is  $m$  and we require that the equations

$L_i(p) = d_i, 1 \leq i \leq m$ , for given  $d_i \in \mathbb{R}$  have a solution. So, in examples, to prove existence we have to prove only uniqueness. To prove the uniqueness it is enough to show that

$$L_i(p) = 0, 1 \leq i \leq m, \quad \text{implies} \quad p \equiv 0.$$

**REMARK 2.** If  $p_j \in P_K, 1 \leq j \leq m$ , are such that

$$L_i(p_j) = \delta_{ij}, 1 \leq i \leq m, 1 \leq j \leq m,$$

then  $\{p_j\}$  form a basis for  $P_K$  and any  $p \in P_K$  can be written as

$$p = \sum_{i=1}^m L_i(p)P_i.$$

**EXAMPLE 2.** (Finite Element of Degree 2). Let  $K = a$  triangle,

$$P_K = P_2(K) = \text{Span} \{1, x, y, x^2, xy, y^2\},$$

$$\sum_K = \{\delta_{a_i}, 1 \leq i \leq 3, \delta_{a_{ij}} : 1 \leq i < j \leq 3\},$$

where  $a_i$  denote the vertices of  $K$  and  $a_{ij}$  denote the mid point of the side  $a_i a_j$ .

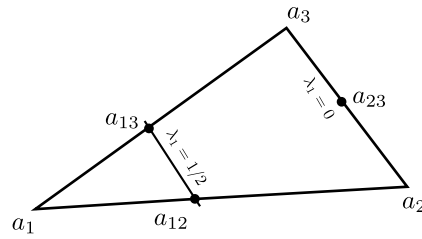


Figure 3.4:

- 42 The equations of the lines  $a_3a_2$  and  $a_{13}a_{12}$  are  $\lambda_1 = 0$  and  $\lambda_1 = 1/2$  respectively. Hence the function  $\lambda_1(\lambda_1 - 1/2)$  vanishes at the points  $a_2, a_3, a_{12}, a_{23}, a_{13}$ . The value of  $\lambda_1(\lambda_1 - 1/2)$  at  $a_1$  is  $1/2$ . Hence  $\lambda_1(2\lambda_1 - 1)$  takes the value 1 at  $a_1$  and 0 at other nodes.

The equations of the lines  $a_1a_3$  and  $a_2a_3$  are  $\lambda_2 = 0$  and  $\lambda_1 = 0$  respectively. Therefore the function  $\lambda_1\lambda_2$  vanishes at  $a_1, a_2, a_{13}, a_{23}, a_3$  and takes the value  $1/4$  at  $a_{12}$ . Thus  $4\lambda_1\lambda_2$  is 1 at  $a_{12}$  and zero at the other nodes.

Thus any  $p \in P_2(K)$  can be written in the form

$$p = \sum_{i=1}^3 p(a_i)\lambda_i(2\lambda_i - 1) + \sum_{\substack{i < j \\ i, j=1}}^3 4p(a_{ij})\lambda_i\lambda_j.$$

**EXAMPLE 3.** (Finite Element of Degree 3). Let  $K = a$  triangle,

$$P_K = P_3(K) = \text{Span} \{1, x, y, x^2, xy, y^2, x^3, x^2y, xy^2, y^3\}.$$

Thus  $\dim P_3 = 10$ .

$$\sum_K = \{\delta_{a_i}, 1 \leq i \leq 3; \delta_{a_{ij}}, 1 \leq i \leq j \leq 3, \delta_{a_{123}}\}.$$

where  $a_i$  denote the vertices of  $K$  and  $a_{ij} = \frac{2}{3}a_i + \frac{1}{3}a_j$ .

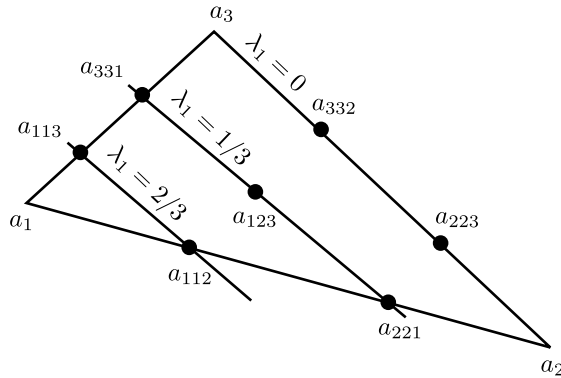


Figure 3.5:

It is easy to see that

43

$$\begin{aligned} p_i &= 1/2 \lambda_i(3\lambda_i - 1)(3\lambda_i - 2), \\ p_{ij} &= 9/2 \lambda_i \lambda_j (3\lambda_i - 1), \\ p_{123} &= 27 \lambda_1 \lambda_2 \lambda_3, \end{aligned}$$

$1 \leq i, j \leq 3$ , is a basis of  $P_3(K)$ .

Moreover,  $p_i$  is 1 at the node  $a_i$  and zero at the other nodes;  $p_{ij}$  is 1 at the node  $a_{ij}$  and vanishes at the other nodes;  $p_{123}$  is zero at all nodes except  $a_{123}$  where its value is 1.

**REMARK 3.** In the above three examples  $\Sigma_K$  contains only Dirac masses and not derivatives of Dirac masses. All the above three finite elements are called Lagrange finite elements.

Let  $\Omega$  be a polygonal domain and let  $T_h$  be a triangulation of  $\Omega$ , i.e.  $\overline{\Omega} = \bigcup_{K \in T_h} \overline{K}$ . Let  $P_K = P_\ell(K)$  consists of polynomials of degree  $\leq \ell$ . Let  $(K, P_K, \Sigma_K)$  be a finite element for each  $K \in T_h$ . Let  $\Sigma_h = \bigcup_{K \in T_h} \Sigma_K$  and

$$V_h = \{v_h : v_h|_K \in P_K, K \in T_h\}$$

From the definition of finite element it follows that a function in  $V_h$  is uniquely determined by the distributions in  $\Sigma_h$ . 44

**REMARK 4.** *In the example*

$\Sigma_h = \{\delta_{a_i} : a_i \text{ is a vertex of a triangle in the triangulation}\}$ . We proved in Sec. 3.2 that  $V_h \subset H^1(\Omega)$ . In this case we say that the finite element is conforming.

**REMARK 5.** *The  $V_h$  so constructed above need not be contained in  $H^1(\Omega)$ . If  $V_h \not\subset H^1(\Omega)$  we say that the finite element method is non-conforming.*

Let  $K$  be a triangle,  $P_K = P_1(K)$  and

$$\Sigma_K = \{\delta_{a_{ij}} : 1 \leq i < j \leq 3\}.$$

Then  $(K, P_K, \Sigma_K)$  will be a finite element, but the space  $V_h \not\subset H^1(\Omega)$ .

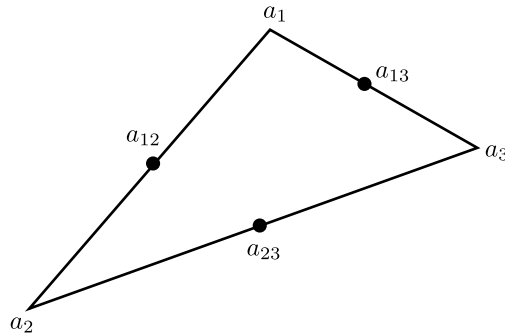


Figure 3.6:

45 When  $(K, P_K, \Sigma_K)$  is as in Example 2, we will prove that  $V_h \subset C^0(\overline{\Omega})$ . This together with theorem 3.2 implies  $V_h \subset H^1(\Omega)$ . Thus the finite element in Example 2 is conforming.

To prove that  $V_h \subset C^0(\overline{\Omega})$ , let  $K_1$  and  $K_2$  be two adjacent triangles in the triangulation.

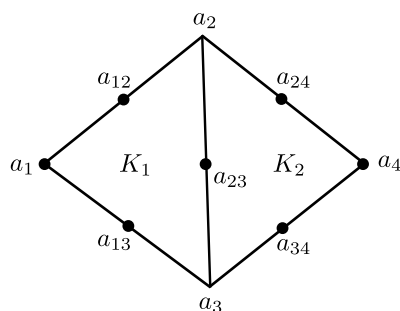


Figure 3.7:

A polynomial of degree 2 in  $x$  and  $y$  when restricted to a line in the plane is a polynomial of degree 2 in a single variable and hence can be determined on the line if the value of the polynomial at three distinct points on the line are known. Let  $v_h \in V_h$ . Let  $\tilde{v}_1$  and  $\tilde{v}_2$  be the continuous extension of  $v_h|_{K_1}$  and  $v_h|_{K_2}$  to  $\overline{K_1}$  and  $\overline{K_2}$  respectively;  $\tilde{v}_1$  and  $\tilde{v}_2$  are polynomials of degree 2 in one variable along the common side;  $\tilde{v}_1$  and  $\tilde{v}_2$  agree at the two common vertices and at the midpoint of the common side. Hence  $\tilde{v}_1 = \tilde{v}_2$  on the common side. This shows  $v_h$  is continuous. Hence  $V_h \subset C^0(\overline{\Omega})$ .

**Exercise 1.** Taking  $(K, P_K, \Sigma_K)$  as in Example 3, show that  $V_h \subset H^1(\Omega)$ .

### 3.4 Internal Approximation of $H^2(\Omega)$ .

In this section we give an example of a finite element which is such that the associated space  $V_h$  is contained in  $H^2(\Omega)$ . This finite element can be used to solve some fourth order problems. We need 46

**THEOREM 3.** *If  $V_h = \{v_h : v_h|_K \in P_K \subset H^2(K), \text{ for all } K \in T_h\}$  is contained in  $C^1(\overline{\Omega})$  then  $V_h$  is contained in  $H^2(\Omega)$ .*

The proof of this theorem is similar to that of theorem 2 of this section.

**EXAMPLE 4.** Let  $K$  be a triangle

$$P_K = P_3(K) = \text{Span} \{1, x, y, x^2, xy, y^2, x^3, x^2y, xy^2, y^3\},$$

$$\sum_K = \left\{ \delta_{a_i}, \frac{\partial}{\partial x} \delta_{a_i}, \frac{\partial}{\partial y} \delta_{a_i}, \delta_{a_{123}}, 1 \leq i \leq 3 \right\}$$

where  $a_i$  are vertices of  $K$ ,  $a_{123}$  is the centroid of  $K$  and

$$\dim P_K = \text{Card} \sum_K = 10.$$

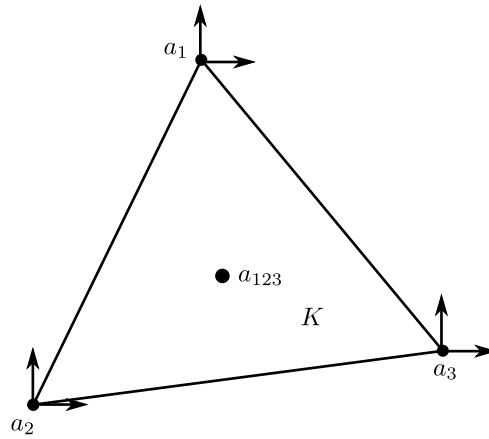


Figure 3.8:

The arrows in the figure denote that the values of the derivatives at the vertices are given. Using the formula

$$p = \sum_{i=1}^3 (-2\lambda_i^3 + 3\lambda_i^2 - 7\lambda_1\lambda_2\lambda_3) p(a_i) + 27\lambda_1\lambda_2\lambda_3 p(a_{123})$$

$$+ \sum_{1 \leq i < j \leq 3} \lambda_i \lambda_j (2\lambda_i + \lambda_j - 1) Dp(a_i)(a_j - a_i) \quad \text{for all } p \in P_3$$

where

$$Dp(a_i) = \left( \frac{\partial p}{\partial x}(a_i), \frac{\partial p}{\partial y}(a_i) \right).$$

We obtain that  $p \equiv 0$  if

$$p(a_i) = p(a_{123}) = \frac{\partial p}{\partial x}(a_i) = \frac{\partial p}{\partial y}(a_i) = 0, 1 \leq i \leq 3.$$

Hence  $(K, P_K, \Sigma_K)$  is a finite element.

The corresponding  $V_h$  is in  $C^0(\bar{\Omega})$  but not in  $C^1(\bar{\Omega})$ . This shows that  $V_h \not\subset H^2(\Omega)$ . Hence this is not a conforming finite element for fourth order problems.

We now give an example of a finite element with  $V_h \subset H^2(\Omega)$ .

**EXAMPLE 5. The Argyris Triangle.** The Argyris triangle has 21 degrees of freedom. Here the values of the polynomial, its first and second derivatives are specified at the vertices; the normal derivative is given at the mid points.

In the figure we denote the derivatives by circles and normal derivative by a straight lines.

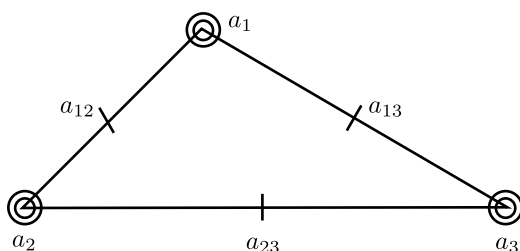


Figure 3.9:

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We take  $P_K = P_5 =$  Space of polynomials of degree less than or equal to 5.

$$\dim P_K = 21;$$

$$\sum_K = \left\{ \delta_{a_i}, \frac{\partial}{\partial x} \delta_{a_i}, \frac{\partial}{\partial y} \delta_{a_i}, \frac{\partial^2}{\partial x^2} \delta_{a_i}, \right.$$

$$\left. \frac{\partial^2}{\partial x \partial y} \delta_{a_i}, \frac{\partial^2}{\partial y^2} \delta_{a_i}, 1 \leq i \leq 3, \frac{\partial}{\partial n} \delta_{a_{ij}} 1 \leq i < j \leq 3 \right\},$$

where  $a_i$  denote the vertices of  $K$ ,  $a_{ij}$  the midpoint of the line joining  $a_i$  and  $a_j$ , and  $\partial/\partial n$ , the normal derivative.

Let  $p \in P_K$  be such that  $L(p) = 0$ ,  $L \in \Sigma_K$ . We will prove that  $p = 0$  in  $K$ .  $p$  is a polynomial of degree 5 in one variable along the side  $a_1 a_2$ . By assumption  $p$ , and its first and second derivatives vanish at  $a_1$  and  $a_2$ . Hence  $p = 0$  along  $a_1 a_2$ . Consider  $\partial p / \partial n$  along  $a_1 a_2$ . By assumption  $\partial p / \partial n$  vanishes at  $a_1 a_2$  and  $a_{12}$ . Since the second derivatives of  $p$  vanish at  $a_1, a_2$  we have the first derivatives of  $\partial p / \partial n$  vanish at  $a_1$  and  $a_2$ ;  $\partial p / \partial n$  is polynomial of degree 4 in one variable along  $a_1 a_2$ . Hence  $\partial p / \partial n = 0$  along  $a_1 a_2$ . Since  $p = 0$  along  $a_1 a_2$ ,  $\partial p / \partial \tau = 0$  along  $a_1 a_2$ , where  $\partial / \partial \tau$  denote the tangential derivative, i.e. derivative along  $a_1 a_2$ . Therefore we have  $p$  and its first derivatives zero along  $a_1 a_2$ .

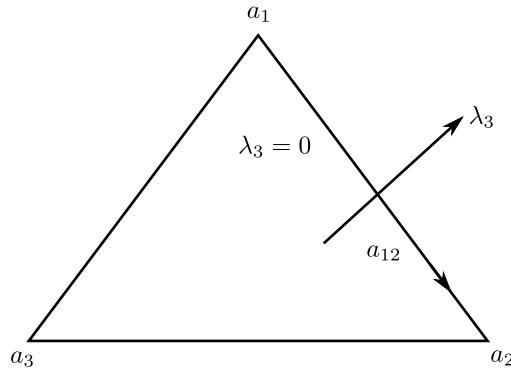


Figure 3.10:

The equation of  $a_1 a_2$  is  $\lambda_3 = 0$ . Hence we can choose a line perpendicular to  $a_1 a_2$  as the  $\lambda_3$  axis. Let  $\tau$  denote the variable along  $a_1 a_2$ . Changing the coordinates from  $(x, y)$  to  $(y_3, \tau)$  we can write the polynomial  $p$  as

$$p = \sum_{i=0}^5 \lambda_3^i q_i(\tau)$$



where  $q_i(\tau)$  is a polynomial in  $\tau$  of degree  $\leq 5 - i$ . Now

$$\frac{\partial p}{\partial \lambda_3} = \sum_{i=1}^5 i \lambda_3^{i-1} q_i(\tau).$$

Since  $p$  and its first derivatives vanish along  $a_1 a_2$  we have

$$\begin{aligned} 0 &= p(0, \tau) = q_0(\tau), \\ 0 &= \frac{\partial p}{\partial \lambda_3}(0, \tau) = q_1(\tau) \end{aligned}$$

Hence

$$p = \lambda_3^2 \sum_{i=2}^5 \lambda_3^{i-2} q_i(\tau).$$

Thus  $\lambda_3^2$  is a factor of  $p$ . By taking the other sides we can prove that  $\lambda_1^2$  and  $\lambda_2^2$  are also factors of  $p$ . Thus

$$p = c \lambda_1^2 \lambda_2^2 \lambda_3^2.$$

But  $\lambda_1^2 \lambda_2^2 \lambda_3^2$  is a polynomial of degree 6 which does not vanish identically in  $K$  and  $p$  is a polynomial of degree 5. Hence  $c = 0$ . Therefore  $p \equiv 0$  in  $K$ . Thus  $(K, P_K, \Sigma_K)$  is a finite element with 21 degrees of freedom. 50

**THEOREM 4.** *If  $V_h$  is the space associated with the Argyris finite element, then  $V_h \subset C^1(\bar{\Omega})$ .*

*Proof.* Let  $K_1, K_2$  be two adjacent finite elements in the triangulation. Let  $v \in V_h$  and let  $p_1 = v|_{K_1}$ ,  $p_2 = v|_{K_2}$ . We denote the continuous extensions of  $p_1$  and  $p_2$  to  $\bar{K}_1$  and  $\bar{K}_2$  also by  $p_1$  and  $p_2$ . We have to show that

$$p_1 = p_2, Dp_1 = Dp_2$$

along the common side  $Q$  of  $K_1$  and  $K_2$ .

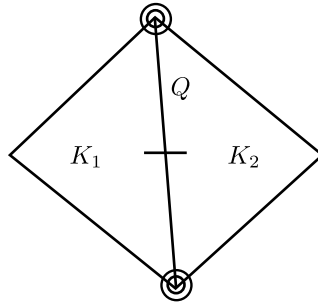


Figure 3.11:

- 51 Since  $v \in V_h$ ,  $p_1, p_2$  their first and second derivatives agree at the two common vertices of  $K_1$  and  $K_2$ .  $p_1$  and  $p_2$  are polynomials of degree 5 in one variable along the common side  $Q$ . Hence, along  $Q$ ,

$$p_1 = p_2. \quad (3.14)$$

Therefore

$$\frac{\partial p_1}{\partial \tau} = \frac{\partial p_2}{\partial \tau} \quad \text{along } Q. \quad (3.15)$$

The normal derivatives  $\frac{\partial p_1}{\partial n}$  and  $\frac{\partial p_2}{\partial n}$  are polynomials of degree 4 in one variable along  $Q$ . Since  $v \in V_h$ ,  $\frac{\partial p_1}{\partial n} = \frac{\partial p_2}{\partial n}$  at the two common vertices and at the midpoint of the common side  $Q$ . Moreover, the first derivatives of  $\frac{\partial p_1}{\partial n}$  and  $\frac{\partial p_2}{\partial n}$  coincide at the common vertices. Hence

$$\frac{\partial p_1}{\partial n} = \frac{\partial p_2}{\partial n} \quad \text{along } Q. \quad (3.16)$$

Equations (3.14) - (3.16) show that

$$p_1 = p_2, Dp_1 = Dp_2 \quad \text{along } Q.$$

□

This proves that  $v$  and its first derivative are continuous across  $Q$ . Therefore  $v \in C^1(\overline{\Omega})$ . Hence  $V_h \subset C^1(\overline{\Omega})$ .

Theorems 3.3 and 3.4 imply that  $V_h \subset H^2(\Omega)$ . For other examples of finite element, the reader can refer to CIARLET [9].

## Chapter 4

# Computation of the Solution of the Approximate Problem

### 4.1 Introduction

THE SOLUTION OF THE APPROXIMATE PROBLEM. Find  $u_h \in V_h$  52  
such that

$$a(u_h, v_h) = L(v_h) \quad \forall v_h \in V_h, \quad (4.1)$$

can be found using either iterative methods or direct methods. We describe these methods in this chapter.

Let  $\{w_i\}_{1 \leq i \leq N(h)}$  be a basis of  $V_h$ . Let  $A = (a(w_i, w_j))$  and  $b = (L(w_i))$ . If  $a(\cdot, \cdot)$  is symmetric, then (4.1) is equivalent to the minimization problem

$$J(u_h) = \min_{v_h \in V_h} J(v_h), \quad (4.2)$$

where  $J(v) = \frac{1}{2}v^T A v - v^T b$ ,  $v \in \mathbb{R}^{N(h)}$ . Here we identify  $V_h$  and  $\mathbb{R}^{N(h)}$  through the basis  $\{w_i\}$  and the natural basis  $\{e_i\}$  of  $\mathbb{R}^{N(h)}$ .  $u_h$  is a solution of (4.2) iff  $Au_h = b$ . Iterative methods are applicable only when  $a(\cdot, \cdot)$  is symmetric.

## 4.2 Steepest Descent Method

Let  $J : \mathbb{R}^N \rightarrow \mathbb{R}$  be differentiable.

In the steepest descent method, at each iteration we move along the direction of the negative gradient to a point such that the functional value of  $J$  is reduced. That is, let  $x^0 \in \mathbb{R}^N$  be given. Knowing  $x^n \in \mathbb{R}^N$ , we define  $x^{n+1} \in \mathbb{R}^N$  by

$$x^{n+1} = x^n - \lambda^n J'(x^n), \quad (4.3)$$

where  $\lambda^n$  minimizes the functional

$$\phi(\lambda) = J(x^n - \lambda J'(x^n)) \quad (4.4)$$

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In the case

$$J(x) = 1/2 x^T A x - x^T b,$$

$\lambda^n$  can be computed explicitly. It is easy to see that

$$J'(x) = Ax - b.$$

Since  $\lambda^n$  minimizes  $\phi(\lambda)$ , we have

$$\phi'(\lambda^n) = (J'(x^n - \lambda^n J'(x^n)), -J'(x^n)) = 0 \quad (4.5)$$

Let

$$r^n = J'(x^n) = Ax^n - b.$$

Then (4.5) implies

$$\lambda^n = \frac{(r^n, r^n)}{(Ar^n, r^n)}. \quad (4.6)$$

For proving an optimal error estimate for this scheme we need Kantorovich's inequality which is left as an exercise.

**Exercise 1.** (See LUENBERGER [30]).

Prove the Kantorovich's inequality

$$\frac{(Ax, x)(A^{-1}x, x)}{\|x\|^4} \leq \frac{(M+m)^2}{4mM} \quad (4.7)$$

where  $A$  is symmetric, positive definite matrix with

$$m = \inf_{x \neq 0} \frac{(Ax, x)}{\|x\|^2} > 0, M = \sup_{x \neq 0} \frac{(Ax, x)}{\|x\|^2}$$

**THEOREM 1.** For any  $x_0 \in X$  the sequence  $\{x_n\}$  defined by

$$x_{n+1} = x_n + \frac{(r_n, r_n)}{(r_n, Ar_n)} r_n,$$

where

$$r_n = b - Ax_n,$$

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converges to the unique solution  $\bar{x}$  of  $Ax = b$ . Furthermore, defining

$$E(x) = ((x - \bar{x}), A(x - \bar{x}))$$

we have the estimate

$$\|x_n - \bar{x}\|^2 \leq \frac{1}{m} E(x_n) \leq \frac{1}{m} \left( \frac{M-m}{M+m} \right)^{2n} E(x_0).$$

*Proof.* We have

$$\begin{aligned} E(x) &= (x - \bar{x}, A(x - \bar{x})) \\ &= 2J(x) + (\bar{x}, A\bar{x}), \end{aligned}$$

where

$$J(x) = 1/2(x, Ax) - (b, x).$$

It is easy to see that

$$\frac{E(x_n) - E(x_{n+1})}{E(x_n)} = \frac{(r_n, r_n)^2}{(r_n, Ar_n)(r_n, A^{-1}r_n)} \geq \frac{4Mm}{(M+m)^2}$$

by Kantorovich inequality. Therefore

$$\frac{E(x_{n+1})}{E(x_n)} \leq \left( \frac{M-m}{M+m} \right)^2.$$

This implies

$$E(x_{n+1}) \leq \left( \frac{M-m}{M+m} \right)^{2(n+1)} E(x_0).$$

From the definition of  $m$  we obtain

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$$\|x_n - \bar{x}\|^2 \leq \frac{1}{m} E(x_n) \leq \frac{1}{m} \left( \frac{M-m}{M+m} \right)^{2n} E(x_0).$$

The condition number of  $A$  is defined by  $\text{Cond}(A) = \frac{M}{m}$ . We have

$$\left( \frac{M-m}{M+m} \right)^2 \sim 1 - \frac{2m}{M} = 1 - \frac{2}{\text{Cond}(A)}$$

If the condition number of  $A$  is smaller, then the convergence is faster. The steepest descent method is not a very good method for finite elements, since  $\text{Cond}(A) \sim C/h^2$  when  $V_h \subset H^1(\Omega)$ .  $\square$

### 4.3 Conjugate Gradient method

**DEFINITION.** The directions  $w_1, w_2 \in \mathbb{R}^N$  are said to be conjugate with respect to the matrix  $A$  if  $w_1^T A w_2 = 0$ .

In the conjugate gradient method, we construct conjugate directions using the gradient of the functional. Then the functional is minimized by proceeding along the conjugate direction. We have

**THEOREM 2.** Let  $w^1, w^2, \dots, w^N$  be  $N$  mutually conjugate directions. Let

$$x^{k+1} = x^k - \lambda^k w^k$$

where  $\lambda^k$  minimizes

$$\phi(\lambda) = J(x^k - \lambda w^k), \quad \lambda \in \mathbb{R}.$$

56 When  $x^1 \in \mathbb{R}^N$  is given, we have

$$x^{N+1} = x^*$$

where

$$Ax^* = b.$$

*Proof.* Let

$$r^n = -J'(x^n) = b - Ax^n.$$

Since  $\lambda^k$  minimizes  $\phi(\lambda)$  we have

$$\phi(\lambda^k) = (J'(x^k - \lambda^k w^k), -w^k) = 0.$$

This gives

$$\lambda^k = \frac{(r^k)^T w^k}{(w^k)^T A w^k} \quad (4.8)$$

Since  $w^1, w^2, \dots, w^N$  are mutually conjugate directions, they are linearly independent. Therefore there exist  $\alpha_i$ ,  $1 \leq i \leq N$ , such that

$$x^1 - x^* = \sum_{k=1}^N \alpha_k w^k.$$

From this, using the fact that  $w^j$  are mutually conjugate, we obtain

$$(x^1 - x^*)^T A w^j = \alpha_j (w^j)^T A w^j.$$

This gives

$$\alpha_j = \frac{(x^1 - x^*)^T A w^j}{(w^j)^T A w^j}. \quad (4.9)$$

Using induction we show that

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$$\alpha_k = \lambda^k.$$

Since  $Ax^* = b$ , we have

$$r^1 = Ax^1 - b = A(x^1 - x^*).$$

This shows that

$$\alpha_1 = \lambda^1.$$

Let  $\alpha_i = \lambda^i$  for  $1 \leq i \leq k-1$ .

From the definition of  $x^k$  we obtain

$$x^k = x^1 - \sum_{i=1}^{k-1} \lambda^i w^i = x^1 - \sum_{i=1}^{k-1} \alpha_i w^i,$$

(by induction hypothesis). Since

$$(w^i)^T Aw^k = 0 \quad \text{for } 1 \leq i \leq k-1,$$

we get

$$(x^k - x^1)^T Aw^k = 0$$

This together with (4.8) and (4.9) shows that

$$\alpha_k = \lambda^k$$

Thus  $\alpha_k = \lambda^k$  for  $1 \leq k \leq N$ .

The definition of  $x^k$  implies

$$x^{N+1} = x^1 - \sum_{i=1}^N \lambda^i w^i = x^1 - \sum_{i=1}^N \alpha_i w^i = x^*$$

□

### Algorithm for Conjugate Gradient Method

**58 THEOREM 3.** Let  $x_0 \in \mathbb{R}^N$ . Define  $w^1 = b - Ax^1$ . Knowing  $x^n$  and  $w^{n-1}$  we define  $x^{n+1}$  and  $w^n$  by

$$\begin{aligned} x^{n+1} &= x^n + \alpha_n w^n \\ w^n &= r^n + \beta_n w^{n-1}, \end{aligned}$$

where

$$r^n = b - Ax^n, \alpha_n = \frac{(r^n, w^n)}{(w^n, Aw^n)}, \beta_n = \frac{(r^n, r^n)}{(r^{n-1}, r^{n-1})}$$

Then  $w^n$  are mutually conjugate directions and  $x^{N+1}$  is the unique solution of  $Ax = b$ .

A proof of this theorem can be found in LUENBERGER [31]. It can be shown that

$$x^n - x^{N+1} \sim \left( \frac{1 - \sqrt{c}}{1 + \sqrt{c}} \right)^n,$$

where  $c = m/M$ . Thus the convergence rate in the conjugate gradient method is faster than in the steepest descent method, at least for quadratic functionals.



### 4.4 Computer Representation of a Triangulation

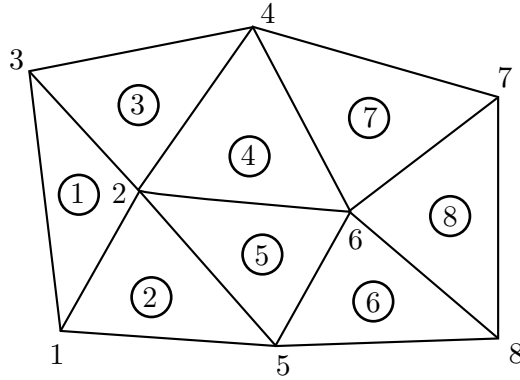


Figure 4.1:

Let  $T_h$  be a triangulation of the domain  $\Omega$ . We number the nodes of the triangulation and the triangles in  $T_h$ . Let

$$NS = \# \text{ of nodes of } T_h,$$

$$NT = \# \text{ of triangles in } T_h.$$

The triangulation is uniquely determined by the two matrices

$$Q(2, NS) = (q_{ij}) \quad \text{and} \quad ME(3, NS) = (m_{jk}),$$

where  $q_{ij}$  denotes the  $i^{\text{th}}$  coordinate of the  $j^{\text{th}}$  node and  $m_{jk}$  denotes the  $j^{\text{th}}$  vertex of the  $k^{\text{th}}$  triangle.

The matrix  $ME$ , corresponding to the triangulation in the above figure is

$$ME = \begin{pmatrix} 1 & 1 & 2 & 2 & 2 & 6 & 4 & 6 \\ 2 & 5 & 4 & 6 & 5 & 5 & 6 & 8 \\ 3 & 2 & 3 & 4 & 6 & 8 & 7 & 7 \end{pmatrix}$$

In some problems it is better to know the boundary nodes. The array  $NG(NS)$  defined by

$$NG(i) = \begin{cases} 1 & \text{if } i \in \Gamma \\ 0 & \text{otherwise} \end{cases}$$

is used for picking boundary nodes.

**Exercise 2.** Draw the triangulation

$$Q = \begin{pmatrix} 0 & 1 & 1 & 0 & 0.5 & 0.5 & 0.5 & 1 & 0 \\ 0 & 0 & 1 & 1 & 0.5 & 1 & 0 & 0 & 0.5 \end{pmatrix}$$

$$ME = \begin{pmatrix} 2 & 2 & 6 & 3 & 4 & 4 & 1 & 5 \\ 7 & 5 & 5 & 6 & 6 & 9 & 7 & 7 \\ 5 & 8 & 8 & 8 & 9 & 5 & 9 & 9 \end{pmatrix}$$

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## 4.5 Computation of the Gradient.

In the Neumann problem we have

$$(J'(u), v) = \int_{\Omega} (\nabla u \cdot \nabla v + a_0 uv) - \int_{\Omega} f v.$$

We now give a practical way of computing the gradient  $J'(u)$ .

Let  $w_i$  be the basis function in  $V_h$  which takes the value 1 at the  $i^{\text{th}}$  node and zero at the other nodes. Let  $V_h$  be defined by  $\mathbb{P}_1$  Lagrange finite element. Let

$$u = (u_i)_{1 \leq i \leq NS}$$

be given. We want to compute  $J'(u)$ . Let  $(J'(u))_i = (J'(u); w_i)$  and

$$B_i = \int_{\Omega} \nabla u \cdot \nabla w_i dx = \sum_{K \in T_h} \int_K \nabla u \cdot \nabla w_i dx.$$

We know that

$$w_i = \begin{cases} \lambda_j^k & \text{if } i = m_{jk}, \text{ for some } j \text{ and } k \\ 0 & \text{otherwise,} \end{cases}$$

where  $\lambda_j^k$  is the  $j^{\text{th}}$  barycentric coordinate of the  $k^{\text{th}}$  triangle. Therefore

$$B_i = \sum_{1 \leq j \leq 3, 1 \leq k \leq NT, i = m_{jk}} a_j^k,$$

$$a_j^k = \int_{K_k} \nabla u \cdot \nabla \lambda_j^k dx, \quad \text{where } K_k$$

is the triangle corresponding to the  $k^{\text{th}}$  element.

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### Algorithm to Compute B.

```

Set   B = 0
for   k = 1, 2, ..., NT;
for   j = 1, 2, 3,
do    Bmjk = Bmjk + ajk.

```

### Computation of the $a_j^k, s$ .

Since  $u = (u_i)_{1 \leq i \leq NS}$  and we take  $\mathbb{P}_1$  Lagrange finite elements, we have

$$\begin{aligned}
 u &= \sum_{j=1}^3 \lambda_j^k u_{m_{jk}} \quad \text{in } K_k, \\
 &= ax + by + c, \quad \text{say.}
 \end{aligned}$$

Let  $(\xi_i, \eta_i), 1 \leq i \leq 3$ , be the coordinates of the vertices of  $K_k$ . Let  $w_j = u_{m_{jk}}$ . Then  $a, b$  are found from the equations

$$\begin{aligned}
 a\xi_1 + b\eta_1 + c &= w_1, \\
 a\xi_2 + b\eta_2 + c &= w_2, \\
 a\xi_3 + b\eta_3 + c &= w_3,
 \end{aligned} \tag{4.10}$$

as

$$a = \frac{(w_1 - w_3)(\eta_2 - \eta_3) - (w_2 - w_3)(\eta_1 - \eta_3)}{C_2} \tag{4.11}$$

$$b = -\frac{(w_1 - w_3)(\xi_2 - \xi_3) - (w_2 - w_3)(\xi_1 - \xi_3)}{C_2} \tag{4.12}$$

where

$$C_2 = (\xi_1 - \xi_3)(\eta_2 - \eta_3) - (\xi_1 - \xi_2)(\eta_1 - \eta_3).$$

Hence  $\nabla u = \begin{pmatrix} a \\ b \end{pmatrix}$  is determined.

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If  $\lambda_j^k = a^j x + b^j y + c$ , then  $a^j$  and  $b^j$  are got from (4.11) and (4.12) by taking  $w_i = \delta_{ij}$ . Note that  $C_2 = 1/2 \text{ area } K_k$ . Then  $a_j^k = C_2/2 (aa_j + bb_j)$ .

## 4.6 Solution by Direct Methods

In Section 4.6.2 and 4.6.3 we gave algorithms to solve the equation  $Ax = b$  when  $A$  is symmetric and positive definite. When  $A$  is not symmetric or  $A$  is sparse, direct methods can be used to solve the equation  $Ax = b$ .

### 4.6.1 Review of the Properties of Gaussian Elimination

The principle of Gaussian elimination is to decompose  $A$  into a product  $LU$  where  $L$  is lower triangular and  $U$  is upper triangular so that the linear system

$$Ax = b$$

is reduced to solving 2 linear systems with triangular matrices

$$\begin{aligned} Ly &= b, \\ Ux &= y. \end{aligned}$$

Each of these is very easy to solve. For  $Ly = b, y = (y_i), y_1$  is given by the first equation, hence  $y_2$  is given by the second since  $y_1$  is already known, etc.

To decompose  $A$  into  $LU$ , one proceeds iteratively. Let

$$A^{(k)} = \left( a_{ij}^{(k)} \right)_{1 \leq i, j \leq N} \quad ,$$

63 be such that

$$a_{ij}^{(k)} = 0 \quad \text{for} \quad 1 \leq j \leq k-1 \quad \text{and} \quad i > j.$$

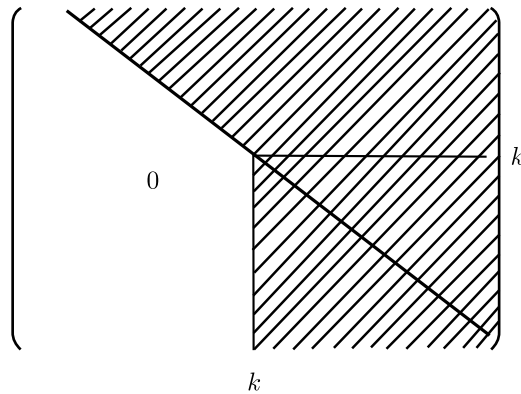


Figure 4.2:

Then  $U$  is  $A^{(N)}$ .

To get  $A^{(k+1)}$  from  $A^{(k)}$  one adds the  $k^{\text{th}}$  equation multiplied by a scaling factor of the  $i^{\text{th}}$  equation in order to have  $a_{ik}^{(k+1)} = 0$ . The scaling factor has to be

$$-\frac{a_{ik}^{(k)}}{a_{kk}^{(k)}}, \quad \text{and hence} \quad (4.13)$$

$$a_{ij}^{(k+1)} = a_{ij}^{(k)} - \frac{a_{ik}^{(k)} a_{kj}^{(k)}}{a_{kk}^{(k)}}, \quad j = k+1, \dots, n.$$

For a full matrix the order of operations for this process is  $N^3/3$ . For a band matrix, i.e. a matrix such that

$$a_{ij} = 0 \quad \text{if} \quad |i - j| \geq w,$$

We see that if  $|i - j| \geq w$ , then either  $|k - i|$  or  $|k - j|$  is greater than  $w + 1$ , provided that  $i, j \geq k + 1$ . Hence in the formula (4.13) an element which is outside the band is never modified since the corrective term in (4.13) is always zero.

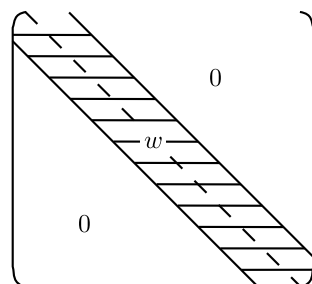


Figure 4.3:

64 More precisely we have

**PROPOSITION 4.** For a band matrix  $A$  with bandwidth  $w$ , at the  $k^{\text{th}}$  step of the Gaussian elimination, only  $w^2$  “corrective elements” have to be computed and added to the submatrix.

$$\begin{pmatrix} a_{k+1,k+1} & \cdots & a_{k+1,k+w} \\ \vdots & & \\ a_{k+w,k+1} & \cdots & a_{k+w,k+w} \end{pmatrix}$$

Note that we get an evaluation of the number of operations for the process which is about  $Nw^2$  (instead of  $N^3/3$ ).

#### 4.6.2 Stiffness Matrix and Stiffness Submatrix

For simplicity we consider the Neumann problem. Find  $u_h \in V_h$  such that

$$a(u_h, v) = L(v) \quad \forall v \in V_h,$$

where  $V_h$  is a finite dimensional subspace of  $H^1(\Omega)$  constituted with functions which are continuous and piecewise linear on the elements of the triangulation  $T_h$  and

$$a(u, v) = \int_{\Omega} (\nabla u \cdot \nabla v + uv) dx.$$

65 For  $(w_i)_{1 \leq i \leq N}$ , a basis of  $V_h$  (where  $N$  denotes the number of vertices of

$T_h$ ), we have

$$a(u, v) = \sum_{i,j=1}^N u_i a_{ij} v_j,$$

where  $v_i$  (respectively  $u_i$ ) denote the value of  $v$  (respectively  $u$ ) at the  $i^{\text{th}}$  vertex of  $T_h$  and

$$a_{ij} = \int_{\Omega} (\nabla w_i \cdot \nabla w_j + w_i w_j) dx.$$

But this is not a practical way to compute the elements  $a_{ij}$  of the matrix  $A$  of the linear system to be solved, since the support of the  $w_i$  involves several elements of  $T_h$ . Instead one writes

$$a(u, v) = \sum_{K \in T_h} \int_K (\nabla u \cdot \nabla v + uv) dx.$$

Hence, as

$$u(x) = \sum_{\alpha=1}^3 u_{m_{\alpha K}} \lambda_{\alpha}^K(x),$$

$$v(x) = \sum_{\beta=1}^3 v_{m_{\beta K}} \lambda_{\beta}^K(x),$$

in the element  $K$ , where  $(m_{\alpha K})_{\alpha=1,2,3}$  denotes the 3 vertices of the element  $K$  and  $\lambda_{\alpha}^K(x)$  the associated barycentric coordinates. One has

$$\sum_{i,j=1}^N u_i a_{ij} v_j = \sum_{K \in T_h} \sum_{\alpha,\beta=1}^3 u_{m_{\alpha K}} v_{m_{\beta K}} a_{\alpha\beta}^K,$$

where

$$a_{\alpha\beta}^K = \int_K (\nabla \lambda_{\alpha}^K \cdot \nabla \lambda_{\beta}^K + \lambda_{\alpha}^K \lambda_{\beta}^K) dx.$$

The matrix  $A^K = (a_{\alpha\beta}^K)_{1 \leq \alpha, \beta \leq 3}$  is called the *element stiffness matrix* of  $K$ .

A convenient algorithm to compute  $A$  is then the following *Assembling algorithm*.

$$\left\{ \begin{array}{l} 1. \text{ Set } A = 0 \\ 2. \text{ For } K \in T_h, \text{ compute } A^K \text{ and for } \alpha, \beta = 1, 2, 3 \text{ make} \\ \quad a_{m_{\alpha K}, m_{\beta K}} = a_{m_{\alpha K}, m_{\beta K}} + a_{\alpha\beta}^K. \end{array} \right.$$

**Exercise 3.** Write a Fortran subroutine performing the assembling algorithm (without the computation of the element stiffness matrices  $A^K$  which will be assumed to be computed in another subroutine).

### 4.6.3 Computation of Element Stiffness Matrices

We shall consider more sophisticated elements, e.g. the triangular, quadratic element with 6 nodes.

The midside points have to be included in the numbering of the vertices to describe properly the triangulation. For each element  $K$ , one has to give the 6 numbers of its 6 nodes in the global numbering,

$$m_{\alpha K}, \alpha = 1, \dots, 6.$$

67

The assembling algorithm of last section is still valid except that  $\alpha$  and  $\beta$  range now from 1 to 6 and that  $\lambda_{\alpha}^K$  has to be replaced by  $p_{\alpha}^K$ ,  $\alpha = 1, 2, \dots, 6$ , the local basis functions of the interpolation (see chapter 3).

To compute the element stiffness matrix

$$A^K = (a_{\alpha\beta}^K)_{\alpha, \beta=1, \dots, 6},$$

one introduces the mapping,

$$F : \hat{K} \rightarrow K$$

where  $\hat{K}$  is the triangle  $(0, 0), (1, 0), (0, 1)$ . Since  $F$  is affine, we have

$$F(\xi) = B\xi + b,$$

where  $B$  is  $2 \times 2$  matrix and  $b \in \mathbb{R}^2$ .



Let  $\hat{u}(\xi) = u(F(\xi))$  and  $\hat{v}(\xi) = v(F(\xi))$ . One has

$$\int_K uv \, dx = \int_{\hat{K}} \hat{u}\hat{v} \, \det(B) \, d\xi \quad (4.14)$$

In the same way one has

$$\nabla \hat{u}|_{\xi} = B^T \nabla u|_{F(\xi)},$$

since  $B^T$  is the Jacobian matrix of  $F$ . Therefore,

$$\int_K \nabla u \cdot \nabla v \, dx = \int_{\hat{K}} (B^{-T} \nabla \hat{u}) \cdot (B^{-T} \nabla \hat{v}) \det(B) \, d\xi \quad (4.15)$$

Finally to compute the coefficients  $a_{\alpha\beta}^K$  of the element stiffness matrix  $A^K$ , one notices that

$$\hat{u}(\xi) = \sum_{\alpha=1}^6 u_{m_{\alpha K}} p_{\alpha}(\xi),$$

where  $(p_{\alpha}(\xi), \alpha = 1, \dots, 6)$  are the basis functions of  $K$  which are easily computed once for all. Note that **68**

$$\lambda_1 = 1 - \xi_1 - \xi_2, \quad \lambda_2 = \xi_1, \lambda_3 = \xi_2.$$

As  $p$  are polynomials (even for higher degree elements) the integrals in (4.14) and (4.15) can be computed by noticing that

$$\int_{\hat{K}} \xi_1^i \xi_2^j \, d\xi = \frac{i! \, j!}{(i+j+2)!}$$

However, for the simplicity of the programming they are usually computed by numerical integration: every integral of the type  $\int_K f(\xi) \, d\xi$  is replaced by

$$\sum_{\ell=1}^L w_{\ell} f(b_{\ell})$$

where  $(b_\ell)_{\ell=1,\dots,L}$  are called the nodes of the numerical integration formula and  $(w_\ell)_{\ell=1,\dots,L}$  the coefficients.

The programming is easier since one may compute (in view of (4.14) and (4.15) only the values of  $p_\alpha$  and  $\partial p_\alpha / \partial \xi_i$  at the points  $b_\ell$ . For more details and model programs we refer to Mercier – Pironneau [32].

## Chapter 5

# Review of the Error Estimates for the Finite Element Method

THE PURPOSE OF this chapter is to state the theorems on error estimates which are useful for our future analysis. The proof of the theorems can be found in CIARLET [9]. 69

**DEFINITION.** Let  $\Omega \subset \mathbb{R}^n$  be an open subset,  $m \geq 0$  be an integer and  $1 \leq p \leq +\infty$ . Then the Sobolev Space  $W^{m,p}(\Omega)$  is defined by

$$W^{m,p}(\Omega) = \{v \in L^p(\Omega) : \partial^\alpha v \in L^p(\Omega), \text{ for all } |\alpha| \leq m\}.$$

On the space  $W^{m,p}(\Omega)$  we define a norm  $\|\cdot\|_{m,p,\Omega}$  by

$$\|v\|_{m,p,\Omega} = \left( \int_{\Omega} \sum_{|\alpha| \leq m} |D^\alpha v|^p dx \right)^{1/p},$$

and a semi norm  $|\cdot|_{m,p,\Omega}$  by

$$|v|_{m,p,\Omega} = \left( \int_{\Omega} \sum_{|\alpha|=m} |D^\alpha v|^p dx \right)^{1/p}.$$

If  $k$  is an integer, then we consider the quotient space

$$W^{-k+1,p}(\Omega) = W^{k+1,p}(\Omega)/p_k(\Omega)$$

with the quotient norm

$$\|\tilde{v}\|_{k+1,p,\Omega} = \inf_{\ell \in p_k} \|v + \ell\|_{k+1,p,\Omega},$$

- 70 where  $\tilde{v}$  is the equivalence class containing  $v$ .  
We introduce a semi norm in  $\tilde{W}^{k+1,p}(\Omega)$  by

$$|\tilde{v}|_{k+1,p,\Omega} = |v|_{k+1,p,\Omega}.$$

Then we have

**THEOREM 1. (CIARLET - RAVIART).** In  $\tilde{W}^{k+1,p}(\Omega)$  the semi norm  $|\tilde{v}|_{k+1,p,\Omega}$  is a norm equivalent to the quotient norm  $\|v\|_{k+1,p,\Omega}$ .

Using this theorem it is easy to prove

**THEOREM 2.** Let  $W^{k+1,p}(\Omega)$  and  $W^{m,q}(\Omega)$  be such that  $W^{k+1,p}(\Omega) \hookrightarrow W^{m,q}(\Omega)$  (continuous injection). Let

$$\pi \in \mathcal{L}(W^{k+1,p}(\Omega), W^{m,q}(\Omega))$$

be such that for each  $p \in p_k$ ,  $\pi p = 0$ . Then there exists a  $c = c(\Omega, \pi)$  such that for each  $v \in W^{k+1,p}(\Omega)$

$$|v - \pi v|_{m,q,\Omega} \leq c |v|_{k+1,p,\Omega}.$$

**DEFINITION.** Two open subsets  $\hat{\Omega}, \Omega$  of  $\mathbb{R}^n$  are said to be affine equivalent if there exists an affine map  $F$  from  $\hat{\Omega}$  onto  $\Omega$  such that  $F(\hat{x}) = B\hat{x} + b$ , where  $B$  is a  $n \times n$  non singular matrix and  $b \in \mathbb{R}^n$ .

We have

- 71 **THEOREM 3.** Let  $\hat{\Omega}, \Omega$  be affine equivalent with  $F$  as their affine map.  
Then there exist constants  $\hat{c}, c$  such that for all  $v \in W^{m,p}(\Omega)$ ,

$$|\hat{v}|_{m,p,\hat{\Omega}} \leq c \|B\|^m |\det B|^{-1/p} |v|_{m,p,\Omega},$$

and for all  $\hat{v} \in W^{m,p}(\hat{\Omega})$ ,

$$|v|_{m,p,\Omega} \leq \hat{c} \|B^{-1}\|^m |\det B|^{1/p} |\hat{v}|_{m,p,\hat{\Omega}},$$

where

$$\hat{v} = v \bullet F.$$

If  $h$  (resp.  $\hat{h}$ ) is the diameter of  $\Omega$  (resp.  $\hat{\Omega}$ ) and  $p$  (resp.  $\hat{p}$ ) is the supremum of the diameters of all balls that can be inscribed in  $\Omega$  (resp.  $\hat{\Omega}$ ), then we have

**THEOREM 4.**  $\|B\| \leq h/\hat{p}$  and  $\|B^{-1}\| \leq \hat{h}/\rho$ .

**DEFINITION.** Two finite elements  $(\hat{K}, \hat{\Sigma}, \hat{P})$  and  $(K, \Sigma, P)$  are said to be affine equivalent if there exists an affine map  $F\hat{x} = B\hat{x} + b$  on  $\mathbb{R}^n$ , where  $B$  is an  $n \times n$  non singular matrix, and  $b \in \mathbb{R}^n$  such that

- (i)  $F(\hat{K}) = K$
- (ii)  $\hat{p} = \{\hat{p} = p \circ F : p \in P\}$ ,
- (iii)  $\hat{\Sigma} = \{\hat{\phi} = F^{-1} \circ \phi : \phi \in \Sigma\}$

where

$$F^{-1} \circ \phi(\hat{p}) = \phi(\hat{p} \circ F^{-1}).$$

**DEFINITION.** Let  $(K, \Sigma, P)$  be a finite element and  $v : K \rightarrow \mathbb{R}$  be a smooth function on  $K$ . Then by virtue of the  $P$ -unisolvency of  $\Sigma$  there exists a unique element, say,  $\pi_K v \in P$ , such that  $\phi(\pi_K v) = \phi(v)$  for all  $\phi \in \Sigma$ . The function  $\pi_K v$  is called the  $P$ -interpolate function of  $v$  and the operator  $\pi_K : C^\infty(K) \rightarrow P$  is called the  $P$ -interpolation operator. 72

Now we state an important theorem which is often used.

**THEOREM 5.** Let  $(\hat{K}, \hat{\Sigma}, \hat{P})$  be a finite element. Let  $s (= 0, 1, 2)$  be the maximal order of derivatives occurring in  $\Sigma$ . Assume that

- (i)  $W^{k+1,p}(\hat{K}) \hookrightarrow C^s(\hat{K})$ ,
- (ii)  $W^{k+1,p}(\hat{K}) \hookrightarrow W^{m,q}(\hat{K})$ ,

$$(iii) P_k \subset \hat{P} \subset W^{m,q}(\hat{K}),$$

Then there exists a constant  $C = C(\hat{K}, \hat{\Sigma}, \hat{P})$  such that for all affine equivalent finite elements  $(K, \Sigma, P)$  we have

$$|v - \pi_K v|_{m,q,K} \leq C (\text{meas } K)^{1/q-1/p} \frac{h_K^{k+1}}{\rho_K^m} |v|_{k+1,p,K}$$

for all  $v \in W^{k+1,p}(K)$ , where  $\pi_K$  is a  $P$ -interpolate operator,  $h_K$  is the diameter of  $K$  and  $\rho_K$  is the supremum of diameter of all balls inscribed in  $K$ .

**DEFINITION.** A family  $(T_h)$  of triangulations of  $\Omega$  is regular if

- 73 (i) for all  $h$  and for each  $K \in T_h$  the finite elements  $(K, \Sigma, P)$  are all affine equivalent to a single finite element  $(\hat{K}, \hat{\Sigma}, \hat{P})$ ;
- (ii) there exists a constant  $\sigma$  such that for all  $T_h$  and for each  $K \in T_h$  we have

$$\frac{h_K}{\rho_K} \leq \sigma;$$

- (iii) for a given triangulation  $T_h$ , if

$$h = \max_{K \in T_h} h_K,$$

then  $h \rightarrow 0$ .

**Exercise 1.** Prove that there exists a constant  $C$  independent of  $h$  such that

$$|p|_{1,k} \leq C/h |p|_{0,K} \quad \text{for all } p \in \mathbb{P}_k.$$

A theorem which gives a global error bound is the following.

**THEOREM 6.** Let us assume that

- (i)  $\pi_h : H^{k+1}(\Omega) \rightarrow V_h$ , the restriction of  $\pi_h$  to  $V_h$  being the identity,
- (ii)  $V_h \subset \pi_{K \in T_h} \mathbb{P}_k(K)$ ,

74 (iii)  $V_h \subset H^m(\Omega)$ ,

(iv)  $u \in H^m(\Omega)$  (regularity assumption),

Then we have

$$\|u - \pi_h u\|_{m,\Omega} \leq C h^{k+1-m} |u|_{k+1,\Omega},$$

where  $C$  is a constant independent of  $h$  and  $(T_h)$  is a regular family of triangulations.

For stating a theorem on  $L_2$  -error estimates we need the definition of a regular adjoint problem.

**DEFINITION.** Let  $V = H^1(\Omega)$  or  $H^1_\circ(\Omega)$ ,  $H = L^2(\Omega)$ . The adjoint problem:

$$\begin{cases} \text{Find } \phi \in V \text{ such that} \\ a(v, \phi) = (g, v) \text{ for all } v \in V, \end{cases}$$

is said to be regular if

(i) for all  $g \in L^2(\Omega)$ , the solution  $\phi$  of the adjoint problem for  $g$  belongs to  $H^2(\Omega) \cap V$ ;

(ii) there exists a constant  $C$  such that

$$\|\phi\|_{2,\Omega} \leq C |g|_{0,\Omega}.$$

We now have

**THEOREM 7.** Let  $(T_h)$  be a regular family of triangulations on  $\Omega$  with reference finite element  $(\hat{K}, \hat{\Sigma}, \hat{P})$ . Let  $s = 0$  and  $n \leq 3$ . Suppose there exists an integer  $k \geq 1$  such that  $u \in H^{k+1}(\Omega)$   $P_k \subset \hat{p} \subset H^1(\hat{K})$ . Assume further that the adjoint problem is regular. Then there exists a constant  $C$  independent of  $h$  such that 75

$$|u - u_h|_{0,\Omega} \leq C h^{k+1} |u|_{k+1,\Omega}.$$





## Chapter 6

# Problems with an Incompressibility Constraint

### 6.1 Introduction

We recall the variational formulation of the Stokes problem (see Chapter 76 2).

Find  $u \in V$  such that

$$a(u, v) = L(v) \quad \forall v \in V,$$

where

$$a(u, v) = \int_{\Omega} \nabla u \cdot \nabla v \, dx,$$
$$L(v) = \int_{\Omega} f \cdot v \, dx, \quad f \in (L^2(\Omega))^n,$$

and

$$V = \{v \in (H_0^1(\Omega))^n : \operatorname{div} v = 0\}.$$

It is difficult to construct internal approximations of  $V$  because of the constraint  $\operatorname{div} v = 0$ . In two dimensional problem we know that

$v \in V \Leftrightarrow$  there exists  $\psi \in H_0^2(\Omega)$  such that

$$v = \text{rot } \psi,$$

where

$$\text{rot } \psi = \left( \frac{\partial \psi}{\partial x_2}, -\frac{\partial \psi}{\partial x_1} \right)$$

Therefore, it seems logical that the difficulties we encountered in approximating  $H^2(\Omega)$  in a conforming way are transferred to the conforming approximation of  $V$ .

## 6.2 Approximation Via Finite Elements of Degree 1

77 Let  $W_h = \{v_h \in (C^0(\bar{\Omega}))^2 : v_h|_K \in (\mathbb{P}_1(K))^2, \text{ for } K \in T_h, v_h = 0 \text{ on } \partial\Omega\}$  It is natural to try for  $V_h$  the space

$$\{v_h \in W_h : \text{div } v_h = 0\}.$$

But for most triangulations,  $V_h = \{0\}$ . This is due to the fact that the number of equations due to the constraint  $\text{div } v_h = 0$  is greater than the number of degrees of freedom of  $W_h$ . In fact,

$$\text{Dimension of } W_h = 2 \quad (\# \text{ internal vertices}).$$

Number of equations due to the constraint  $\text{div } v_h = 0$  is equal to number of triangles.

Hence  $V_h$  cannot be a good approximation to  $V$ . However, if the triangulation  $T_h$  is obtained by first taking quadrilaterals and then dividing each quadrilateral into four triangles by joining the diagonals (see figure 6.1), we obtain a 'good space'  $V_h$ . In this case only 3 of the four equations  $\text{div } v_h = 0$  are independent.

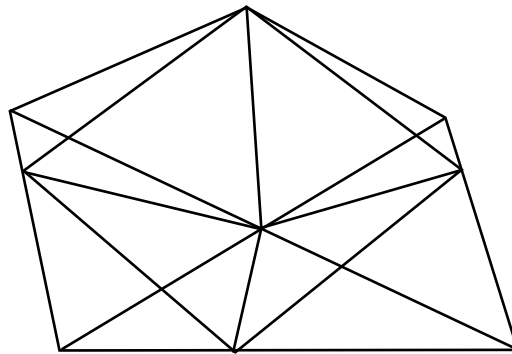


Figure 6.1:

**Exercise 1.** Let  $K$  be a quadrilateral. Let it be divided into four triangles  $T_i, 1 \leq i \leq 4$ , by

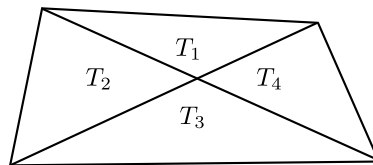


Figure 6.2:

joining the diagonals of  $K$ . Let  $v \in C^0(\bar{K})$  such that

$$v|_{T_i} \in \mathbb{P}_1(T_i), 1 \leq i \leq 4. \quad \text{Let } d_i = \text{div } v|_{T_i}, 1 \leq i \leq 4.$$

Then show that  $d_1 + d_3 = d_2 + d_4$ .

When the mesh is *uniform*, it is possible to prove convergence and to construct an interpolation operator,

$$\pi_h : V \rightarrow V_h$$

such that

$$\|v - \pi_h v\|_{1,\Omega} \leq ch \|v\|_{2,\Omega}.$$

Therefore the solution  $v_h$  of the approximate problem

$$a(u_h, v_h) = L(v_h) \quad \forall v_h \in V_h,$$

satisfies

$$\|u - u_h\|_{1,\Omega} \leq ch; \|u - u_h\|_{0,\Omega} \leq ch^2.$$

When the domain is a square and the mesh is uniform we define  $\pi_h$  as follows.

We choose for simplicity the diagonals of the square to be the coordinate axes.

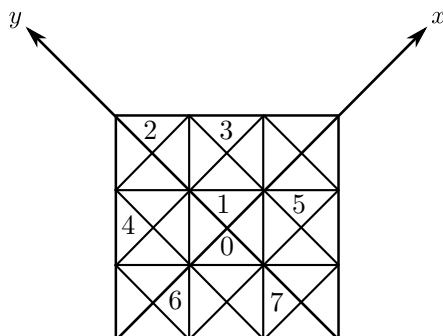


Figure 6.3:

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$\pi_h u$  at each of the main nodes (like 1, 5, 6, 7) is chosen as an average of  $u$ . For example, at the node 1 the two components of  $\pi_h u$  are given by

$$(\pi_h u)_1 = \text{average of } u_1 \text{ on } 02,$$

$$(\pi_h u)_2 = \text{average of } u_2 \text{ on } 34.$$

At each of the secondary nodes (like 0, 2, 3, 4)  $\pi_h u$  is chosen as an average of the values of  $\pi_h u$  at the main nodes. For example, at the node 0

$$(\pi_h u)_1(0) = \frac{(\pi_h u)_1(1) + (\pi_h u)_1(7)}{2},$$

$$(\pi_h u)_2(0) = \frac{(\pi_h u)_2(5) + (\pi_h u)_2(6)}{2},$$

Numerically the method works even for irregular, not too distorted, meshes.

### 6.3 The Fraeijs De Veubeke - Sander Element

We shall first describe a  $C^1$  element of a non standard type: the Fraeijs de Veubeke - Sander element.

Let  $K$  be a quadrilateral. We divide  $K$  into four triangles  $T_i$ ,  $1 \leq i \leq 4$ , by joining the diagonals of  $K$ .

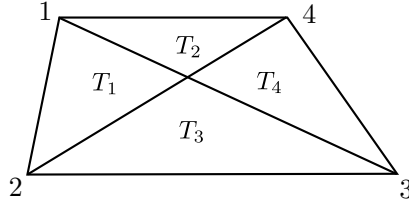


Figure 6.4:

Let

$$Q(K) = \{p \in C^1(K) : p \in \mathbb{P}_3(T_i), 1 \leq i \leq 4\}.$$

We have

**LEMMA 1.**  $\dim Q(K) = 16$ .

*Proof.* Indeed we choose  $p = p_1$  on  $T_1$  where  $p_1 \in \mathbb{P}_3$ . Then  $p_1$  depends on 10 parameters. Let  $p_2 \in \mathbb{P}_3$  be such that  $p_2$  and  $\partial p_2 / \partial n$  vanish on the diagonal 13. Hence  $p_2$  depends on  $10 - 4 - 3 = 3$  parameters. We choose  $p = p_1 + p_2$  on  $T_2$  and it is easy to see that  $p$  is  $C^1$  across 13.

In the same way we choose  $p = p_1 + p_3$  on  $T_3$  with  $p_3 = \partial p_3 / \partial n = 0$  on 24.  $p_3$  depends again on 3 parameters and  $p$  will be  $C^1$  across 24.

Taking  $p = p_1 + p_2 + p_3$  on  $T_4$ , it is easy to see that  $p$  is  $C^1$  across 13 and 24. Finally,  $p$  depends on  $10 + 3 + 3 = 16$  parameters. Hence the result.  $\square$

We choose,

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$$\Sigma_K = \left\{ \delta_i, \frac{\partial}{\partial x} \delta_i, \frac{\partial}{\partial y} \delta_i, 1 \leq i \leq 4; \frac{\partial}{\partial n} \delta_{b_i}, 1 \leq i \leq 4 \right\}$$

as the set of degrees of freedom. Here  $b_i$  denotes the midpoint of a side of the quadrilateral  $K$ .

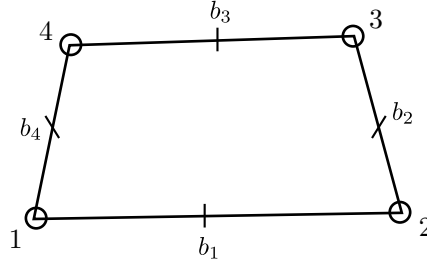


Figure 6.5:

We refer the reader to CIAVALDINI-NEDELEC [11] to prove that this is an admissible choice of degrees of freedom.

Let  $Q_h$  denote a regular family of *quadrangulations* of the polygonal domain  $\Omega$ . Let

$$\chi_h = \{\phi_h \in C^1(\bar{\Omega}) : \phi_h|_K \in Q(K), \phi_h = \frac{\partial \phi_h}{\partial n} = 0 \text{ on } \partial\Omega\}.$$

Then the following result has been proved by CIAVALDINI-NEDELEC [11].

**THEOREM 2.** *The operator  $\pi_h : H^3(\Omega) \rightarrow \chi_h$  defined by the above choice of degrees of freedom satisfies*

$$\|\phi - \pi_h \phi\|_{2,\Omega} \leq \text{ch}^s \|\phi\|_{s+2,\Omega},$$

where  $s = 1, 2$ .

- 82 As a consequence of this result the space  $\chi_h$  can be used to approximate any variational problem on  $H^2_0(\Omega)$  and in particular the biharmonic problem associated to the Stokes problem. However, as is the case with any Hermite finite element, the programming of the Fde V-S element is difficult and it is easier to use the corresponding element in the velocity formulation, which is quadratic.

## 6.4 Approximation of the Stokes Problem Via Quadratic Elements

Let  $T_h$  be the triangulation associated to  $Q_h$  by dividing each quadrilateral of  $Q_h$  into four triangles in the usual way. Let

$$W_h = \{v_h \in C^0(\bar{\Omega})^2 : v_h|_T \in (\mathbb{P}_2(T))^2, T \in T_h, v_h = 0 \text{ on } \partial\Omega\}$$

and

$$V_h = \{v_h \in W_h : \operatorname{div} v_h = 0\}.$$

Obviously  $v_h \in V_h$  implies that there exists  $\psi_h \in \chi_h$  such that  $v_h = \operatorname{rot} \psi_h$ . Therefore we can state

**THEOREM 3.** *There exists a constant  $c$  such that, for  $u \in V \cap H^{s+1}(\Omega)$  ( $s = 1$  or  $2$ ),*

$$\inf_{v_h \in V_h} \|u - v_h\|_{1,\Omega} \leq ch^s \|u\|_{s+1,\Omega}.$$

*Proof.* As  $u \in V$ , there exists  $\psi \in H^{s+2}(\Omega) \cap H_0^2(\Omega)$  such that  $u = \operatorname{rot} \psi$ . Then choosing  $v_h = \operatorname{rot} \pi_h \psi$  we obtain the result.  $\square$

Numerically, one may solve the approximate problem:

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$$\begin{cases} \text{find } u_h \in V_h \text{ such that} \\ a(u_h, v_h) = L(v_h) \quad \forall v_h \in V_h, \end{cases}$$

via a penalty method.

Let  $\varepsilon > 0$  and

$$a_\varepsilon(u, v) = a(u, v) + 1/\varepsilon \int_{\Omega} \operatorname{div} v \operatorname{div} u \, dx.$$

The penalised problem is:

$$\begin{cases} \text{Find } u_h \in W_h \text{ such that} \\ a_\varepsilon(u_h, v_h) = L(v_h) \quad \forall v_h \in W_h. \end{cases}$$

This problem is much easier to solve. We shall see in Chapter 7 that the order of the error due to penalization is only

$$O(\varepsilon) : \|u_{\varepsilon h} - u_h\|_1 \leq c\varepsilon,$$

where  $c$  may depend on  $h$ .

## 6.5 Penalty Methods.

We now come back to the case of finite elements of degree 1 where the space  $W_h$  is

$$W_h = \{v_h \in (C^\circ(\bar{\Omega}))^2 : v_h|_K \in (\mathbb{P}_1(K))^2, K \in T_h, v_h = 0 \text{ on } \partial\Omega\}.$$

The approximate problem can be taken as

$$\left\{ \begin{array}{l} \text{Find } u_{\varepsilon h} \in W_h \text{ such that} \\ a_\varepsilon(u_{\varepsilon h}, v_h) = L(v_h) \text{ for all } v_h \in W_h. \end{array} \right.$$

We have

**84 THEOREM 4.** *Assume that  $u \in H^2(\Omega)$ . Then there exists a constant  $c$  such that*

$$\|u_{\varepsilon h} - u\|_1^2 \leq c[h^2(1 + 1/\varepsilon) + h + \sqrt{\varepsilon}].$$

Hence by choosing  $\varepsilon = h^{4/3}$  we obtain

$$\|u_{\varepsilon h} - u\|_1 \leq ch^{1/3}.$$

**Exercise 2.** Prove the above theorem.

Note that the above convergence rate is very poor, which is confirmed by the poor numerical results obtained with this method.

## 6.6 The Navier-Stokes Equations.

The stationary flow of a viscous, Newtonian fluid subjected to gravity loads in a bounded region  $\Omega$  of  $\mathbb{R}^3$  is governed by the following dimensionless equations.

$$\left. \begin{array}{l} -\gamma\Delta u + \sum_{i=1}^3 u_i \frac{\partial u}{\partial x_i} + \nabla p = f \quad \text{in } \Omega, \\ \operatorname{div} u = 0 \quad \text{in } \Omega, \\ u = 0 \quad \text{in } \partial\Omega, \end{array} \right\} \quad (6.1)$$



where  $u$  represents the velocity,  $p$  the pressure and  $f$  is the body force. All these quantities are in dimensionless form and  $\gamma = \frac{\mu}{DV\rho} = \frac{1}{Re}$  where  $Re$  is called the Reynolds number. Here  $\mu$  is the viscosity of the fluid  $D$  a length characterizing the domain  $\Omega$ ,  $V$  a characteristic velocity of the flow and  $\rho$  the density of the fluid (For more details the reader is referred to BIRD-STEWART-LIGHTFOOR 'Transport Phenomena, Wiley Ed. p. 108).

The Reynolds number is the only parameter in the equation and it measures how far the Navier-Stokes model is from the Stokes model. The limiting case  $\gamma = 0$  corresponds to Euler's equations for inviscid fluids. However, at high Reynolds number, the flow develops a boundary layer near the boundary. Moreover, instability and bifurcation phenomena can be observed which correspond physically to turbulence. We are going to study only the flows at *low* Reynolds number.

## 6.7 Existence and Uniqueness of Solutions of Navier-Stokes

EQUATIONS AT LOW REYNOLDS NUMBERS. The variational formulation of Navier-Stokes equations is:

$$\begin{cases} \text{Find } u \in V \text{ such that} \\ a(u, v) + b(u, u, v) = (f, v) \quad \forall v \in V, \end{cases} \quad (6.2)$$

where

$$a(u, v) = \gamma \int_{\Omega} \nabla u \cdot \nabla v \, dx,$$

$$b(u, v, w) = \int_{\Omega} \sum_{i,j=1}^3 u_i \frac{\partial v_j}{\partial x_i} w_j \, dx$$

and

$$V = \{v \in (H_0^1(\Omega))^3 : \operatorname{div} v = 0\}.$$

**Exercise 3.** Show that if  $u$  is a solution of (6.2), then there exists a  $p \in L^2(\Omega)$  such that  $\{u, p\}$  is a solution of (6.1) in the sense of distributions. 86

**Exercise 4.** Show that for all  $u, v, w \in V$  one has

$$\begin{aligned} b(u, v, w) &= -b(u, w, v), \\ b(u, v, v) &= 0. \end{aligned}$$

As  $H^1(\Omega) \hookrightarrow L^4(\Omega)$  (see LADYZHENSKAYA [27] for a proof of this), using Schwarz inequality twice we obtain

$$b(u, v, w) \leq \|u\|_{0,4} \|v\|_1 \|w\|_{0,4} \leq c \|u\|_1 \|v\|_1 \|w\|_1.$$

Hence

$$\beta = \sup_{u, v, w \in V} \frac{b(u, v, w)}{\|u\|_1 \|v\|_1 \|w\|_1} < \infty.$$

where  $\|u\| = |u|_1$ . We have

**THEOREM 5.** Assume that  $\beta/\gamma^2 \|f\|_* < 1$ . Then the problem (6.2) has a unique solution.

*Proof.* Let  $u^i \in V, i = 1, 2$ . Let  $v^i, i = 1, 2$  be the solution

$$a(v^i, w) + b(u^i, v^i, w) = (f, w) \quad \forall w \in V \quad i = 1, 2 \quad (6.3)$$

Note that (6.3) is a linear problem and has a unique solution by virtue of the Lax-Milgram Lemma.

87 Choosing  $w = v^i$  in (6.3) it is easy to see that

$$\gamma \|v^i\|^2 \leq (f, v^i) \leq \|f\|_* \|v^i\|.$$

Thus

$$\|v^i\| \leq \frac{\|f\|_*}{\gamma}$$

Taking  $w = v^2 - v^1$ , we obtain

$$\begin{aligned} \gamma \|w\|^2 &\leq a(v^2 - v^1, w) = b(u^1, v^1, w) - b(u^2, v^2, w) = b(u^1, -w, w) + \\ &+ b(u^1 - u^2, v^2, w) \leq 0 + \beta \|v^2\| \cdot \|w\| \|u^1 - u^2\|. \end{aligned}$$

Hence we obtain

$$\|v^1 - v^2\| \leq \frac{\beta}{\gamma^2} \|f\|_* \|u^1 - u^2\|.$$

Since  $\frac{\beta}{\gamma^2} \|f\|_* < 1$ , the mapping  $T : u^i \rightarrow v^i$  is a strict contraction and has a fixed point which obviously is the unique solution of (6.2)  $\square$

**REMARK 1.** *The proof is constructive in the sense that the algorithm  $u^{n+1} = Tu^n$  gives a sequence which converges to the solution. At each step of this algorithm one has to solve the linearised problem (6.3).*

**REMARK 2.** *If  $\gamma^2 \geq \beta \|f\|_*$  then there exists atleast one solution to (6.2). The solution of (6.2) in this case may not be unique (see LIONS [28]). Note that problem (6.2) is equivalent to solving a non-linear equation  $F(u) = 0$  where  $F : V \rightarrow V'$  is given by* 88

$$F(u)(v) = a(u, v) + b(u, u, v) - (f, v).$$

Let  $G_u$  be the linear operator which is tangent to  $F$ , i.e.

$$(G_u w, v) = \lim_{\theta \rightarrow 0} \frac{1}{\theta} (F(u + \theta w) - F(u), v) = a(w, v) + b(u, w, v) + b(w, u, v).$$

If  $G_u$  is not singular, then  $u$  is an isolated solution, otherwise there may be a bifurcation. The eigenvalue problem associated to the linearised problem is:

$$\begin{cases} \text{Find } w \in V, \lambda \in \mathcal{C} \text{ such that} \\ a(w, v) + b(u, w, v) + b(w, u, v) = \lambda(w, v) \quad v \in V \end{cases}$$

and a study of this problem is of fundamental interest.

We refer the reader to BREZZI-RAPPAZ-RAVIART [6] for a study of the convergence in the case where  $u$  is an isolated solution. In the next section we restrict ourselves to the case where  $u$  is unique.

## 6.8 Error Estimates for Conforming Method

Let  $V_h \subset V$ . We consider the approximate problem:

$$\begin{cases} \text{Find } u_h \in V_h \text{ such that} \\ a(u_h, v_h) + b(u_h, u_h, v_h) = (f, v_h) \quad \forall v_h \in V_h \end{cases} \quad (6.4)$$

Let

$$\beta_h = \sup_{u, v, w \in V_h} \frac{b(u, v, w)}{\|u\| \|v\| \|w\|}$$

89 and

$$\|f\|_{h^*} = \sup_{v \in V_h} \frac{(f, v)}{\|v\|}.$$

Then it is easy to see that (6.4) has a unique solution when

$$\frac{\beta_h}{\gamma^2} \|f\|_{h^*} < 1.$$

The iterative method mentioned in Remark 1 converges for all  $\gamma$  satisfying  $\gamma^2 > \beta_h \|f\|_{h^*}$ . Note that  $\beta_h \leq \beta$  and  $\|f\|_{h^*} \leq \|f\|_*$ ; however, JAMET-RAVIART [23] proved that  $\beta_h \rightarrow \beta$  and  $\|f\|_{h^*} \rightarrow \|f\|_*$  as  $h \rightarrow 0$ .

**THEOREM 6.** *Assume that*

$$\frac{\beta}{\gamma^2} \|f\|_* < 1 - \delta$$

with  $0 < \delta < 1$ ; then one has

$$\|u - u_h\| \leq 3/\delta \|u - v_h\| \quad \forall v_h \in V_h.$$

*Proof.* Let  $w_h = v_h - u_h$ . Then

$$\begin{aligned} \gamma \|w_h\|^2 &\leq a(v_h - u, w_h) + a(u - u_h, w_h), \\ a(u, w_h) &= (f, w_h) - b(u, u, w_h), \\ a(u_h, w_h) &= (f, w_h) - b(u_h, u_h, w_h), \\ a(u - u_h, w_h) &= b(u_h, u_h - u, w_h) \end{aligned}$$

$$+b(u_h - u, u, w_h) \leq \beta \|u_h\| \|v_h - u\| \|w_h\| + \beta \|u_h - u\| \cdot \|u\| \|w_h\| \quad \text{since } b(u_h, u_h - v_h, w_h) = 0.$$

But we know that

$$\|u\| \leq 1/\gamma \|f\|_*, \|u_h\| \leq 1/\gamma \|f\|_*.$$

Therefore we obtain

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$$\|w_h\|^2 \leq (\|v_h - u\| + (1 - \delta) \|v_h - u\| + (1 - \delta) \|u_h - u\|) \cdot \|w_h\|.$$

As

$$\|u - u_h\| \leq \|u - v_h\| + \|w_h\|$$

we get

$$\delta \|w_h\|^2 \leq (3 - 2\delta) \|v_h - u\| \|w_h\|.$$

Hence

$$\|u - u_h\| \leq \|v_h - u\| + \frac{3 - 2\delta}{\delta} \|u - v_h\|.$$

This gives the desired result.  $\square$

An immediate consequence of Theorem 6 is that, when the solution  $u$  of (6.2) is sufficiently regular, we obtain the same error estimate for Navier-Stokes equations at low Reynolds number as for Stokes.

The method described in Section 6.4 is probably one of the best methods for Navier-Stokes equations at low Reynolds number.

However, when the Reynolds number is large, a major disadvantage is that the velocity field is required to be continuous (since our method is conforming) and this is not good to take into account boundary layer phenomena.

Indeed the velocity profile near the boundary has the behaviour as shown in figure. 91

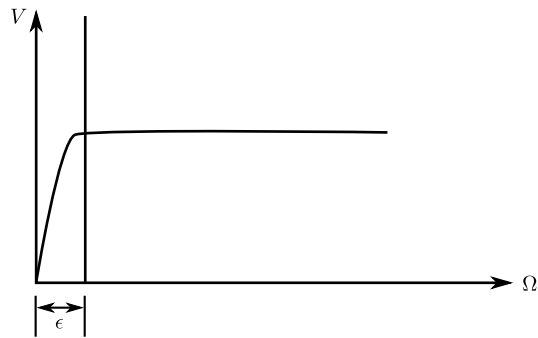


Figure 6.6:

$\epsilon$  is called the thickness of the boundary layer. In fact,  $\epsilon = \gamma^{1/2}$ , so that for high Reynolds number, this requires a very high refinement of the mesh near the boundary and therefore very expensive computer time. This can be partly avoided with mixed finite element since we shall work with discontinuous velocity fields.

## Chapter 7

# Mixed Finite Element Methods

### 7.1 The Abstract Continuous Problem

Let  $V, M, H$  be Hilbert spaces with  $V \hookrightarrow H$ . The continuous problem is: 92

Find  $\{u, \lambda\} \in V \times M$  such that

$$a(u, v) + b(v, \lambda) = (f, v) \quad \forall v \in V, \quad (7.1)$$

$$b(u, \mu) = (\phi, \mu) \quad \forall \mu \in M, \quad (7.1b)$$

where  $a(\cdot, \cdot) : H \times H \rightarrow \mathbb{R}$  and  $b(\cdot, \cdot) : V \times M \rightarrow \mathbb{R}$  are continuous bilinear forms and  $f \in V', \phi \in M'$ .

Let  $a(\cdot, \cdot)$  and  $b(\cdot, \cdot)$  satisfy

$$a(v, v) \geq \alpha \|v\|_H^2 \quad \forall v \in H \quad \text{ellipticity}, \quad (7.2)$$

$$\sup_{v \in V} \frac{b(v, \mu)}{\|v\|_V} \geq \beta \|\mu\|_M \quad \text{Brezzi's condition} \quad (7.3)$$

$|a(u, v)| \leq |a| \|u\| \|v\|, |b(v, \mu)| \leq |b| \|v\| \|\mu\|$ . We have

**THEOREM 1.** *If  $H = V$  then, under the above assumptions problem (7.1) has a unique solution.*

*Proof.* Let us consider the regularised problem:

$$a(u_\varepsilon, v) + b(v, \lambda_\varepsilon) = (f, v) \quad \forall v \in V, \quad (7.4)$$

$$-b(u_\varepsilon, \mu) + \varepsilon(\lambda_\varepsilon, \mu) = -(\phi, \mu) \quad \forall \mu \in M, \quad (7.4b)$$

Let  $\Phi = \{u, \lambda\}$ ,  $\Psi = \{v, \mu\} \in V \times M$ . We define

$$A_\varepsilon(\Phi, \Psi) = a(u, v) + b(v, \lambda) - b(u, \mu) + \varepsilon(\lambda, \mu), \quad L(\Psi) = (f, v) - (\phi, \mu)$$

- 93 Then  $A_\varepsilon(\Phi, \Psi)$  is  $V \times M$  coercive and  $L(\Psi)$  is a continuous, linear form on  $V \times M$ .

It is easy to see that problem (7.4) is equivalent to:

$$\left. \begin{array}{l} \text{Find } \Phi \in V \times M \text{ such that} \\ A_\varepsilon(\Phi, \Psi) = L(\Psi) \quad \forall \Psi \in V \times M \end{array} \right\} \quad (7.5)$$

By Lax-Milgram Lemma, problem (7.5) has a unique solution which implies that the regularized problem (7.4) has a unique solution.

Taking  $v = u_\varepsilon$  in (7.4),  $\mu = \lambda_\varepsilon$  in (7.4b) and adding and using the continuity of bilinear forms and  $H$ -ellipticity of  $a(\cdot, \cdot)$ , we get

$$\alpha \|u_\varepsilon\|^2 + \varepsilon \|\lambda_\varepsilon\|^2 \leq C(\|u_\varepsilon\| + \|\lambda_\varepsilon\|), \quad (7.6)$$

where  $C$  is a constant.

Since

$$b(v, \lambda_\varepsilon) = (f, v) - a(u_\varepsilon, v) \leq (\|f\|_* + |a| \|u_\varepsilon\|) \|v\|$$

we obtain, using Brezzi's condition,

$$\beta \|\lambda_\varepsilon\| \leq \sup_{v \in V} \frac{b(v, \lambda_\varepsilon)}{\|v\|} \leq \|f\|_* + |a| \|u_\varepsilon\|.$$

This implies

$$\|\lambda_\varepsilon\| \leq C(1 + \|u_\varepsilon\|). \quad (7.7)$$

From (7.6) and (7.7) we obtain

$$\|u_\varepsilon\| \leq C \quad \text{and} \quad \|\lambda_\varepsilon\| \leq C,$$



94 where  $C$  is a constant. Hence there exists a subsequence  $\{\epsilon'\}$ ,  $u \in V$ ,  $\lambda \in M$ , such that

$$u_{\epsilon'} \rightharpoonup u \quad \text{and} \quad \lambda_{\epsilon'} \rightharpoonup \lambda.$$

Obviously  $\{u, \lambda\}$  is a solution of (7.1).

If  $\{u_1, \lambda_1\}$  and  $\{u_2, \lambda_2\}$  are solutions of (7.1), then

$$\begin{aligned} a(u_1 - u_2, v) + b(v, \lambda_1 - \lambda_2) &= 0 \quad \forall v \in V, \\ b(u_1 - u_2, \mu) &= 0 \quad \forall \mu \in M. \end{aligned} \quad (7.8)$$

Taking  $v = u_1 - u_2$  and  $\mu = \lambda_1 - \lambda_2$ , we obtain

$$\alpha \|u_1 - u_2\|^2 \leq a(u_1 - u_2, u_1 - u_2) = 0.$$

Therefore

$$u_1 = u_2.$$

Since  $u_1 = u_2$ , using Brezzi's condition, we obtain from (7.8) that

$$\lambda_1 = \lambda_2.$$

Hence the solution of (7.1) is unique.  $\square$

**REMARK 1.** We now give an error estimate for the solution of (7.1) and the regularized problem (7.4). We have

$$a(u - u_\epsilon, v) + b(v, \lambda - \lambda_\epsilon) = 0 \quad \forall v \in V.$$

Hence

$$\beta \|\lambda - \lambda_\epsilon\| \leq |a| \cdot \|u - u_\epsilon\|. \quad (7.9)$$

From (7.1b) and (7.4b) we obtain

$$b(u - u_\epsilon, \mu) + \epsilon (\lambda_\epsilon, \mu) = 0 \quad \forall \mu \in M.$$

Choosing  $v = u - u_\epsilon, \mu = \lambda - \lambda_\epsilon$ , we get

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$$\alpha \|u - u_\epsilon\|^2 \leq a(u - u_\epsilon, u - u_\epsilon) = \epsilon (\lambda_\epsilon, \lambda - \lambda_\epsilon) \leq \epsilon \|\lambda_\epsilon\| \|\lambda - \lambda_\epsilon\|$$

Thus, using (7.7)

$$\|u - u_\epsilon\| \leq C \epsilon.$$

Thus from (7.9) we get

$$\|\lambda - \lambda_\epsilon\| \leq C \epsilon.$$

So we obtain

$$\|u - u_\epsilon\| = O(\epsilon) \quad \text{and} \quad \|\lambda - \lambda_\epsilon\| = O(\epsilon).$$

**REMARK 2.** Let  $B : V \rightarrow M$  be such that

$$(Bv, \mu)_M = b(v, \mu) \quad \forall \mu \in M.$$

One has, from (7.4b) and (7.4),

$$\begin{aligned} \lambda_\epsilon &= 1/\epsilon (Bu_\epsilon - \phi), \\ a(u_\epsilon, v) + 1/\epsilon (Bv, Bu_\epsilon - \phi) &= (f, v) \quad \forall v \in V, \end{aligned}$$

which correspond to penalization of (7.1b).

**REMARK 3.** We proved the existence and uniqueness of (7.1) only under the assumption  $V = H$ . If  $V \neq H$  then we do not have a general existence theorem. But existence theorems for particular examples when  $V \neq H$  are proved.

We now give examples of (7.1).

**96 EXAMPLE 1. The Stokes Problem.** We recall that the Stokes problem is:

Find  $\{u, p\}$  such that

$$\begin{aligned} -\gamma \Delta u + \nabla p &= f \quad \text{in } \Omega, \\ \operatorname{div} u &= 0 \quad \text{in } \Omega, \\ u &= 0 \quad \text{on } \Gamma. \end{aligned}$$

Without loss of generality we can take  $\gamma = 1$ . Using the standard technique of integration by parts we find that this corresponds to the problem:

Find

$$\{u, p\} \in (H_0^1(\Omega))^n \times (L^2(\Omega)/\mathbb{R})$$

such that

$$\begin{aligned} a(u, v) + b(v, p) &= (f, v) \quad \forall v \in (H_0^1(\Omega))^n, \\ b(u, \mu) &= 0 \quad \forall \mu \in L^2(\Omega)/\mathbb{R}; \end{aligned}$$

where

$$\begin{aligned} a(u, v) &= \int_{\Omega} \nabla u \cdot \nabla v \, dx, \\ b(v, \mu) &= - \int_{\Omega} \mu \operatorname{div} v \, dx. \end{aligned}$$

We take

$$\begin{aligned} V &= H = (H_0^1(\Omega))^n, \\ M &= (L^2(\Omega)/\mathbb{R}). \end{aligned}$$

Clearly  $a(\cdot, \cdot)$  is  $H$ -elliptic, continuous and bilinear. Let us prove Brezzi's condition,

$$\begin{aligned} \sup_{v \in V} \frac{b(v, \lambda)}{\|v\|_V} &= \sup_{v \in (H_0^1(\Omega))^n} \frac{- \int_{\Omega} \lambda \operatorname{div} v}{\|v\|_V} = \sup_{v \in (H_0^1(\Omega))^n} \frac{\langle \nabla \lambda, v \rangle}{\|v\|_V} \\ &= \|\nabla \lambda\|_{(H^{-1}(\Omega))^n} \geq \frac{1}{C} \|\lambda\|_{L^2(\Omega)/\mathbb{R}} \end{aligned}$$

where  $\langle \cdot, \cdot \rangle$  denotes the duality between  $(H_0^1(\Omega))^n$  and  $(H^{-1}(\Omega))^n$ . 97

Thus Stokes problem has a unique solution.

**EXAMPLE 2. The Biharmonic Problem.** Let

$$V = H^1(\Omega), \quad M = H_0^1(\Omega), \quad H = L^2(\Omega).$$

Consider the problem:

Find

$$\{u, \lambda\} \in H^1(\Omega) \times H_0^1(\Omega)$$

such that

$$a(u, v) + b(v, \lambda) = 0 \quad \forall v \in H^1(\Omega), \quad (7.10)$$

$$b(u, \mu) = \int_{\Omega} \phi \mu \, dx \quad \forall \mu \in H_0^1(\Omega). \quad (7.11)$$

where

$$a(u, v) = \int_{\Omega} uv \, dx, \quad b(v, \mu) = \int_{\Omega} \nabla v \cdot \nabla \mu \, dx$$

and  $\phi \in L^2(\Omega)$ .

Using integration by parts, we obtain from (7.10) that

$$\int_{\Omega} uv \, dx - \int_{\Omega} \Delta \lambda \cdot v \, dx + \int_{\Gamma} \frac{\partial \lambda}{\partial n} v \, d\Gamma = 0,$$

which implies

$$\left. \begin{aligned} u - \Delta \lambda &= 0 && \text{in } \Omega, \\ \frac{\partial \lambda}{\partial n} &= 0 && \text{on } \Gamma. \end{aligned} \right\} \quad (7.12)$$

98 From (7.11) we get

$$-\Delta u = \phi. \quad (7.13)$$

Thus we have, from (7.12) and (7.13), the biharmonic problem

$$\left\{ \begin{aligned} \Delta^2 \lambda &= 0 && \text{in } \Omega, \\ \frac{\partial \lambda}{\partial n} &= 0 && \text{on } \Gamma, \\ \lambda &= 0 && \text{on } \Gamma. \end{aligned} \right. \quad (7.14)$$

In the variational form of the biharmonic equation, we notice that  $V = H^1(\Omega) \neq L^2(\Omega) = H$ . It is easy to see when (7.10), (7.11) has one solution. If a solution  $\lambda$  of (7.14) is in  $H^3(\Omega) \cap H_0^2(\Omega)$  then  $u$  defined by (7.13) is in  $H^1(\Omega)$ . Moreover, one can check that this  $\{u, \lambda\}$  is a solution of (7.10), (7.11).

## 7.2 The Approximate Problem

Let  $V_h \subset V$  and  $M_h \subset M$  be two families of finite-dimensional spaces approximating  $V$  and  $M$ . We shall study the approximate problem:

Find

$$\{u_h, \lambda_h\} \in V_h \times M_h$$

such that

$$a(u_h, v_h) + b(v_h, \lambda_h) = (f, v_h) \quad \forall v_h \in V_h, \quad (7.15)$$

$$b(u_h, \mu_h) = (\phi, \mu_h) \quad \forall \mu_h \in M_h. \quad (7.15b)$$

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**Exercise 1.** Show that the problem (7.15) leads to solving a linear system with matrix

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

where  $A$  is a  $m \times m$  positive definite matrix and  $B$  is  $n \times m$  matrix where  $m = \dim V_h$  and  $n = \dim M_h$ .

Since  $V_h$  is finite-dimensional, the norms  $\|\cdot\|_V$  and  $\|\cdot\|_H$  on  $V_h$  are equivalent, that is there exists a function  $S(h)$  such that

$$\|v_h\|_V \leq S(h) \|v_h\|_H \quad \forall v_h \in V_h. \quad (7.16)$$

We introduce the affine spaces

$$Z_h(\phi) = \{v_h \in V_h : b(v_h, \mu_h) = (\phi, \mu_h) \quad \forall \mu_h \in M_h\}$$

$$Z(\phi) = \{v \in V : b(v, \mu) = (\phi, \mu) \quad \forall \mu \in M\}$$

**Exercise 2.** Let  $\phi = 0$ . Show that (7.1) is equivalent to the problem:

Find

$$\left. \begin{array}{l} u \in Z = Z(0) \quad \text{such that} \\ a(u, v) = (f, v) \quad \forall v \in Z \end{array} \right\} \quad (7.17)$$

In the same way show that (7.15) is equivalent to:

Find

$$\left\{ \begin{array}{l} u_h \in Z_h = Z_h(0) \quad \text{such that} \\ a(u_h, v_h) = (f, v_h) \quad \forall v_h \in Z_h \end{array} \right. \quad (7.18)$$

The present framework allows us to deal with  $Z_h \not\subset Z$  and hence we can consider non-conforming approximations of (7.17).

We now give an error bound in  $H$ -norm.

**THEOREM 2.** *Assuming that the continuous problem has at least one solution  $\{u, \lambda\}$  one has the error bound*

$$\|u - u_h\|_H \leq \left(1 + \frac{|a|}{\alpha}\right) \inf_{v_h \in Z_h(\phi)} \|u - v_h\|_H + \frac{|b|}{\alpha} S(h) \inf_{\mu_h \in M_h} \|\lambda - \mu_h\|_M \quad (7.19)$$

*Proof.* Let  $w_h = v_h - u_h$  and we have

$$a(w_h, w_h) = a(v_h - u, w_h) + a(u - u_h, w_h)$$

From (7.1) and (7.15), we obtain

$$a(u - u_h, w_h) = b(w_h, \lambda_h - \lambda) \quad \forall w_h \in V_h.$$

We notice that for  $v_h \in Z_h(\phi)$ ,

$$b(v_h - u_h, \mu_h) = 0 \quad \forall \mu_h \in M_h$$

and hence

$$a(u - u_h, v_h - u_h) = b(v_h - u_h, \mu_h - \lambda) \quad \forall v_h \in Z_h(\phi), \mu_h \in M_h.$$

Using the  $H$ -coerciveness of  $a(\cdot, \cdot)$  and the continuity of  $a(\cdot, \cdot)$  and  $b(\cdot, \cdot)$  we obtain

$$\begin{aligned} \alpha \|v_h - u_h\|_H^2 &\leq |a| \|v_h - u_h\|_H \|v_h - u\|_H + |b| \|\lambda - \mu_h\|_M \|v_h - u_h\|_V \\ &\leq |a| \|v_h - u_h\|_H \|v_h - u\|_H + |b| \|\lambda - \mu_h\|_M S(h) \|v_h - u_h\|_H \end{aligned}$$

101 Hence

$$\|v_h - u_h\|_H \leq \frac{|a|}{\alpha} \|v_h - u\|_H + \frac{|b|}{\alpha} S(h) \|\lambda - \mu_h\|_M.$$

We get the desired result by noticing that

$$\|u - u_h\|_H \leq \|u - v_h\|_H + \|v_h - u_h\|_H$$

□

**REMARK 4.** If  $Z_h(0) \subset Z(0)$  then the error estimate (7.19) reduces to

$$\|u - u_h\|_H \leq \left(1 + \frac{|a|}{\alpha}\right) \inf_{v_h \in Z_h(\phi)} \|u - v_h\|_H, \quad (7.20)$$

using the fact  $b(v_h - u_h, \lambda_h - \lambda) = 0$  as  $v_h - u_h \in Z_h(0) \subset Z(0)$ .

When  $\phi = 0$ , the error estimate (7.20) is obvious, in view of exercise 2. Then (7.20) is the error bound obtained in the conforming case.

### 7.3 Application to the Stokes Problem

In what follows  $T_h$  will denote a regular family of triangulations of the polygonal domain  $\Omega$ ,  $\vartheta_h$  will denote the set of vertices of  $T_h$ ,  $m_h$  the set of mid side points and  $\mathcal{E}_h$  the set of edges.

We consider the Stokes problem (example 1) where  $u$  is the velocity and  $\lambda$  is the pressure. 102

We shall choose for  $V_h$  a conforming  $\mathbb{P}_2$  space and for  $M_h$  a piecewise constant space, namely

$$V_h = \{v_h \in (C^0(\overline{\Omega}))^n : v_h|_K \in (\mathbb{P}_2(K))^n, K \in T_h, v_h = 0 \text{ on } \partial\Omega\}$$

and

$$M_h = \{\mu_h \in L^2(\Omega) : \mu_h|_K \in \mathbb{P}_0(K), K \in T_h\}.$$

We notice that

$\dim V_h = n$  (# internal vertices + # internal edges) and choose  $\Sigma_h$  for the set of degrees of freedom in each component where

$$\Sigma_h = \{\delta_N, N \in \vartheta_h; M_\gamma, \gamma \in \mathcal{E}_h\},$$

$M_\gamma(p) = \frac{1}{|\gamma|} \int_\gamma p ds$  denoting the average on the edge.

With this choice of  $\Sigma_h$ , the interpolation operator

$$\pi_h : (H^2(\Omega))^n \rightarrow V_h$$

defined by

$$\delta_N(\pi_h u) = \delta_N(u), \quad N \in \vartheta_h;$$

$$M_\gamma(\pi_h u) = M_\gamma(u), \quad \gamma \in \mathcal{E}_h,$$

will have some nice properties. Note that  $\pi_h$  is defined only on a subset of  $V$  since  $u(N)$  is not defined for all  $u$  in  $(H^1_0(\Omega))^n$ ;  $\pi_h$  is defined on  $V \cap (H^2(\Omega))^n$  since the functions in  $(H^2(\Omega))^n$  are continuous.

**Exercise 3.** Let  $K$  be a triangle.

Let  $P_K = \mathbb{P}_2(K)$  and

$$\Sigma_K = \{\delta_{a_i}, M_{\gamma_i}, 1 \leq i \leq 3\},$$

where  $a_i$  are the vertices of  $K$  and  $\gamma_i$  are edges of  $K$ .  $M_{\gamma_i}$  is defined by

$$M_{\gamma_i}(p) = \frac{1}{|\gamma_i|} \int_{\gamma_i} p \, ds.$$

Show that  $\Sigma_K$  is  $P_K$  unisolvent.

We have

**LEMMA 3.** One has

$$\int_K \operatorname{div}(\pi_h v) \, dx = \int_K \operatorname{div} v \, dx \quad \forall v \in (H^2(\Omega))^n$$

*Proof.* Indeed, by Green's formula,

$$\begin{aligned} \int_K \operatorname{div}(\pi_h v) \, dx &= \int_{\partial K} (\pi_h v \cdot n) \, ds \\ &= \int_{\partial K} v \cdot n \, ds \end{aligned}$$

since  $n$  is constant on each side of  $K$ .

Applying again Green's formula we get the desired result.  $\square$

**104 Error Estimates for  $\|u - u_h\|_1$**  If  $u \in H^2(\Omega)$  (which is true when  $\Omega$  is convex) then, since  $u \in Z(0)$  (i.e.  $\operatorname{div} u = 0$ ) and  $M_h$  contains functions



which are piecewise constant, we have, by Lemma 3,  $\pi_h u \in Z_h(0)$ . Hence we obtain

$$\inf_{v_h \in Z_h(0)} \|u - v_h\|_1 \leq \|u - \pi_h u\|_1 \leq ch \|u\|_2.$$

If the solution  $u \in H^3(\Omega)$  (which is unlikely since  $\Omega$  is a polygon) the error bound becomes  $ch^2 \|u\|_3$ . Indeed, the interpolation operator  $\pi_h$  leaves invariant the polynomial space  $\mathbb{P}_2(K)$  on each element  $K$  and the above error bound follows from Chapter 5.

However, we have

$$\inf_{\mu_h \in M_h} \|\lambda - \mu_h\|_0 \leq ch \|\lambda\|_1,$$

provided that the pressure  $\lambda \in H^1(\Omega)$ .

Finally, Theorem 2 gives

$$\|u - u_h\|_1 \leq ch(\|u\|_2 + \|\lambda\|_1),$$

which is only  $O(h)$  due to the low degree of approximation for the multiplier.

#### Other Choices for $V_h$ and $M_h$ .

Let

$$Q(K) = \mathbb{P}_2(K) + \{\lambda_1 \lambda_2 \lambda_3\},$$

where  $\lambda_1 \lambda_2 \lambda_3$  is called the *bubble* function. Note that the dimension of  $Q(K)$  is 7.

The choice

$$V_h = \{v_h \in (C^0(\bar{\Omega}))^2, v_h|_K \in (Q(K))^2, K \in T_h, v_h = 0 \text{ on } \partial\Omega\}$$

$$M_h = \prod_K \mathbb{P}_1(K)$$

leads to the error estimates

$$\|u - u_h\|_1 \leq ch^2, \quad \|u - u_h\|_0 \leq ch^3.$$

(See CROUZEIX - RAVIART [14].

The choice

$$V_h = \{v_h \in (C^0(\bar{\Omega}))^2 : v_h|_K \in (\mathbb{P}_2(K), K \in T_h, v_h = 0 \text{ on } \partial\Omega)\}$$

$$M_h = \{\mu_h \in C^0(\bar{\Omega}) : \mu_h|_K \in \mathbb{P}_1(K), K \in T_h\},$$

in which  $M_h$  contains continuous piecewise linear functions, leads to the same error estimates. (See BERCOVIER-PIRONNEAU [3]).

This last method, due to TAYLOR-HOOD [42], is widely used by engineers.

## 7.4 Dual Error Estimates for $u - u_h$

Let

$$V_2 \hookrightarrow V \hookrightarrow V_0$$

$$M_1 \hookrightarrow M.$$

- 106** We denote by  $\|\cdot\|_i$  the norms in  $V_i$  ( $i = 0, 2$ ) and  $\|\cdot\|_1$  the norm in  $M_1$ . We assume that  $V_0 \equiv V'_0$ . (In practical applications  $V_0$  will be a  $L^2$  space) and that  $H = V$  (i.e.  $a(\cdot, \cdot)$  is  $V$ -coercive).

Let  $g \in V_0$  and  $\{w, \psi\} \in V \times M$  satisfy

$$a(v, w) + b(v, \psi) = (g, v) \quad \forall v \in V, \quad (7.21)$$

$$b(w, \mu) = 0 \quad \forall \mu \in M. \quad (7.21b)$$

We assume (regularity result) that

$$\|w\|_2 + \|\psi\|_1 \leq c \|g\|_0 \quad (7.22)$$

and

$$\inf_{v_h \in Z_h(\phi)} \|w - v_h\|_V \leq e(h) \|w\|_2, \quad (7.23)$$

$$\inf_{\mu_h \in M_h} \|\mu - \mu_h\|_M \leq e(h) \|\psi\|_1. \quad (7.24)$$

A  $V_0$ -error estimate is given by

**THEOREM 4.** *Under the above assumptions, we have*

$$\|u - u_h\|_0 \leq ce(h) \left( \|u - v_h\|_V + \inf_{\mu_h \in M_h} \|\lambda - \mu_h\|_M \right).$$

*Proof.* We know that

$$\|u - u_h\|_0 = \sup_{g \in V_0} \frac{(g, u - u_h)}{\|g\|_0} \quad (7.25)$$

From (7.21) we have

$$(g, u - u_h) = a(u - u_h, w) + b(u - u_h, \psi) \quad (7.26)$$

Moreover, we have

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$$\begin{aligned} a(u - u_h, v_h) + b(v_h, \lambda - \lambda_h) &= 0 \quad \forall v_h \in V_h, \\ a(u - u_h, v_h) + b(v_h, \lambda - \mu_h) &= 0 \quad \forall v_h \in Z_h(\phi), \quad \forall \mu_h \in M_h, \end{aligned} \quad (7.27)$$

and

$$b(u - u_h, v_h) = 0 \quad \forall v_h \in M_h. \quad (7.28)$$

From (7.26), (7.21b), (7.27) and (7.28) we obtain

$$\begin{aligned} (g, u - u_h) &= a(u - u_h, w - v_h) + b(w - v_h, \lambda - \mu_h) + b(u - u_h, \psi - v_h) \\ &\quad \forall v_h \in Z_h(\phi), \\ &\quad \forall \mu_h, v_h \in M_h. \end{aligned}$$

Using the continuity of  $a(\cdot, \cdot)$  and  $b(\cdot, \cdot)$  we get

$$\begin{aligned} (g, u - u_h) &\leq c(\|u - u_h\|_V \|w - v_h\|_V + \|w - v_h\|_V \|\lambda - \mu_h\|_M + \\ &\quad + \|u - u_h\|_V \|\psi - v_h\|_M) \quad \forall v_h \in Z_h(\phi), \quad \forall \mu_h, v_h \in M_h. \end{aligned}$$

Taking the infimum over  $v_h \in Z_h(\phi)$  and  $\mu_h, v_h \in M_h$  and using (7.23), (7.24), we obtain

$$\begin{aligned} (g, u - u_h) &\leq c[\|w\|_2 \left( \|u - u_h\|_V + \inf_{\mu_h \in M_h} \|\lambda - \mu_h\|_M \right) + \\ &\quad + \|u - u_h\|_V \|\psi\|_1] e(h) \end{aligned}$$

Finally, using the regularity result (7.22) and (7.25) we get the desired result.  $\square$

**Application to Stokes Problem.**

We choose

$$V_0 = (L^2(\Omega))^n, V_2 = (H^2(\Omega))^n, M_1 = H^1(\Omega).$$

108 From the error estimates of section 7.3, we have

$$\inf_{v_h \in Z_h(0)} \|w - v_h\|_1 \leq ch \|w\|_2$$

and

$$\inf_{\mu_h \in M_h} \|\psi - \mu_h\|_0 \leq ch \|\psi\|_1.$$

The regularity result (7.22) is nothing but the regularity result for the Stokes problem.

Hence applying Theorem 4, we obtain  $L^2$ -error estimate.

$$\|u - u_h\|_0 \leq ch^2 (\|u\|_2 + \|\lambda\|_1).$$

## 7.5 Nonconforming Finite Element Method for Dirichlet Problem

We recall that the variational formulation of the Dirichlet problem

$$\left. \begin{array}{l} -\Delta u = f \quad \text{in } \Omega \\ u = 0 \quad \text{on } \partial\Omega \end{array} \right\} \quad (7.29)$$

is:

Find  $u \in H_0^1(\Omega)$  such that

$$a(u, v) = (f, v) \quad \forall v \in H_0^1(\Omega), \quad (7.30)$$

where

$$a(u, v) = \int_{\Omega} \nabla u \cdot \nabla v = \sum_{K \in T_h} \int_K \nabla u \cdot \nabla v,$$

$$(f, v) = \int_{\Omega} f v,$$

109 and  $T_h$  is a triangulation of  $\Omega$ .

We like to consider the nonconforming finite element approximation of (7.30), namely

Find  $u_h \in Z_h$  such that

$$a(u_h, v_h) = (f, v_h) \quad \forall v_h \in Z_h, \quad (7.31)$$

where

$Z_h = \{v_h : v_h|_K \in \mathbb{P}_1(K), K \in T_h, v_h \text{ is continuous}$

across the midside points of internal edges,  
 $v_h = 0 \text{ on } \partial\Omega\}$ .

Notice that  $Z_h \not\subset H^1_0(\Omega)$ .

We will construct a mixed finite element which is equivalent to (7.31). Multiplying the first equation in (7.29) by  $v \in \Pi^1_K(K)$  and integrating we obtain, using integration by parts in each triangle  $K$ ,

$$\sum_K \int_K \nabla u \cdot \nabla v - \sum_K \int_{\partial K} (\nabla u \cdot n) v = \int_{\Omega} f v.$$

This suggests

$$a(u, v) = \sum_K \int_K \nabla u \cdot \nabla v, \quad (7.32)$$

$$b(v, \mu) = - \sum_K \int_K (\mu \cdot n) v, \quad (7.33)$$

where  $\mu$  belongs to some *suitable space*.

We have to construct finite-dimensional subspaces  $V_h$  and  $M_h$  such that the problem

$$\left. \begin{array}{l} \text{Find } \{u_h, \lambda_h\} \in V_h \times M_h \text{ with} \\ a(u_h, v_h) + b(v_h, \lambda_h) = (f, v_h) \quad \forall v_h \in V_h \\ b(u_h, \mu_h) = 0 \quad \forall \mu_h \in M_h \end{array} \right\} \quad (7.34)$$

is equivalent to (7.31). Here  $a(\cdot, \cdot), b(\cdot, \cdot)$  are as in (7.32) and (7.33). 110

We take  $V_h = \prod_K \mathbb{P}_1(K)$ .

It is easy to see that if  $(\mu_h, n)$  is constant and continuous along internal edges, then  $b(v_h, \mu_h) = 0 \quad \forall v_h \in Z_h$ .

Define

$$Q(K) \subset (\mathbb{P}_1(K))^2$$

by

$$Q(K) = \{q = (q_1, q_2) : q_1 = a + bx, q_2 = c + by\}.$$

If  $\alpha x + \beta y = \ell$  is the equation of an edge  $\gamma$  then  $q \cdot n$  is constant on  $\gamma$  for  $q \in Q(K)$ , where  $n$  is normal to  $\gamma$ . The set

$\sum_K = \{(q \cdot n)(a_{ij}) : a_{ij} \text{ are mid points of the sides of } K\}$  is  $Q(K)$ -unisolvent. Hence

$$\begin{aligned} M_h &= \{q \in (L^2(\Omega))^2 : q|_K \in Q(K), K \in T_h, \operatorname{div} q \in L^2(\Omega)\} \\ &= \{q \in (L^2)^2 : q|_K \in Q(K), K \in T_h, q \cdot n \\ &\quad \text{is continuous across the edges of } T_h\} \end{aligned}$$

serves our purpose.

- 111 **Exercise 4.** With the above constructed  $V_h$  and  $M_h$  show that (7.31) and (7.34) are equivalent. Further show that  $Z_h(0) = Z_h$ .

$$\text{(Recall } Z_h(0) = \{v_h \in V_h : b(v_h, \mu_h) = 0 \quad \forall \mu_h \in M_h\})$$

The continuous problem corresponding to (7.34) can be obtained as follows:

It is natural to take

$$V = \prod_K H^1(K).$$

When  $\mu$  is smooth we can write

$$\begin{aligned} b(v, \mu) &= - \sum_K \int_{\partial K} (\mu \cdot n) v \\ &= - \sum_K \left( \int_K \operatorname{div} \mu \cdot v + \int_K \mu \cdot \nabla v \right) \end{aligned}$$

Hence we take

$$M = \{\mu \in (L^2(\Omega))^2 : \operatorname{div} \mu \in L^2(\Omega)\}.$$

Thus the continuous problem is:

$$\left. \begin{aligned} &\text{Find } \{u, \lambda\} \in V \times M \text{ such that} \\ &a(u, v) + b(v, \lambda) = (f, v) \quad \forall v \in V, \\ &b(u, \mu) = 0 \quad \forall \mu \in M, \end{aligned} \right\} \quad (7.35)$$

where

$$a(u, v) = \sum_K \int_K \nabla u \cdot \nabla v.$$

We have the characterisation:

**LEMMA 5.**

$$Z = \{v \in V : b(v, \mu) = 0 \quad \forall \mu \in M\} = H^1_0(\Omega).$$

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*Proof.* Let  $v \in \mathcal{D}(\Omega)$ . Then

$$\begin{aligned} b(v, \mu) &= - \sum_K \left( \int_K \operatorname{div} \mu \cdot v + \int_K \mu \cdot \nabla v \right) \\ &= - \left( \int_{\Omega} \operatorname{div} \mu \cdot v + \int_{\Omega} \mu \cdot \nabla v \right) \\ &= - \langle \operatorname{div} \mu, v \rangle - \langle \mu, \nabla v \rangle \\ &= - \langle \operatorname{div} \mu, v \rangle + \langle \operatorname{div} \mu, v \rangle \\ &= 0. \end{aligned}$$

Since  $b(\cdot, \cdot)$  is continuous on  $V \times M$  and  $\mathcal{D}(\Omega)$  is dense in  $H^1_0(\Omega)$  in the  $\|\cdot\|_1$  norm topology, we obtain

$$H^1_0(\Omega) \subset Z$$

We have to prove the other inclusion. Let  $v \in Z$ . Define

$$v_i \text{ by } v_i|_K = \frac{\partial}{\partial x_i}(v|_K), \forall K \in T_h.$$

Then  $v_i \in L^2(\Omega)$ .

Let  $\phi \in \mathcal{D}(\Omega)$ . Then

$$\begin{aligned} \langle v_1, \phi \rangle &= \sum_K \int_K v_1 \phi = \sum_K \int_K \frac{\partial v}{\partial x_1} \phi \\ &= \sum_K \left( - \int_K v \frac{\partial \phi}{\partial x_1} + \int_{\partial K} v \phi n_1 \right) \\ &= - \left\langle v, \frac{\partial \phi}{\partial x_1} \right\rangle + \sum_K \int_{\partial K} v \phi n_1. \end{aligned}$$

113 Since  $v \in Z$ ,  $b(v, \mu) = 0 \forall \mu \in M$ . Taking  $\mu = (\phi, 0)$ , we obtain

$$\begin{aligned} 0 = b(v, \mu) &= - \left( \sum_K \int_K \operatorname{div} \mu \cdot v + \int_K \nabla v \cdot \mu \right) \\ &= - \sum_K \int_{\partial K} (\mu \cdot n) v \text{ since } \mu \text{ is smooth} \\ &= - \sum_K \int_{\partial K} \phi v n_1. \end{aligned}$$

Therefore,

$$\langle v_1, \phi \rangle = - \left\langle v, \frac{\partial \phi}{\partial x_1} \right\rangle.$$

Hence  $v_1 = \frac{\partial v}{\partial x_1}$  in  $\mathcal{D}'$ . Similarly we have  $v_2 = \frac{\partial v}{\partial x_2}$  in  $\mathcal{D}'$ . Therefore  $v \in H^1(\Omega)$ . Further

$$0 = b(v, \mu) = - \sum_K \int_{\partial K} (\mu \cdot n) v \forall \mu \in (H^1(\Omega))^2 \subset M$$



This implies  $v = 0$  on  $\partial\Omega$ . Thus  $v \in H_0^1(\Omega)$ . Hence  $Z \subset H_0^1(\Omega)$ .

If  $\{u, \lambda\}$  is a solution of (7.35) then Lemma 5 and the second equation in (7.35) imply that  $u \in H_0^1(\Omega)$ . The first equation in (7.35) gives that  $-\Delta u = f$  in  $\mathcal{D}'$ . Thus if  $\{u, \lambda\}$  is a solution of (7.35), then  $u$  is the solution of the Standard Dirichlet problem:

$$\begin{cases} -\Delta u = f & \text{in } \Omega \\ u = 0 & \text{on } \partial\Omega, \end{cases} \quad (7.36)$$

and  $\lambda$  is the Lagrange multiplier for the continuity constraint for  $u$  from one element to the other; (7.35) is called the *primal hybrid formulation* of the Dirichlet problem. 114

If  $u$  is the solution of the Dirichlet problem (7.36) then  $\{u, -\nabla u\}$  is a solution of (7.35). Note that when  $\lambda$  is smooth then  $b(\cdot, \cdot)$  contains only the trace of  $\lambda.n$  on the edges  $\gamma$ , so that  $\lambda$  is certainly not unique and the Brezzi condition does not hold.  $\square$

#### Error Estimates for $u - u_h$ .

We notice that the interpolation operator

$$\pi_h : H^2(\Omega) \rightarrow Z_h$$

such that  $\pi_h u = u$  at the mid side points of  $T_h$  satisfies

$$\|u - \pi_h u\|_{1,K} \leq ch \|u\|_{2,K}.$$

Therefore

$$\|u - \pi_h u\|_V \leq ch \|u\|_2. \quad (7.37)$$

On the other hand, we define

$$\pi'_h : (H^1(\Omega))^2 \rightarrow M_h$$

on  $\gamma$  by

$$(\pi'_h \mu).n = \frac{1}{|\gamma|} \int_{\gamma} \lambda.n$$

We now state a theorem whose proof can be found in RAVIART-THOMAS [38].

**THEOREM 6.** *There exists a constant  $c$  such that*

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$$\| \mu - \pi'_h \mu \|_M \leq ch (\| \mu \|_1 + \| \operatorname{div} \mu \|_1), \quad (7.38)$$

$$\| \pi'_h \mu \|_M \leq c \| \mu \|_1 .$$

*Using Theorem 2, (7.37) and (7.38) we obtain*

$$\| u - u_h \|_V \leq ch \| u \|_2 . \quad (7.39)$$

*The dual error estimate of section 7.4 together with (7.39) gives*

$$\| u - u_h \|_0 \leq ch^2 \| u \|_2 \quad (7.40)$$

**REMARK 5.** *In practice, one solves the approximate problem;*

*Find  $u_h \in Z_h$  such that*

$$a(u_h, v_h) = (f, v_h) \quad \forall v_h \in Z_h,$$

*since  $Z_h$  is a simple finite element space. The fact that  $Z_h \not\subset H^1_0(\Omega)$  has no importance for practical purposes. The same  $Z_h$  can be used to approximate the Stokes problem. One chooses  $V_h = (Z_h)^2$  (for the velocities) and piecewise constant pressure; i.e.  $M_h$  as in section 7.3. For the error analysis we refer the reader to CROUZEIX-RAVIART [14].*

**REMARK 6.** *Other primal hybrid finite elements can be obtained with other choices for  $V_h \in \mathcal{E}M_h$ , some of them giving other well known non-conforming finite elements. (See RAVIART-THOMAS [37]).*

## 7.6 Approximate Brezzi Condition

116 We assume that the approximate Brezzi condition

$$\sup_{v_h \in V_h} \frac{b(v_h, \mu_h)}{\| v_h \|_V} \geq \gamma \| \mu_h \|_M \quad \forall \mu_h \in M_h, \quad (7.41)$$

where  $\gamma$  is independent of  $h$ , holds. The approximate Brezzi condition guarantees a unique solution for the approximate problem. Notice that the continuous Brezzi condition need not imply the approximate Brezzi condition.

We have the following result:

**THEOREM 7.** *Under the assumption (7.41) there exists a constant  $c_1$  such that*

$$\inf_{v_h \in V_h(\phi)} \|u - v_h\|_V \leq c_1 \inf_{w_h \in V_h} \|u - w_h\|_V \quad (7.42)$$

*Proof.* We will show that for each  $w_h \in V_h$  there exists a  $v_h \in Z_h(\phi)$  such that

$$\|u - v_h\| \leq c \|u - w_h\|,$$

where  $c$  is a constant. This will imply (7.42).

Let  $w_h \in V_h$ . Let  $\{y_h, v_h\} \in V_h \times M_h$  be the solution of

$$\begin{aligned} (y_h, v_h)_V + b(v_h, v_h) &= 0 \quad \forall v_h \in V_h, \\ b(y_h, \mu_h) &= b(u - w_h, \mu_h) \quad \forall \mu_h \in M_h. \end{aligned}$$

Using the approximate Brezzi condition and the continuity of  $b(\cdot, \cdot)$ , 117 it is easy to prove that

$$\begin{aligned} \|v_h\|_M &\leq 1/\gamma \|y_h\|_V, \\ \|y_h\|_V &\leq |b|/\gamma \|u - w_h\|_V. \end{aligned}$$

Let  $v_h = y_h + w_h$ . Then

$$\begin{aligned} b(v_h, \mu_h) &= b(u - w_h, \mu_h) + b(w_h, \mu_h) \\ &= b(u, \mu_h). \end{aligned}$$

Hence  $v_h \in Z_h(\phi)$ . Further

$$\begin{aligned} \|u - v_h\|_V &\leq \|u - w_h\|_V + \|y_h\|_V \\ &\leq (1 + |b|/\gamma) \|u - w_h\|_V. \end{aligned}$$

Hence the theorem.  $\square$

We now give an error estimate for the multiplier.

**THEOREM 8.** *If (7.41) holds then one has the error estimate*

$$\|\lambda - \lambda_h\|_M \leq c \left( \|u - u_h\|_H + \inf_{\mu_h \in M_h} \|\lambda - \mu_h\|_M \right) \quad (7.43)$$

*Proof.* We have

$$\begin{aligned} b(v_h, \lambda_h - \mu_h) &= b(v_h, \lambda_h - \lambda) + b(v_h, \lambda - \mu_h) \\ &= a(u - u_h, v_h) + b(v_h, \lambda - \mu_h) \end{aligned}$$

From (7.41) we obtain

$$\begin{aligned} \gamma \|\lambda_h - \mu_h\|_M &\leq \sup_{v_h \in V_h} \frac{b(v_h, \lambda_h - \mu_h)}{\|v_h\|_V} \\ &\leq |a|c \|u - u_h\|_H + |b| \|\lambda - \mu_h\|_M. \end{aligned} \quad (7.44)$$

118 Further

$$\|\lambda - \lambda_h\|_M \leq \|\lambda - \mu_h\|_M + \|\lambda_h - \mu_h\|_M \quad (7.45)$$

From (7.44) and (7.45) we obtain (7.43).  $\square$

We now give a practical way of verifying the approximate Brezzi condition (7.41).

**LEMMA 9.** *If  $b(\cdot, \cdot)$  satisfies continuous Brezzi condition and if*

$$\Lambda_h : V \rightarrow V_h$$

*is such that*

$$b(v - \Lambda_h v, \mu_h) = 0 \quad \forall \mu_h \in M_h \quad (7.46)$$

*(in fact  $\Lambda_h$  maps  $Z(\phi)$  into  $Z_h(\phi)$ ) and*

$$\|\Lambda_h v\|_V \leq c \|v\|_V, \quad (7.47)$$

*then (7.41) holds with  $\gamma = \beta/c$ .*

**Exercise 5.** Prove Lemma 9.

For further details see FORTIN [18].

## 7.7 Dual Error Estimate for the Multiplier

**119** This section is an analogue to Section 7.4. We assume that  $V_2 \hookrightarrow V$  and  $M_3 \hookrightarrow M \hookrightarrow M_0$ . Further, we assume that  $M'_0 = M_0$  and that for  $g \in M_0$ , the problem:

Find  $\{w, \psi\} \in V \times M$  such that

$$a(v, w) + b(v, \psi) = 0 \quad \forall v \in V \quad (7.48)$$

$$b(w, \mu) = (g, \mu) \quad \forall \mu \in M \quad (7.48b)$$

has one solution such that

$$\|w\|_2 + \|\psi\|_3 \leq c \|g\|_0 \quad (7.49)$$

We also assume that

$$\inf_{v_h \in V_h} \|w - v_h\|_V \leq e(h) \|w\|_2, \quad (7.50)$$

$$\inf_{\mu_h \in M_h} \|\psi - \mu_h\|_M \leq \varepsilon(h) \|\psi\|_3. \quad (7.51)$$

Here  $\|\cdot\|_0$ ,  $\|\cdot\|_2$ ,  $\|\cdot\|_3$  denote the norms in  $M_0, V_2, M_3$  respectively. The dual error estimate is given by

**THEOREM 10.** *One has*

$$\begin{aligned} \|\lambda - \lambda_h\|_0 &\leq ce(h) (\|\lambda - \lambda_h\|_M + \|u - u_h\|_H) + \\ &+ c \varepsilon(h) \left( S(h) \|u - u_h\|_H + \inf_{y_h \in V_h} (\|u - y_h\|_V + S(h) \|u - y_h\|_H) \right) \end{aligned} \quad (7.52)$$

*Proof.* We have

$$\|\lambda - \lambda_h\|_0 = \sup_{g \in M_0} \frac{(g, \lambda - \lambda_h)}{\|g\|_0}$$

and

$$\begin{aligned} (g, \lambda - \lambda_h) &= b(w, \lambda - \lambda_h) \\ &= b(w - v_h, \lambda - \lambda_h) - a(u - u_h, v_h). \end{aligned}$$

where the above is obtained from (7.1) and (7.15). Using  $b(u - u_h, v_h) = 0 \quad \forall v_h \in M_h$  and (7.48), we obtain

$$(g, \lambda - \lambda_h) = b(w - v_h, \lambda - \lambda_h) + a(u - u_h, w - v_h) + b(u - u_h, \psi - v_h)$$

This implies,

$$(g, \lambda - \lambda_h) \leq c(e(h)(\|\lambda - \lambda_h\|_M + \|u - u_h\|_H) + \|u - u_h\|_V \varepsilon(h)) \|g\|_0.$$

Here (7.49) - (7.51) are used.

To estimate  $\|u - u_h\|_V$ , we remark that

$$\begin{aligned} \|u - u_h\|_V &\leq \|u - y_h\|_V + \|y_h - u_h\|_V, \\ \|y_h - u_h\|_V &\leq S(h) \|y_h - u_h\|_H \\ &\leq S(h)(\|y_h - u\|_H + \|u - u_h\|_H). \end{aligned}$$

Finally,

$$\begin{aligned} \|\lambda - \lambda_h\|_0 &\leq c[e(h)(\|\lambda - \lambda_h\|_M + \|u - u_h\|_H) + \\ &\varepsilon(h)(S(h) \|u - u_h\|_H + \inf_{y_h \in V_h} (\|u - y_h\|_V + S(h) \|u - y_h\|_H))] \end{aligned}$$

□

## 7.8 Application to Biharmonic Problem

- 121 We shall now study a finite element approximation to the biharmonic problem. (Example 2, Section 7.1). We recall that in the variational formulation of the biharmonic problem

$$\begin{aligned} \Delta^2 \lambda &= -\phi \quad \text{in } \Omega, \\ \lambda &= \frac{\partial \lambda}{\partial n} = 0 \quad \text{on } \partial\Omega, \end{aligned}$$

we have

$$V = H^1(\Omega), \quad M = H_0^1(\Omega), \quad H = L^2(\Omega),$$

$$a(u, v) = \int_{\Omega} uv \, dx, \quad b(v, \mu) = - \int_{\Omega} \nabla v \cdot \nabla \mu$$

and  $f = 0$ .

In the terminology of hydrodynamics  $\lambda = \psi$  is called the *stream function* and  $w = -\Delta\psi = u$  is called the *Vortex function*.

For both  $V_h$  and  $M_h$  we shall use standard Lagrange finite element space of degree  $k$ :

$$V_h = \{v_h \in C^0(\bar{\Omega}) : v_h|_K \in \mathbb{P}_k(K), \forall K \in T_h\}$$

and

$$M_h = \{\mu_h \in V_h : \mu_h = 0 \text{ on } \partial\Omega\}.$$

We immediately notice that the approximate Brezzi condition (7.41) holds. Indeed,

$$\sup_{v_h \in V_h} \frac{\int_{\Omega} \nabla v_h \cdot \nabla \mu_h}{\|v_h\|_1} \geq \frac{\int_{\Omega} |\nabla \mu_h|^2}{\|\mu_h\|_1} \geq \gamma \|\mu_h\|_1.$$

since  $M_h \subset V_h$ . Here  $\gamma$  is the constant occurring in Poincaré's inequality. **122**

Moreover, assuming that the triangulation satisfies the uniformity condition

$$h \leq c\rho_K, \quad \forall K \in T_h,$$

where  $\rho_K$  denotes the diameter of the largest ball included in  $K$  and  $c$  is a constant independent of  $h$ , one always has

$$\|v\|_1 \leq \frac{c}{\min_{K \in T_h} \rho_K} \|v_h\|_0.$$

Therefore one has the inverse inequality  $\|v_h\|_1 \leq \frac{c}{h} \|v_h\|_0$  which gives an evaluation of  $S(h)$ .

We can state a convergence result for  $k > 2$ .

**THEOREM 11.** *If  $\lambda \in H^{m+1}(\Omega)$  and  $u \in H^m(\Omega)$  then one has*

$$\|u - u_h\|_0 + \|\lambda - \lambda_h\|_1 \leq ch^{m-1} (\|u\|_m + \|\lambda\|_{m+1})$$

for  $m = 2, \dots, k$  and

$$\|\lambda - \lambda_h\|_0 \leq ch^m (\|u\|_m + \|\lambda\|_{m+1})$$

The above result is a consequence of Theorems 2, 7, 8, 10. We note that the result is not optimal since with polynomials of degree  $k$ , we should get an error bound in  $h^k$  for  $\|\lambda - \lambda_h\|_1$  provided that  $\lambda \in H^{k+1}$ .

- 123 These results have been recently improved by SCHOLTZ [41] who is able to give an error estimate in the case  $k = 1$ .

Note that the matrix  $A$  of the bilinear form on  $V_h$  has to be computed exactly in this case. (The use of a 1-point formula for the computation of  $\int_K uv \, dx$  leads to a diagonal matrix  $A$  but there is no convergence).

**REMARK 7.** Due to the inclusion  $M_h \subset V_h$ , the matrix  $B$  of the bilinear form  $b(\cdot, \cdot)$  on  $V_h \times M_h$  has the particular form  $B = (B_1, B_2)$  where  $B_1$  is the matrix of  $b(\cdot, \cdot)$  on  $M_h \times M_h$  and therefore  $B_1$  is invertible. The linear system corresponding to the approximate problem can be written as

$$\begin{pmatrix} A_1 & A_2 & B_1^T \\ A_2^T & A_3 & B_2^T \\ B_1 & B_2 & 0 \end{pmatrix} \begin{pmatrix} u_1 \\ u_2 \\ \lambda \end{pmatrix} = \begin{pmatrix} 0 \\ 0 \\ \phi \end{pmatrix}$$

One can eliminate  $u_1$  from the last equation ( $u_1 = -B_1^{-1} B_2 u_2$ ) and  $\lambda$  from the first one. This gives a linear system of equations in  $u_2$  which can be solved by any of the standard method.

The advantage is that the size of the linear system is  $p \times p$  where  $p$  is the number of boundary points which is relatively small.

## 7.9 General Numerical Methods for the Solution of the Approximate Problem

- 124 As we have noticed in Exercise 1, the approximate problem is equivalent to solving

$$\begin{pmatrix} A & B^T \\ B & 0 \end{pmatrix} \begin{pmatrix} u \\ \lambda \end{pmatrix} = \begin{pmatrix} f \\ \phi \end{pmatrix} \quad (7.53)$$

where the matrix on the left is invertible (provided that  $B$  has a maximal rank, i.e. the approximate Brezzi condition holds at least for  $\gamma$  dependent on  $h$ ) and symmetric if  $a(\cdot, \cdot)$  is symmetric but not positive.



Alternatively one can use the matrix

$$\begin{pmatrix} A & B^T \\ -B & 0 \end{pmatrix}$$

which is positive but not symmetric.

- a) **Solution of the Linear System (7.53) by Direct Methods.** The Gaussian elimination of (7.53) can be performed without pivoting. However, due to storage considerations, it is often much better to permute rows and columns to get a band matrix after a suitable ordering of the unknowns. (For finite elements, we order the unknowns following the ordering that we give to the associated nodes, with no distinction between  $\lambda_i$ 's and  $u_i$ 's). In this case a strategy of partial pivoting may be necessary.
- b) **Solution of the Approximate Problem by Penalty Methods.** 125  
First solve

$$(A + 1/\epsilon B^T B) u^\epsilon = f + 1/\epsilon B^T \phi,$$

and then find

$$\lambda^\epsilon = 1/\epsilon (Bu^\epsilon - \phi).$$

If  $B$  has maximal rank the error is only  $O(\epsilon)$  (see proof of Theorem 1), which is certainly small compared to the discretization error if  $\epsilon = 10^{-4}$  or  $10^{-6}$ .

However, the condition number of the matrix

$$A + 1/\epsilon B^T B$$

might be quite big and it is wise in such a case to use direct methods and then to compute explicitly the matrix  $B^T B$ , which is easy only if

$$M_h = \pi_K(\mathbb{P}_k(K))^d$$

for some positive  $k$  and  $d$ ; in other words, no continuity requirements has to be asked for  $\lambda_h$  between two elements. (Then  $B^T B$

can be computed by assembling some stiffness matrices). Otherwise the computation of  $B^T B$  is too costly. Note that this method is possible even if  $B$  has not a maximal rank (and the error is then only  $O(\sqrt{\epsilon})$ ).

- 126 c) **Solution of the Problem by Iterative Methods.** The conjugate gradient method has been successfully extended to the case of matrices such as (7.53) by PAIGE and SAUNDERS [35].

Another way of applying the conjugate gradient method is to notice that  $u$  can be eliminated from (7.53):

$$u = A^{-1} f - A^{-1} B^T \lambda.$$

Then we get

$$C \lambda = b$$

where

$$C = B A^{-1} B^T, \quad b = B A^{-1} f - \phi.$$

As  $C$  is symmetric and positive definite (if  $a(\cdot, \cdot)$  is symmetric) then the conjugate gradient method can be applied to the matrix  $C$ , but *one has never to compute the matrix  $C$  explicitly*. (This is too costly *unless*  $A$  is block diagonal and therefore  $V_h$  is a finite element space with no continuity requirements between 2 elements. Let  $A_K$  denote the element matrix of  $a(\cdot, \cdot)$  on  $K$  and  $B_K$  that of  $b(\cdot, \cdot)$ ; the matrix  $C$  is computed by assembling the  $B_K A_K^{-1} B_K^T$ 's. This is the case for hybrid elements).

Indeed what one needs for the conjugate gradient method is to be able to compute  $y = Cz$  for any column vector  $z$  and this is done in the following way:

$$\begin{aligned} \text{Compute } z_1 &= B^T ; \\ \text{Solve } A z_2 &= z_1 ; \\ \text{Compute } y &= B z_2 . \end{aligned}$$

- 127 Note that it is not necessary for  $B$  to have maximal rank.

## 7.10 Equilibrium Elements for the Dirichlet Problem

Let us consider the following problem:

Find  $\{u, \lambda\}$  such that

$$\begin{cases} u - \nabla \lambda = 0 & \text{in } \Omega \\ \operatorname{div} u = \phi & \text{in } \Omega \\ \lambda = 0 & \text{in } \partial\Omega \end{cases} \quad (7.54)$$

If  $\{u, \lambda\}$  is a solution of (7.54) then  $\lambda$  is the solution of the Standard Dirichlet Problem:

$$\begin{cases} \Delta \lambda = \phi & \text{in } \Omega \\ \lambda = 0 & \text{on } \partial\Omega \end{cases} \quad (7.55)$$

Multiplying (7.54) by  $v \in (L^2(\Omega))^2$  and using integration by parts we obtain an equivalent problem:

Find  $\{u, \lambda\} \in (L^2(\Omega))^2 \times H_0^1(\Omega)$  such that

$$\begin{cases} a(u, v) + b(v, \lambda) = 0 & \forall v \in (L^2(\Omega))^2, \\ b(u, \mu) = \int_{\Omega} \phi \mu & \forall \mu \in H_0^1(\Omega), \end{cases}$$

where

$$\begin{aligned} a(u, v) &= \int_{\Omega} u \cdot v \, dx, \\ b(v, \mu) &= - \int_{\Omega} v \cdot \nabla \mu \, dx. \end{aligned}$$

If  $M_h = C^\circ(\bar{\Omega}) \cap \pi_K^{\mathbb{P}_k}(K)$  then a natural choice for  $V_h$  is

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$$V_h = \pi_K^{\mathbb{P}_{k-1}}(K)^2.$$

Note that the operator  $\nabla$  maps  $M_h$  into  $V_h$ . This implies that the approximate Brezzi condition holds.

Problem (7.54) can be formulated in another way also.

Find  $\{u, \lambda\} \in H(\text{div}, \Omega) \times L^2(\Omega)$  such that

$$\int_{\Omega} uv + \int_{\Omega} \lambda \text{div } v = 0 \quad \forall v \in H(\text{div}, \Omega),$$

$$\int_{\Omega} \mu \text{div } u = (\mu, \phi) \quad \forall \mu \in L^2(\Omega),$$

where

$$H(\text{div}, \Omega) = \{v \in (L^2(\Omega))^2 : \text{div } v \in L^2(\Omega)\}.$$

We notice that

$$a(u, v) = \int_{\Omega} u \cdot v \, dx$$

is coercive on

$$H = (L^2(\Omega))^2.$$

To prove that  $b(\cdot, \cdot)$  satisfies Brezzi condition, we use the fact that if  $\mu \in L^2(\Omega)$  then the solution  $\psi$  of the Dirichlet problem:

$$\begin{aligned} \Delta \psi &= \mu \quad \text{in } \Omega \\ \psi &= 0 \quad \text{on } \partial\Omega, \end{aligned}$$

**129** satisfies  $\|\psi\|_1 \leq c \|\mu\|_0$ .

Therefore,

$$\sup_{v \in (L^2(\Omega))^2} \frac{\int_{\Omega} \mu \cdot \text{div } v}{\|v\|_V} \geq \frac{\int_{\Omega} \mu \text{div } (\nabla \psi)}{\|\nabla \psi\|_V} \geq \frac{\int_{\Omega} \mu^2 \, dx}{c \|\mu\|_0} = \frac{1}{c} \|\mu\|_0.$$

As approximate spaces, we choose

$$V_h = \{v \in V : v|_K \in Q(K), v.n \text{ is continuous across the sides of } T_h\}.$$

(See Section 7.5 for the definition of  $Q(K)$ ).

$$M_h = \prod_K \mathbb{P}_0(K).$$

As  $\text{div} : V_h \rightarrow M_h$ , we see that  $Z_h(0) \subset Z(0)$ . (Hence the name *equilibrium* elements:  $u_h$  will satisfy equilibrium equations  $\text{div } u_h = \phi$  for  $\phi$  piecewise constant).

We may then apply the error estimate derived in Section 7.2 and use the improvement given in Remark 4, since  $Z_h(0) \subset Z(0)$ . We shall choose  $v_h = \pi'_h v$  where  $\pi'_h$  is the interpolation operator defined in Section 7.5.

Indeed, we have

$$\int_{\gamma} (\pi'_h v) \cdot n \, ds = \int_{\gamma} v \cdot n \, ds \quad \text{for each edge } \gamma \text{ of } T_h.$$

Therefore,

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$$\begin{aligned} \int_{\Omega} \mu_h \text{div } \pi'_h v &= \sum_K \int_{\partial K} (\pi'_h v) \cdot n \mu_h \, d\Gamma \\ &= \sum_K \int_{\partial K} (v \cdot n) \mu_h \, d\Gamma = \int_{\Omega} \mu_h \text{div } v \, dx \quad \forall \mu_h \in M_h. \end{aligned}$$

This shows that  $v \in Z(\phi) \Rightarrow \pi'_h v \in Z_h(\phi)$ .

Finally we get

$$\|u - u_h\|_0 \leq c \|u - \pi'_h u\|_0 \leq ch,$$

where we have used the estimate given in Theorem 6.

To get an error estimate for  $\|\lambda - \lambda_h\|_0$ , we shall make use of the results in Section 7.6 and construct the operator  $\pi_h$  occurring in Lemma 9.

Let  $v$  in  $V$  be given and  $\phi$  satisfy

$$\begin{aligned} \Delta \phi &= \text{div } v \quad \text{in } \Omega, \\ \phi &= 0 \quad \text{on } \partial\Omega. \end{aligned}$$

Let

$$\Lambda : V \rightarrow (H^1(\Omega))^2$$

be defined by

$$\Lambda v = \nabla \phi.$$

We have (regularity result)

$$\|\Lambda v\|_1 \leq c \|\operatorname{div} v\|_0,$$

so that

$$\Lambda_h = \pi'_h \Lambda$$

131 satisfies

$$\|\Lambda_h v\|_M \leq c \|\Lambda v\|_1 \leq c \|\operatorname{div} v\|_0$$

and

$$b(\mu_h, \pi'_h \Lambda v) = b(\mu_h, \Lambda v) = b(\mu_h, v),$$

where we used the definitions of  $\pi'_h$  and  $\Lambda v$ .

Thus  $\Lambda_h$  satisfies the conditions required in Lemma 9. Hence we have

$$\|\lambda - \lambda_h\|_0 \leq ch,$$

by Theorem 8.

For further details about equilibrium elements the reader can refer the thesis of J.M. THOMAS, 1977.

**REMARK 8.** *If we replace  $Q(K)$  by  $(\mathbb{P}_1(K))^2$  we get a finite element with 6 degrees of freedom instead of 3 (2 values of v.n on each side). The interpolation operator  $\pi'_h$  is defined with the help of the degrees of freedom and has the same properties.*

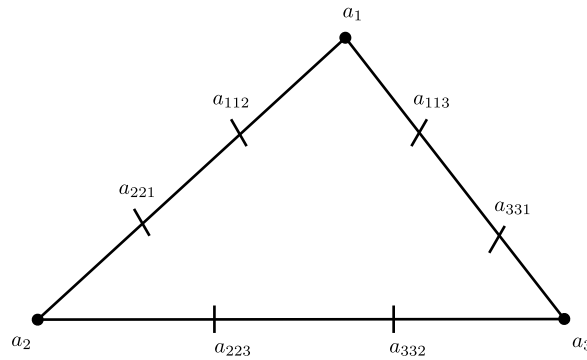


Figure 7.1:

In fact,  $\pi'_h$  is defined by

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$$\int_{\gamma} p(\pi'_h v) \cdot n \, ds = \int_{\gamma} p(v \cdot n) \, ds \quad \forall p \in \mathbb{P}_1(\gamma)$$

However, the error estimates now become

$$\begin{aligned} \|u - u_h\|_0 &\leq ch^2 \\ \|\lambda - \lambda_h\|_0 &\leq ch \end{aligned}$$

**REMARK 9.** The present finite element method can be extended to the elasticity equation where  $v$  represents the stress tensor  $\sigma_{ij}$ . The difficulty lies in the required symmetry of  $\sigma_{ij}$  but can be surmounted. See C. JOHNSON-B. MERCIER [25] and AMARA-THOMAS[2].

**REMARK 10. A posteriori Error Estimate.** Let us consider the following optimization problem:

$$\begin{aligned} &\text{Inf } J(v, \mu) \\ &v \in Z(\phi) \\ &\mu \in M \end{aligned}$$

where  $J(v, \mu) = 1/2 \|v - \nabla \mu\|^2$ . Clearly the optimal value is zero and corresponds to  $v = u$  and  $\mu = \lambda$  solution of (7.54). Since  $v \in Z(\phi)$ , we

also have

$$J(v, \mu) = 1/2 \|\nabla \mu\|^2 + (\phi, \mu) + 1/2 \|v\|^2 .$$

Since  $J(u, \lambda) \leq J(v_h, \lambda) \forall v_h \in Z(\phi)$ , we obtain

$$1/2 \|u\|^2 \leq 1/2 \|v_h\|^2 .$$

**133** Adding

$$1/2 \|\nabla \mu_h\|^2 + (\phi, \mu_h)$$

to both sides, where  $\mu_h \in M$ , we get

$$J(u, \mu_h) \leq J(v_h, \mu_h) \quad \forall \mu_h \in M, v_h \in Z(\phi).$$

That is,

$$\|\nabla(u - \mu_h)\|_0 \leq \|v_h - \nabla \mu_h\|_0 \quad \forall v_h \in Z(\phi).$$

Suppose that  $\mu_h$  is a solution of Dirichlet problem with a conforming finite element method, then an upper bound for the error in the energy norm is given by

$$\|v_h - \nabla \mu_h\|_0$$

where

$$v_h \in Z_h(\phi)$$

is arbitrary. One can choose  $v_h = u_h$ , a solution of the present equilibrium finite element approximation to Dirichlet problem.

## 7.11 Equilibrium Elements for the Plate Problem

We recall that the equations of the plate problem are:

Find  $\sigma_{ij}, w$  such that

$$\sigma_{ij} = \lambda \Delta w \delta_{ij} + 2\mu \frac{\partial^2 w}{\partial x_i \partial x_j} \quad \text{in } \Omega \subset \mathbb{R}^2 \quad (7.57)$$

$$\frac{\partial^2 \sigma_{ij}}{\partial x_i \partial x_j} = f \quad \text{in } \Omega \quad (7.57b)$$



$$w = 0 \text{ on } \partial\Omega \quad (7.57c)$$

134 and

$$\begin{cases} \frac{\partial w}{\partial n} = 0 & \text{on } \partial\Omega \text{ clamped case} \\ \sigma_{ij} n_i n_j = 0 & \text{on } \partial\Omega \text{ simply supported plate problem} \end{cases} \quad (7.57d)$$

The summation convention is used in the above equations.

**Exercise 6.** Show that (7.57) is equivalent to a biharmonic problem for  $w$  alone, and write the variational formulation of that biharmonic problem. Notice that in the clamped case several bilinear forms may be chosen unlike in the simply supported case.

It is easy to check that (7.57) is equivalent to

$$\alpha \sigma_{ij} + \beta (\sigma_{kk}) \delta_{ij} = \frac{\partial^2 w}{\partial x_i \partial x_j}, \quad (7.58)$$

where

$$\alpha = 1/2\mu \quad \text{and} \quad \beta = \frac{-\lambda}{2\mu(2\lambda + 2\mu)}.$$

Let

$$a(\sigma, \tau) = \int_{\Omega} (\alpha \sigma_{ij} \tau_{ij} + \beta (\sigma_{kk})(\tau_{\ell\ell})) dx$$

and

$$Dw = \left( \frac{\partial^2 w}{\partial x_i \partial x_j} \right).$$

Then problem (7.57) is equivalent to:

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Find  $\sigma \in (L^2(\Omega))_s^4$ ,  $w \in H_0^2(\Omega)$  such that

$$a(\sigma, \tau) - \int_{\Omega} \tau Dw dx = 0 \quad \forall \tau \in (L^2(\Omega))_s^4,$$

$$\int_{\Omega} \sigma Dv = \int_{\Omega} fv dx. \quad \forall \begin{cases} v \in H_0^2(\Omega). \\ v \in H^2(\Omega) \cap H_0^1(\Omega). \end{cases}$$

$(L^2(\Omega))_s^4$  denotes the set of symmetric  $2 \times 2$  tensors which are in  $L^2(\Omega)$ .

However, as noticed in Section 7.10, to approximate this problem in the usual way does not represent any progress on the usual conforming approximations since one has to approximate  $H_0^2(\Omega)$ .

We have, by Green's formula,

$$\begin{aligned} \int_K \tau Dv &= \int_K \tau_{ij} \frac{\partial^2 v}{\partial x_i \partial x_j} dx \\ &= \int_{\partial K} \tau_{ij} n_j \frac{\partial v}{\partial x_i} - \int_K \frac{\partial \tau_{ij}}{\partial x_j} \frac{\partial v}{\partial x_i} dx \end{aligned}$$

Since

$$\frac{\partial v}{\partial x_i} = \frac{\partial v}{\partial n} n_i + \frac{\partial v}{\partial s} s_i$$

where  $(s_i)$  are the components of the unit tangent vector, we obtain

$$\int_{\Omega} \tau Dv = \sum_K \int_{\partial K} M_n(\tau) \frac{\partial v}{\partial n} + \sum_K \left( \int_{\partial K} M_{ns}(\tau) \frac{\partial v}{\partial s} - \int_K \frac{\partial \tau_{ij}}{\partial x_j} \frac{\partial v}{\partial x_i} dx \right)$$

Here

$$M_n(\tau) = \tau_{ij} n_i n_j, \quad M_{ns}(\tau) = \tau_{ij} n_j s_i.$$

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$$b(\tau, v) = \sum_K \left( \int_K \frac{\partial \tau_{ij}}{\partial x_j} \frac{\partial v}{\partial x_i} dx - \int_{\partial K} M_{ns}(\tau) \frac{\partial v}{\partial s} ds \right),$$

$$V = \{ \tau : \tau|_K \in (H^1(K))_s^4, K \in T_h, M_n(\tau) \}$$

is continuous across inter element boundaries}.

Then

$$\int_{\Omega} \tau Dv = b(v, \tau) \quad \forall v \in H_0^2(\Omega), \tau \in V.$$

In fact,  $b(\cdot, \cdot)$  is continuous over  $V \times M$  where  $M = W_0^{1,p}(\Omega)$  ( $p > 2$ ) so that (7.57) (clamped case) is equivalent to:

Find  $\{\sigma, w\} \in V \times M$  such that

$$\begin{aligned} a(\sigma, \tau) + b(\tau, w) &= 0 \quad \forall \tau \in V, \\ b(\sigma, v) &= -(f, v) \quad \forall v \in M. \end{aligned} \quad (7.59)$$

We take

$$H = (L^2(\Omega))_s^4.$$

The Brezzi condition holds only on  $H_0^1(\Omega)$ , since if  $v$  is smooth and  $\tau_{ij} = v\delta_{ij}$ , then

$$\begin{aligned} M_{ns}(\tau) &= v (\delta_{ij} n_j s_i) = 0, \\ M_n(\tau) &= v, \end{aligned}$$

and

$$b(v, \tau) = \int_{\Omega} |\nabla v|^2 dx \geq \alpha \|v\|_1^2 \geq c \|v\|_1 \|\tau\|. \quad (7.59)$$

For the proof of existence of solutions of (7.59) and modified error estimates see BREZZI-RAVIART [7].

We choose

$$V_h = \{\tau : \tau \in (\mathbb{P}_0(K))_s^4, K \in T_h, M_n(\tau) \text{ is continuous}\}.$$

Since

$$v \in W_0^{1,p}(\Omega) \subset C^0(\bar{\Omega})$$

we find that after integration by parts on each of  $\partial K$ ,  $b(\tau, v)$  involves only the values of  $v$  at the vertices of  $T_h$ :

$$b(\tau, v) = \sum_K \int_{\partial K} M_{ns}(\tau) \frac{\partial v}{\partial s} = \sum_N R(\tau, N) v(N) \quad \forall \tau \in V_h. \quad (7.59b)$$

Notice that only the value of  $v$  at the vertices has to be taken; therefore, we choose

$$M_h = \{v_h \in C^0(\bar{\Omega}) : v_h|_K \in \mathbb{P}_1(K), K \in T_h, v_h = 0 \text{ on } \partial\Omega\}$$

so that, if  $\tau \in V_h$ , then

$$b(\tau, v_h) = 0 \quad \forall v_h \in M_h \Rightarrow b(\tau, v) = 0 \quad \forall v \in M.$$

Therefore,  $Z_h(0) \subset Z(0)$ . Hence

$$\|\sigma - \sigma_h\|_0 \leq \inf_{\tau_h \in Z_h(-f)} \|\sigma - \tau_h\|.$$

Here  $\{\sigma_h, w_h\} \in V_h \times M_h$  is the solution of the approximate problem.

$$\begin{cases} a(\sigma_h, \tau_h) + \sum R(\tau_h, N) w_h(N) = 0 \quad \forall \tau_h \in V_h, \\ \sum R(\sigma_h, N) v_h(N) = -(f, v_h) \quad \forall v_h \in M_h \end{cases} \quad (7.60)$$

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### Interpolation Operator.

The interpolation operator

$$\pi_h : (H^1(\Omega))^4 \rightarrow V_h$$

is defined by

$$\int_{\gamma} M_n(\pi_h, \sigma) ds = \int_{\gamma} M_n(\sigma) ds,$$

for each edge  $\gamma$  of the triangulation.

We have the estimate

**THEOREM 12.** *There exists a constant  $c$  independent of  $h$  such that*

$$\|\pi_h v\|_V \leq c \|v\|_V \quad (7.61)$$

and

$$\|\pi_h v - v\|_{0,\Omega} \leq ch \|v\|_{1,\Omega} \quad \forall v \in (H^1(\Omega))^4. \quad (7.62)$$

*The proof of this is found in C.JOHNSON [24].*

**Properties of  $T_h$ .**

We have

$$\begin{aligned} b(\pi_h \sigma, v_h) &= \sum_K \int_{\partial K} M_n(\pi_h \sigma) \frac{\partial v_h}{\partial n} - \int_K (\pi_h \sigma) D v_h \\ &= \sum_K \int_{\partial K} M_n(\sigma) \frac{\partial v_h}{\partial n}, \quad \forall v_h \in M_h \\ &= b(\sigma, v_h). \end{aligned}$$

Hence

$$b(\pi_h \sigma - \sigma, v_h) = 0 \quad \forall v_h \in M_h \quad (7.63)$$

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Therefore  $\pi_h$  maps  $Z(\phi)$  into  $Z_h(\phi)$ .

Equations (7.61) and (7.63) imply that the discrete Brezzi condition is satisfied.

We have the error estimate

$$\|w - w_h\|_1 \leq c \left( \|\sigma - \sigma_h\|_0 + \inf_{v_h \in M_h} \|w - v_h\|_{1,p} \right)$$

(See BREZZI-RAVIART [7]).

The above method is called Hermann-Johnson method.

**Morley Nonconforming Method.**

Let

$$W_h = \{v_h : v_h|_K \in \mathbb{P}_2(K), K \in T_h,$$

$v_h$  continuous at the vertices,

$\frac{\partial v_h}{\partial n}$  continuous at the mid side point,

$v_h = 0$  at the boundary vertices,

$\frac{\partial v_h}{\partial n} = 0$  at the mid point of boundary edges}

The space  $W_h$  makes use of the Morley finite element which has 6 degrees of freedom, namely, values at the three vertices and the values of the normal derivatives at the three mid side points.

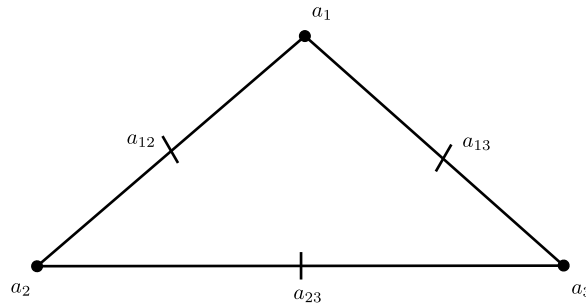


Figure 7.2:

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We consider, for simplicity, the case  $\lambda = 0$  and  $\mu = 1/2$  so that

$$a(\sigma, \tau) = \int_{\Omega} \tau_{ij} \sigma_{ij} dx.$$

Let

$$L(v) = \sum_N f_N v(N);$$

that is  $L$  is a linear combination of Dirac masses (concentrated loads).  
Then

**THEOREM 13.** *The problem:*

*Find  $u_h \in W_h$  such that*

$$\sum_K \int_K D u_h D v_h dx = L(v_h) \quad \forall v_h \in W_h, \quad (7.64)$$

*is equivalent to (7.60) when*

$$(f, v) = \sum_N f_N v(N)$$

in the sense that

$$u_h(N) = W_h(N) \quad \text{at the vertices } N,$$

and

$$\sigma_h|_K = Du_h|_K, \quad K \in T_h.$$

*Proof.* Let  $u_h$  be a solution of (7.64).

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Define

$$\sigma_h^*|_K = Du_h|_K \quad \text{and} \quad w_h^* \in M_h \quad \text{by} \quad w_h^*(N) = u_h(N).$$

We will show that  $\{\sigma_h^*(N), w_h^*\}$  is the solution of (7.60). Since  $u_h$  is a solution (7.64), we have

$$\sum_K \int_K \sigma_h^* Dv_h dx = L(v_h) \quad \forall v_h \in W_h \quad (7.65)$$

Using Green's formula, we obtain

$$\sum_K \left( \int_{\partial K} M_n(\sigma_h^*) \frac{\partial v_h}{\partial n} + \int_{\partial K} M_{ns}(\sigma_h^*) \frac{\partial v_h}{\partial s} \right) = L(v_h) \quad \forall v_h \in W_h \quad (7.66)$$

If  $b_i$  is one mid side point, then substituting  $v_h$  satisfying:

$$v_h = 0 \quad \text{at the vertices}$$

$$\frac{\partial v_h}{\partial n} = \begin{cases} 1 & \text{at } b_i \\ 0 & \text{at the other nodes } b_j, j \neq i \end{cases}$$

in the above equation we obtain that  $M_n(\sigma_h^*)$  is continuous at  $b_i$  (by using 7.59b). This proves that  $\sigma_h^* \in V_h$ .

Since  $\sigma_h^* \in V_h$ , equation (7.66) gives

$$\sum_K \int_{\partial K} M_{ns}(\sigma_h^*) \frac{\partial v_h}{\partial s} = \sum_N f_N v_h(N) \quad \forall v_h \in W_h$$

But

$$\sum_K \int_K M_{ns}(\sigma_h^*) \frac{\partial v_h}{\partial s} = - \sum_N R(\sigma_h^*, N) v_h(N).$$

Hence

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$$\sum_K R(\sigma_h^*, N) v_h(N) = - \sum_N f_N v_h(N) \quad \forall v_h \in W_h. \quad (7.67)$$

Let  $v_h \in M_h$ . Consider  $\tilde{v}_h \in V_h$  defined by

$$\tilde{v}_h(N) = v_h(N), \quad \frac{\partial}{\partial n} \tilde{v}_h(b_i) = 0.$$

Then (7.67) gives

$$\sum_K R(\sigma_h^*, N) \tilde{v}_h(N) = - \sum_N f_N \tilde{v}_h(N)$$

Therefore

$$\sum_N R(\sigma_h^*, N) v_h(N) = - \sum_N f_N v_h(N) \quad \forall v_h \in M_h.$$

This is nothing but the second equation in (7.60) with  $\sigma_h$  replaced by  $\sigma_h^*$ . Now

$$\begin{aligned} a(\sigma_h^*, \tau) &= \sum_K \int_K (\sigma_h^*)_{ij} \tau_{ij} \\ &= \sum_K \int_K \frac{\partial^2 u_h}{\partial x_i \partial x_j} \tau_{ij} \\ &= \sum_K \int_K M_n(\tau) \frac{\partial u_h}{\partial n} + \int_K M_{ns}(\tau) \frac{\partial u_h}{\partial s} \quad \forall \tau \in V_h \end{aligned}$$

by Green's formula.

The first term in the right side is zero since  $u_h \in W_h$  and  $\tau \in V_h$ . The second term equals  $-\sum_N R(\tau, N) u_h(N)$ . Hence we obtain

$$a(\sigma_h^*, \tau) + \sum_N R(\tau, N) w_h^*(N) = 0 \quad \forall \tau \in V_h,$$



since  $u_h(N) = w_h^*(N)$ .

143 Thus  $\{\sigma_h^*, w_h^*\}$  is a solution of (7.60). By uniqueness we have  $\sigma_h = \sigma_h^*$  and  $w_h^* = w_h$ .

Thus we have proved that (7.64)  $\Rightarrow$  (7.60). Let  $\{\sigma_h, w_h\}$  be the solution of (7.60). We will show that  $u_h$  defined by

$$u_h(N) = w_h(N) \quad \text{for each vertex } N \quad (7.68)$$

$$Du_h|_K = \sigma_h|_K \quad \text{for each } K \in T_h \quad (7.69)$$

is the solution of (7.64).

It is easy to see that (7.68) and (7.69) define a unique  $u_h$  such that  $u_h|_K \in \mathbb{P}_2(K)$  for each  $K \in T_h$ . We will prove that this  $u_h \in W_h$ .

From the first equation in (7.60), we obtain

$$\sum_K \int_{\partial K} M_n(\tau) \frac{\partial u_h}{\partial n} = 0 \quad \forall \tau \in V_h.$$

This implies  $\partial u_h / \partial n$  is continuous at mid side points and  $\partial u_h / \partial n = 0$  at the boundary mid side points. Hence  $u_h \in W_h$ .

Let  $v_h \in W_h$ . Then there exists  $v_h \in M_h$  such that  $\tilde{v}_h(N) = v_h(N)$ . Hence the second equation in (7.60) gives

$$\sum R(\sigma_h, N) v_h(N) = - \sum f_N v_h(N).$$

This shows

$$\sum_K \int_K Du_h Dv_h = \sum_N f_N v_h.$$

This proves (7.60)  $\Rightarrow$  (7.64).

Thus the Hermann-Johnson method and the Morley nonconforming 144 method are equivalent, in this particular case where the load is a sum of concentrated loads.  $\square$

**Exercise 7.** Let  $K$  be a triangle. Let  $\mathbb{P}_K = \mathbb{P}_2(K)$  and

$$\sum_K = \left\{ \delta_{a_i}, \frac{\partial}{\partial n} \delta_{a_{ij}}, 1 \leq i < j \leq 3 \right\},$$

where  $a_i$ 's denote the vertices of  $K$  and  $a_{ij}$ 's denote the mid points of the sides of  $K$ . Show that  $\Sigma_K$  is  $\mathbb{P}_K$ -unisolvant. The above finite element is called the Morley finite element.

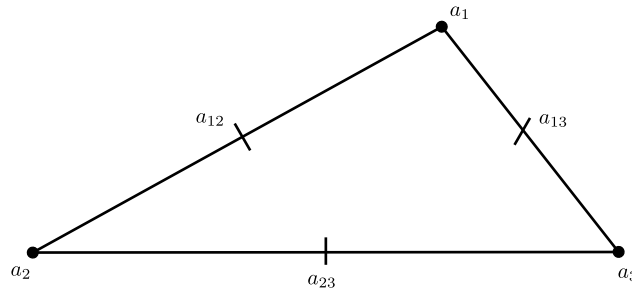


Figure 7.3:

**REMARK 11.** We note that the Morley element has advantage over Herrmann-Johnson method, since in Morley's method we get a positive definite matrix and we have no constraints.

## Chapter 8

# Spectral Approximation for Conforming Finite Element Method

### 8.1 The Eigen Value Problem

Let  $V$  and  $H$  be Hilbert spaces such that  $V \hookrightarrow H$ . We also assume that this imbedding is compact. Let  $\|\cdot\|_1$  denote the norm in  $V$ . The norm in  $H$  is denoted by  $\|\cdot\|$  or  $\|\cdot\|_0$  and the scalar product in  $H$  is  $(\cdot, \cdot)$ . We identify  $H$  with its dual  $H'$ .

Let  $a(\cdot, \cdot) : V \times V \rightarrow \mathbb{R}$  be a continuous, symmetric bilinear form which is  $V$ -coercive with  $\alpha$  as the coercive constant.

We shall consider the eigen value problem:

Find  $u \in V, \mu \in \mathbb{R}$  such that

$$a(u, v) = \mu(u, v) \quad \forall v \in V \quad (8.1)$$

In the following, for an operator  $T : H \rightarrow H$ , we write

$$\|T\| = \sup_{f \in H, f \neq 0} \frac{\|Tf\|}{\|f\|}.$$

## 8.2 The Operator $\mathbb{T}$

The operator  $\mathbb{T} : H \rightarrow V$  is defined as follows. If  $f \in H$  then  $\mathbb{T}f$  is defined to be the unique solution of the variational equation

$$a(\mathbb{T}f, v) = (f, v) \quad \forall v \in V.$$

By Lax-Milgram Lemma  $\mathbb{T}f$  is well defined for all  $f \in H$ . As the  
 146 imbedding  $V \hookrightarrow H$  is compact we obtain that  $\mathbb{T}$ , considered as an operator from  $H$  into  $H$ , is compact. The symmetry of  $a(\cdot, \cdot)$  implies that  $\mathbb{T}$  is symmetric. It is easy to see that (8.1) is equivalent to:

Find  $u \in V$  and  $\lambda \in \mathbb{R}$  such that

$$\mathbb{T}u = \lambda u \tag{8.2}$$

The  $\mu$  and  $\lambda$  in (8.1) and (8.2) have the relation

$$\lambda\mu = 1.$$

From the Spectral Theorem for compact self-adjoint operators we have:

$\text{Sp}(\mathbb{T})$  is a countable set with no accumulation point other than zero. Every point in  $\text{Sp}(\mathbb{T})$  other than zero is an eigenvalue of  $\mathbb{T}$  with finite multiplicity.

## 8.3 Example

The model problem for (8.1) is

$$V = H_0^1(\Omega), \quad H = L^2(\Omega)$$

where  $\Omega$  is a smooth bounded open subset of  $\mathbb{R}^n$ ,

$$a(u, v) = \int_{\Omega} \nabla u \cdot \nabla v$$

The compactness of the imbedding  $H^1(\Omega) \hookrightarrow L^2(\Omega)$  is well known. Problem (8.1) corresponds to:

Find  $u \in H_0^1(\Omega)$ ,  $\mu \in \mathbb{R}$  such that

$$\begin{cases} -\Delta u = u & \text{in } \Omega, \\ u = 0 & \text{on } \partial\Omega \end{cases} \quad (8.3)$$

We note that  $\mathbb{T}$  is the inverse of  $-\Delta$ .

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## 8.4 Approximate Problem

Let  $V_h \hookrightarrow V$  be a finite element subspace of  $V$ . We consider the approximate eigen value problem:

Find  $u_h \in V_h$ ,  $\mu_h \in \mathbb{R}$  such that

$$a(u_h, v_h) = \mu_h(u_h, v_h) \quad \forall v_h \in V_h \quad (8.4)$$

Here again, we introduce an operator  $\mathbb{T}_h : H \rightarrow H$  where  $\mathbb{T}_h f$  is the unique solution of

$$a(\mathbb{T}_h f, v_h) = (f, v_h) \quad \forall v_h \in V_h. \quad (8.5)$$

As in the continuous case, we have  $\|\mathbb{T}_h\| \leq c/\alpha$  which shows that  $\mathbb{T}_h$  is uniformly bounded. Again (8.4) is equivalent to:

Find  $u_h \in V_h$  and  $\lambda_h = 1/\mu_h$  such that

$$\mathbb{T}_h u_h = \lambda_h u_h \quad (8.6)$$

It is obvious that  $\mathbb{T}_h$  is a self-adjoint, compact operator.

We assume that

$$\|\mathbb{T} - \mathbb{T}_h\| \leq \varepsilon(h) \quad (8.7)$$

and

$$\|(\mathbb{T} - \mathbb{T}_h) f\| \leq e(h) \quad (8.8)$$

for all smooth  $f$  and  $\mathbb{T}f$ . Further, we assume

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$$0 \leq e(h) \leq \varepsilon(h) \quad \text{and} \quad \varepsilon(h) \rightarrow 0 \quad (8.9)$$

**EXAMPLE.** Let

$$V_h = \{v_h \in H^1_0(\Omega) : v_h|_K \in \mathbb{P}_k(K), K \in T_h\}$$

where  $T_h$  is a regular family of triangulations of  $\Omega$ . We have (cf. Chapter 5)

$$\|\tau f - \tau_h f\|_{0,\Omega} \leq ch^{s+1} \|f\|_{s-1,\Omega}, \quad 1 \leq s \leq k, \quad (8.10)$$

provided that  $\Omega$  is a convex polygon and that

$$\|\tau f\|_{s+1} \leq c \|f\|_{s-1,\Omega}. \quad (8.11)$$

From GRISVARD [22] this is atleast true for  $s = 1$ , which shows that  $\varepsilon(h) = ch^2$  and  $e(h) = 0(h^{k+1})$ .

## 8.5 Convergence and Error Estimate for the Eigen Space.

Assumption (8.7) (8.9) show that  $\tau_h \rightarrow \tau$  in norm.

From KATO [26] (Chapter 5. Section 4.3) we know that the spectrum of  $\tau_h$  converges to the spectrum of  $\tau$  in the following sense: For all non-zero  $\lambda \in \text{Sp}(\tau)$  with multiplicity  $m$  and for each  $h$  such that  $\varepsilon(h) < d/2$ , where

$$d = \min_{\lambda' \in \text{Sp}(\tau)} |\lambda - \lambda'|,$$

there exist exactly  $m$  eigen values  $\lambda_{ih} \in \text{Sp}(\tau_h)$  (counted according to multiplicity) such that

$$|\lambda - \lambda_{ih}| \leq \varepsilon(h).$$

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Let  $\Gamma = \{z \in \mathbb{C} : |z - \lambda| = d/2\}$ . We know that

$$P = -\frac{1}{2\pi i} \int_{\Gamma} R_z(\tau) dz, \quad (8.12)$$

$$P_h = -\frac{1}{2\pi i} \int_{\Gamma} R_z(\tau_h) dz, \quad \varepsilon(h) < d/2, \quad (8.13)$$

where  $R_z(\mathbb{T}) = (\mathbb{T} - z)^{-1}$ , are the spectral projections on to the eigen-spaces  $E$  and  $E_h$  associated with  $\lambda$  and  $\lambda_{ih}$ 's. The dimension of each of the spaces  $E$  and  $E_h$  is  $m$  (See KATO [26], Chapter 4, Section 4.3).

**LEMMA 1.** For  $u \in E$ , we have

$$\|u - P_h u\| \leq c_2 \|(\mathbb{T} - \mathbb{T}_h) u\|. \quad (8.14)$$

*Proof.* We consider

$$\begin{aligned} R_z(\mathbb{T}) - R_z(\mathbb{T}_h) &= R_z(\mathbb{T}_h)(\mathbb{T}_h - z)R_z(\mathbb{T}) - R_z(\mathbb{T}_h)(\mathbb{T} - z)R_z(\mathbb{T}) \\ &= R_z(\mathbb{T}_h)(\mathbb{T}_h - \mathbb{T})R_z(\mathbb{T}). \end{aligned}$$

Hence

$$P - P_h = -\frac{1}{2\pi i} \int_{\Gamma} R_z(\mathbb{T}_h)(\mathbb{T}_h - \mathbb{T})R_z(\mathbb{T}) dz.$$

Let  $u \in E$ . Then we have

$$Pu = u, \mathbb{T}u = \lambda u \quad \text{and} \quad R_z(\mathbb{T})u = \frac{1}{\lambda - z}u.$$

Therefore

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$$u - P_h u = -\frac{1}{2\pi i} \int_{\Gamma} \frac{R_z(\mathbb{T}_h)}{\lambda - z} dz (\mathbb{T} - \mathbb{T}_h) u.$$

We show that the integral on the right is bounded. Indeed, for  $z \in \Gamma$  and  $\varepsilon(h) < d/2$  we have

$$\begin{aligned} \mathbb{T}_h - z &= \mathbb{T}_h - \mathbb{T} + \mathbb{T} - z \\ &= ((\mathbb{T}_h - \mathbb{T})R_z(\mathbb{T}) + I)(\mathbb{T} - z), \end{aligned}$$

which implies

$$R_z(\mathbb{T}_h) = R_z(\mathbb{T})(I + A_h)^{-1},$$

where

$$A_h = (\mathbb{T}_h - \mathbb{T})R_z(\mathbb{T}).$$

If  $P(\tau)$  denotes the resolvent set of  $\tau$  then, as  $R_z(\tau)$  is continuous in  $z \in P(\tau)$  and  $\Gamma$  is a compact subset of  $P(\tau)$ , we obtain

$$\|R_z(\tau)\| \leq c_1 \quad \text{for all } z \in \Gamma,$$

and  $\|A_h\| \leq c_1 \varepsilon(h)$ , where  $c_1$  is a constant. This implies

$$\|(I + A_h)^{-1}\| \leq 2 \quad \text{for } \varepsilon(h) \leq \frac{1}{2c_1}.$$

Thus

$$\|u - P_h u\| \leq c_2 \|(\tau - \tau_h) u\|.$$

□

**LEMMA 2.** *If the eigen vectors in  $E$  are smooth enough, we have*

$$\delta(E, E_h) \leq c \varepsilon(h), \quad (8.15)$$

151 *where*

$$\delta(E, E_h) = \text{Sup} \{d(u, E_h) : u \in E, \|u\| = 1\}.$$

**Remark 1.** In this chapter we follow closely OSBORN [34]. We will use the result: “In a Hilbert space  $H$ ,  $\delta(E, E_h) = \delta(E_h, E)$ ”. Osborn, however, considers the more general case of a non-self-adjoint operator in a Banach space, which involves more complicated arguments.

## 8.6 Error Estimates for the Eigen Values.

Let  $Q_h = P_h|_E$ , the restriction of  $P_h$  to  $E$ . Then  $Q_h$  maps  $E$  into  $E_h$ . We prove that  $Q_h$  is invertible for small  $h$ .

Indeed, for  $h$  small enough we have  $\dim E = \dim E_h$ .

Let  $f \in E$  be such that  $Q_h f = 0$ . Then

$$\|f\| = \|f - Q_h f\| = \|f - P_h f\| \leq c_2 \varepsilon(h) \|f\|,$$

where we have used Lemma 1. Therefore, for  $c_2 \varepsilon(h) < 1$  we have  $\|f\| = 0$ . Hence  $Q_h$  is invertible for  $\varepsilon(h) < \min(d/2, 1/c_2)$ .



Let us evaluate  $\|Q_h^{-1}\|$ . If  $f \in E$  with  $\|f\| = 1$ , then

$$1 - \|Q_h f\| \leq \|f - P_h f\| \leq c_2 \varepsilon(h).$$

Therefore

$$\|Q_h f\| \geq 1/2, \quad \text{if } \varepsilon(h) \leq 1/2c_2$$

and

$$\|Q_h^{-1}\| \leq 2, \quad \text{for } \varepsilon(h) \leq 1/2c_2.$$

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Let  $\hat{\tau}_h : E \rightarrow E$  be defined by

$$\hat{\tau}_h = Q_h^{-1} \tau_h Q_h.$$

The eigenvalues of  $\hat{\tau}_h$  are again  $\lambda_{jh}, i = 1, 2, \dots, m$  (but the eigenvectors of  $\hat{\tau}_h$  are different from those of  $\tau_h$ ).

Let  $W_{jh} \in E, \|W_{jh}\| = 1$  be an eigen vector of  $\hat{\tau}_h$  associated with the eigen value  $\lambda_{jh}$ . Therefore,

$$\begin{aligned} \lambda - \lambda_{jh} &= ((\lambda - \lambda_{jh}) w_{jh}, w_{jh}) \\ &= ((\tau - \tau_h) w_{jh}, w_{jh}) \\ &\leq \sup_{\phi \in E, \|\phi\|=1} \{((\tau - \tau_h) \phi, \phi)\}. \end{aligned}$$

Now

$$\begin{aligned} \tau - \hat{\tau}_h &= \tau - Q_h^{-1} \tau_h Q_h = \tau - Q_h^{-1} \tau_h P_h \\ &= \tau - Q_h^{-1} P_h \tau_h. \end{aligned}$$

Hence

$$\tau - \hat{\tau}_h = Q_h^{-1} P_h (\tau - \tau_h), \quad (8.16)$$

since  $P_h$  commutes with  $\tau_h$  and  $Q_h^{-1} P_h u = u$  for  $u \in E$ . Hence

$$\|(\tau - \hat{\tau}_h)\phi\| \leq 2 \|(\tau - \tau_h)\phi\| \quad \text{for all } \phi \in E$$

since

$$\|Q_h^{-1}\| \leq 2, \|P_h\| \leq 1 \quad \text{for } \varepsilon(h) \leq 1/2c_1.$$

Therefore

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$$|\lambda - \lambda_{jh}| \leq 2e(h) \quad \text{for} \quad \varepsilon(h) \leq \min(1/2c_1, d/2) \quad (8.17)$$

**Application.** In the example given in Section 8.3, where  $\mathbb{T}$  is the inverse of the negative Laplace operator, we get

$$|\lambda - \lambda_{ih}| \leq ch^{k+1}, \quad 1 \leq i \leq m, \quad (8.18)$$

$$d(E, E_h) \leq ch^{k+1}, \quad (8.19)$$

since  $e(h) \leq ch^{k+1}$  provided that the eigen functions in  $E$  are in  $H^{k+1}$  (This may happen even though  $\Omega$  is a polygon: If  $\Omega = ]0, 1[$  the eigen functions are known explicitly and they are  $C^\infty$ . In fact, the eigen functions are

$$u_{nm} = \text{Sin } n\pi x. \text{Sin } m\pi x.$$

However, the error estimate for  $|\lambda - \lambda_{jh}|$  can be improved as will be shown in the following section. One indeed has

$$|\lambda - \lambda_{jh}| \leq ch^{2k}.$$

## 8.7 Improvement of the Error Estimate for the Eigen Values.

We denote by  $S_h$  the projection on  $E_h$  along  $E^\perp$ . We notice that

$$R_h = Q_h^{-1}P_h$$

is the projection on  $E$  along  $E_h^\perp$ .

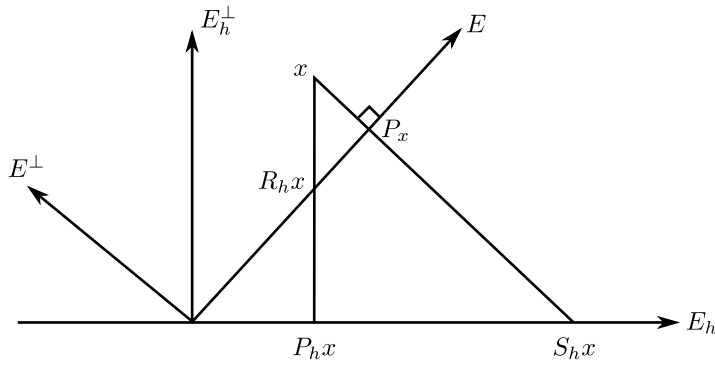


Figure 8.1:

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$S_h$  and  $R_h$  are related by

**LEMMA 3.**  $S_h = R_h^*$ .

*Proof.* For  $u, v \in H$ , we have

$$\begin{aligned}
 (u, S_h \phi) &= (P_h u, S_h \phi), \quad \text{since } P_h S_h = S_h \text{ and } P_h^* = P_h, \\
 &= (P_h R_h u, S_h \phi), \quad \text{since } R_h \text{ is the projection on } E \text{ along } E_h^\perp; \\
 &= (R_h u, S_h \phi) \\
 &= (R_h u, P S_h \phi), \quad \text{since } P R_h = R_h \text{ and } P = P^*, \\
 &= (R_h u, P \phi), \quad \text{since } P S_h = P, \\
 &= (R_h u, \phi).
 \end{aligned}$$

□

We will now prove a Lemma which will give an upper bound for  $\|P - S_h\|$ .

**LEMMA 4.** *We have*

$$\|P - S_h\| \leq \frac{\delta(E_h, E)}{1 - \delta(E_h, E)} \quad (8.20)$$

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*Proof.* Since  $S_h$  is the projection on  $E_h$  along  $E^\perp$  we have  $P = PS_h$ . If  $x \in E_h$  and  $\|x\| = 1$  then

$$\|x - Px\| = d(x, E) \leq \sup_{x \in E_h, \|x\|=1} d(x, E) = \delta(E_h, E).$$

Therefore,

$$\|y - Py\| \leq \delta(E_h, E) \|y\| \quad \text{for all } y \in E_h \quad (8.21)$$

Now, for any  $x \in H$ ,

$$\begin{aligned} \|S_h x\| &\leq \|S_h x - PS_h x\| + \|PS_h x\| \\ &\leq \delta(E_h, E) \|S_h x\| + \|Px\|. \end{aligned}$$

Thus

$$\|S_h x\| \leq \frac{1}{1 - \delta(E_h, E)} \|x\| \quad (8.22)$$

Hence we obtain

$$\begin{aligned} \|(P - S_h)x\| &= \|PS_h x - S_h x\| \\ &\leq \delta(E_h, E) \|S_h x\|, \quad \text{by (8.21),} \\ &\leq \frac{\delta(E_h, E)}{1 - \delta(E_h, E)} \|x\|, \quad \text{using (8.22)} \end{aligned}$$

Therefore

$$\|P - S_h\| \leq \frac{\delta(E_h, E)}{1 - \delta(E_h, E)}.$$

□

Finally we have from (8.16),

$$\begin{aligned} ((T - \hat{T}_h)\phi, \phi) &= (R_h(T - T_h)\phi, \phi) \\ &= ((T - T_h)\phi, S_h\phi) \\ &= ((T - T_h)\phi, \phi) + ((T - T_h)\phi, S_h\phi - \phi), \end{aligned}$$

156 where  $\phi \in E$  and  $\|\phi\| = 1$ .

For  $\phi \in E$ , using Lemma 4 and Lemma 2, we obtain

$$\|\phi - S_h\| \leq \frac{\delta(E_h, E)}{1 - \delta(E_h, E)} \leq k e(h),$$

for sufficiently small  $h$ , where  $k$  is a constant.

We know that for all  $\phi \in E$  with  $\|\phi\| = 1$

$$\lambda - \lambda_{ih} = ((T - \hat{T}_h)\phi, \phi).$$

Hence

$$|\lambda - \lambda_{ih}| \leq \sup_{\phi \in E, \|\phi\|=1} ((T - T_h)\phi, \phi) + k(e(h))^2.$$

Thus we have proved

**THEOREM 5.** *When  $E$  is a smooth subset of  $H$  and  $h$  is sufficiently small we have*

$$|\lambda - \lambda_{ih}| \leq \sup_{\phi \in E, \|\phi\|=1} ((T - T_h)\phi, \phi) + k(e(h))^2, \quad (8.23)$$

where  $k$  is a constant.

**Application.** In the case of the example in Section 8.3, we give an estimate of 157

$$\alpha_h = \sup_{\phi \in E, \|\phi\|=1} ((T - T_h)\phi, \phi).$$

Let  $w$  be the solution of

$$a(v, w) = (\phi, v) \quad \forall v \in V, \quad (8.24)$$

where  $\phi \in E$  and  $\|\phi\| = 1$ .

We have

$$\begin{aligned} ((T - T_h)\phi, \phi) &= a((T - T_h)\phi, w) \\ &= a((T - T_h)\phi, w - v_h) \quad \text{for all } v_h \in V_h, \end{aligned}$$

since

$$a((T - T_h)\phi, v_h) = 0 \quad \text{for all } v_h \in V_h.$$

If  $w \in H^{k+1}(\Omega)$ , we get

$$\|w - v_h\|_1 \leq ch^k \|w\|_{k+1} \leq ch^k \|\phi\|_{k-1},$$

using regularity theorem and the error estimates in Chapter 5.

Since  $E$  is finite-dimensional there exists a constant  $c$  such that

$$\|\phi\|_{k-1} \leq c \|\phi\| = c.$$

Finally, we have

$$\begin{aligned} \alpha_h &\leq ch^k \|(\mathbb{T} - \mathbb{T}_h)\phi\|_1 \\ &\leq ch^{2k}, \end{aligned}$$

158 provided that  $E \subset H^{k+1}$ .

Thus we have proved

**THEOREM 6.** *For the model problem (See Section 8.3) we have*

$$|\lambda - \lambda_{ih}| \leq ch^{2k},$$

provided  $E \subset H^{k+1}(\Omega)$ , the solution of (8.24) is in  $H^{k+1}(\Omega)$  and  $h$  is small.

**REMARK 2.** *Error estimates for the semi-discrete approximation to parabolic equation of the form*

$$\left(\frac{du}{dt}, v\right) + a(u, v) = (f, v) \quad \text{for all } v \in V,$$

$$u(0) = v_0, \quad v_0 \in V,$$

can be obtained using spectral approximation. The reader is referred to THOMEE [44], [45].

## Chapter 9

# Nonlinear Problems

### Introduction.

We consider here problems of the following type:

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Find  $u \in C$  such that

$$J(u) \leq J(v) \quad \text{for all } v \in V, \quad (9.1)$$

where  $C$  is a closed, convex subset of a Banach space  $V$  and  $J : V \rightarrow \mathbb{R}$  is a convex, lower semi-continuous (*l.s.c.*) function.

We denote by  $(\cdot, \cdot)$  the duality pairing  $V' - V$  and  $\|\cdot\|$  the norm in  $V$ .

We write (9.1) in the form:

Find  $u \in C$  such that

$$J(u) = \inf_{v \in C} J(v). \quad (9.2)$$

The existence and uniqueness of the solutions of (9.2) is given by

**THEOREM 1.** *Assume that  $J$  is coercive on  $C$ , that is*

$$J(v) \rightarrow \infty \quad \text{if} \quad \|v\| \rightarrow \infty, v \in C. \quad (9.3)$$

*Then problem (9.2) has at least one solution provided that  $V$  is reflexive.*

*If  $J$  is strictly convex, then (9.2) has almost one solution.*

The proof of theorem 1 can be found in EMELAND-TEMAN [15].

**Exercise 1.** If  $J$  is Gateaux-differentiable, then  $u$  is a solution of (9.2) 160  
iff

$$(J'(u), v - u) \geq 0 \quad \text{for all } v \in C.$$

If  $C$  is affine linear, (9.1) implies

$$(J'(u), v - u) = 0 \quad \forall v \in C.$$

**Approximation.**

Let  $V_h$  be a finite-dimensional subspace of  $V$  and  $C_h$  a closed, convex subset of  $V_h$ . Then the approximate problem corresponding to (9.2) is:

Find  $u_h \in C_h$  such that

$$J(u_h) = \inf_{v_h \in C_h} J(v_h). \quad (9.4)$$

We assume that  $J$  is strictly convex and coercive. Moreover, we assume that  $C_h$  approximates  $C$ .

$$\left. \begin{array}{l} \text{For all } v \in C \text{ there exists } v_h \in C_h \\ \text{such that } v_h \rightarrow v \text{ (strongly) as } h \rightarrow 0; \end{array} \right\} \quad (9.5)$$

$$\left. \begin{array}{l} \text{If } w_h \rightarrow w \in V \text{ as } h \rightarrow 0 \text{ and } w_h \in C_h \\ \text{then } w \in C. \end{array} \right\} \quad (9.6)$$

Then one can easily prove that the solution  $u_h$  of (9.4) converges weakly to  $u$ , the solution of (9.2).

Note that (9.5) implies that  $C_h$  has to be sufficiently big and (9.6) demands  $C_h$  to be sufficiently small.

161 To get strong convergence of  $u_h$  to  $u$ , one needs some strong monotonicity; there exists  $\alpha, \gamma > 0$  such that

$$(J'(u) - J'(v), u - v) \geq \alpha \|u - v\|^\gamma, \quad \forall u, v \in C \quad (9.7)$$



**Case 1.** We obtain an error estimate when  $C_h = C \cap V_h$ . From Exercise 1, we have

$$(J'(u), v - u) \geq 0 \quad \forall v \in C, \quad (9.8)$$

$$(J'(u_h), v_h - u_h) \geq 0 \quad \forall v_h \in C_h. \quad (9.9)$$

As  $C_h = C \cap V_h$ , choosing  $v = u_h$  in (9.8) and adding to (9.9), we get

$$(J'(u) - J'(u_h), u_h - u) + (J'(u_h), v_h - u) \geq 0.$$

Therefore (assuming  $J$  to be continuously differentiable)

$$\alpha \|u_h - u\|^\gamma \leq (J'(u_h), v_h - u) \leq c \|v_h - u\|,$$

since  $u_h$  is bounded. Finally,

$$\|u_h - u\| \leq c \inf_{v_h \in C_h} \|v_h - u\|^{1/\gamma}, \quad (9.10)$$

provided  $J'$  is weakly continuous.

Note that

$$\inf_{v_h \in C_h} \|v_h - u\|$$

measures how good is the approximation  $C_h$  to  $C$ . Note also that, as  $u_h$  is bounded, it is enough if (9.7) holds on bounded subsets of  $C$ .

**Exercise 2.** Let  $\phi : V \rightarrow \mathbb{R}$  be Lipschitz but not differentiable. Then 162

$$J(u) = \inf_{v \in C} [J(v) + \phi(v)]$$

is equivalent to

$$(J'(u), v - u) + \phi(v) - \phi(u) \geq 0 \quad \forall v \in C.$$

Derive an error estimate similar to (9.10).

**Case 2.** When  $C$  is of the form

$$C = \{v \in V : b(v, \mu) = (\phi, \mu) \quad \forall \mu \in M\}, \quad (9.11)$$

where  $M$  is a Hilbert space,  $b(\cdot, \cdot)$  is a continuous, bilinear form on  $V \times M$  and  $\phi \in M$  satisfying Brezzi's condition (See Chapter 7), Problem (9.1) is equivalent to:

Find  $\{u, \lambda\} \in V \times M$  such that

$$\langle J'(u), v \rangle + b(v, \lambda) = 0 \quad \forall v \in V, \quad (9.12)$$

$$b(u, \mu) = (\phi, \mu) \quad \forall \mu \in M. \quad (9.13)$$

We notice that (9.11) is affine linear and from Exercise 1 we obtain that (9.1) is equivalent to

$$(J'(u), v - u) = 0 \quad \forall v \in C. \quad (9.14)$$

Let  $B : V \rightarrow M$  be defined by

$$(Bv, \mu) = b(v, \mu) \quad \forall v \in V, \mu \in M.$$

**163** Clearly  $C = v + \text{Ker } B$ , where  $v \in C$ . This together with (9.14) implies

$$J'(u) \in (\text{Ker } B)^\perp = \text{Im } B^*,$$

which is closed from Brezzi's condition. Hence there exists  $\lambda \in M$  such that

$$J'(u) = -B^* \lambda.$$

Thus

$$(J'(u), v) + b(v, \lambda) = 0 \quad \forall v \in V.$$

$u \in C$  implies  $b(u, \mu) = (\phi, \mu) \forall \mu \in M$ .

So we have proved that (9.1) implies (9.12) and (9.13).

If (9.12) and (9.13) hold, then  $u \in C$  and

$$(J'(u), u) = -b(u, \lambda) = -b(v, \lambda) \quad \forall v \in C$$

Hence  $\langle J'(u), v - u \rangle = 0 \quad \forall v \in C$ , which is equivalent to (9.1).

Thus we proved the equivalence of (9.1) and (9.12) (9.13).

A natural approximation  $C_h$  to  $C$  will be

$$C_h = \{v \in V : b(v, \mu_h) = (\phi, \mu_h) \quad \forall \mu_h \in M_h\},$$

where  $M_h$  approximates  $M$ . In this case

$$C_h \not\subset C.$$

164 **EXAMPLE 1.** Nonlinear Dirichlet Problem.

$$V = W^{1,p}(\Omega), C = V$$

$$J(v) = \frac{1}{p} \int_{\Omega} |\nabla v|^p dx - \int_{\Omega} f v dx,$$

where

$$f \in L^q, 1/p + 1/q = 1$$

For  $1 < p < \infty$ ,  $W^{1,p}(\Omega)$  is reflexive and  $J(v) \rightarrow \infty$  as  $\|v\| \rightarrow \infty$ . One has

$$(J'(u), v) = \int_{\Omega} |\nabla u|^{p-2} \nabla u \cdot \nabla v dx - \int_{\Omega} f v dx,$$

and some strong monotonicity results of the type (9.7) are proved in GLOWINSKI-MARROCCO [19].

**EXAMPLE 2.** The Obstacle Problem.

$$V = H^1_0(\Omega),$$

$$C = \{v \in H^1_0(\Omega) : v \geq 0 \text{ a. e. on } \Omega\},$$

$$J(v) = \frac{1}{2} \int_{\Omega} |\nabla v|^2 dx - \int_{\Omega} f v dx.$$

Existence and uniqueness of the solution of the minimization problem are straightforward.

Let  $V_h$  be the standard Lagrange finite element space of degree 1 and  $C_h = C \cap V_h$ . One has (9.10) with  $\gamma = 2$ ; therefore it seems that one gets

$$\|u - u_h\| = O(\sqrt{h})$$

since the interpolate  $\pi_h u \in C_h$  as long as  $u \in C$ . However, one has

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$$(J'(u_h), v_h - u) = (J'(u), v_h - u) + (J'(u_h) - J'(u), v_h - u)$$

$$\leq (-\Delta u - f, v_h - u) + \frac{M\varepsilon}{2} \|u_h - u\|_1^2 + \frac{M}{2\varepsilon} \|v_h - u\|_1^2.$$

Hence

$$\|u - u_h\|_1^2 \leq c \left( \|v_h - u\|_0^2 + \|v_h - u\|_1^2 \right)$$

Therefore

$$\|u - u_h\|_1 = O(h).$$

**EXAMPLE 3.** Elasto - Plastic Torsion.  $J$  and  $V$  as is Example 2,

$$C = \{v \in H_0^1(\Omega) : |\nabla v| \leq 1 \text{ a. e. on } \Omega\},$$

$V_h$  same as in Example 2 and  $C_h = C \cap V_h$ .

In this case the interpolate  $\pi_h u$  is not in  $C_h$  whereas  $u$  is in  $C$ . One gets  $O(h^{1/2-\varepsilon})$ .

**EXAMPLE 4.** The Flow of a Bingham Fluid in a Cylindrical Pipe. This is a particular case of Exercise 2 with  $J, V$  as above and

$$\phi(v) = \int_{\Omega} |\nabla v| dx.$$

## 9.2 Generalization

Note that  $J' : V \rightarrow V'$  satisfies

$$(J'(u) - J'(v), u - v) \geq 0 \quad \forall u, v \in V.$$

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An operator  $A : V \rightarrow V'$  is said to be *monotone* if

$$(Au - Av, u - v) \geq 0 \quad \forall u, v \in V. \quad (9.15)$$

$A$  is *bounded* if  $A$  maps bounded sets of  $V$  into bounded sets of  $V'$ .

$A$  is *hemi-continuous* if

$$\lim_{\lambda \rightarrow 0} (A(u + \lambda w), v) = (A(u), v) \quad \forall u, v, w \in V. \quad (9.16)$$

$A$  is *coercive* if

$$\frac{(A(v), v)}{\|v\|} \rightarrow \infty \text{ if } \|v\| \rightarrow \infty \text{ for } v \in C. \quad (9.17)$$

We have

**THEOREM 2.** *If  $A$  is a monotone, bounded hemi continuous and coercive operator then the problem:*

*Find  $u \in C$  such that*

$$(A(u), v - u) \geq 0 \quad \forall v \in C \quad (9.18)$$

*has at least one solution.*

For a proof of this Theorem see LIONS [28]. The problem (9.18) has at most one solution if  $A$  is *strongly monotone*, i.e. there exists  $\alpha, \gamma > 0$  such that

$$\alpha \|u - v\|^\gamma \leq (A(u) - A(v), u - v) \quad \forall u, v \in C \quad (9.19)$$

The error analysis can be carried out in the same way.

### 9.3 Contractive Operators.

Let  $T : C \rightarrow C$  be a mapping, where  $C$  is a closed, convex subset of a Hilbert space  $H$ . The scalar product in  $H$  is denoted by  $(\cdot, \cdot)$ .

We call  $T$  *contractive* iff

$$\|Tx - Ty\| \leq \|x - y\|, \quad \forall x, y \in C. \quad (9.20)$$

$T$  is *strictly contractive* iff there exists a  $\theta$  with  $0 < \theta < 1$  such that

$$\|Tx - Ty\| \leq \theta \|x - y\| \quad \forall x, y \in C. \quad (9.21)$$

We say that  $T$  is *firmly contractive* iff (cf. BROWDER-PETRYSHN [8])

$$\|Tx - Ty\|^2 \leq (Tx - Ty, x - y) \quad \forall x, y \in C \quad (9.22)$$

$T$  is *quasi firmly contractive* iff there exists a  $\theta, 0 < \theta < 1$  such that

$$\|Tx - Ty\|^2 \leq \theta(Tx - Ty, x - y) + (1 - \theta) \|x - y\|^2 \quad (9.23)$$

Note that (9.22)  $\Rightarrow$  (9.23)  $\Rightarrow$  (9.20) and (9.21)  $\Rightarrow$  (9.20).

**Geometrical Interpretation of the Above Definitions** Let  $y \in C$  be a fixed point of  $T$ , i.e.  $Ty = y$ . If  $T$  is contractive and  $x \in C$  then  $Tx$  lies in

the closed ball with  $y$  as centre and  $\|y - x\|$  as radius. If  $T$  is strictly contractive, then  $Tx$  lies in the open ball with  $y$  as centre and  $\|y - x\|$  as radius, for all  $x \in C$ . If  $T$  is firmly contractive then from (9.22) we obtain

$$(Tx - y, x - Tx) \geq 0 \quad \forall x \in C.$$

168 This means that the angle between  $y - Tx$  and  $x - Tx$  is obtuse.

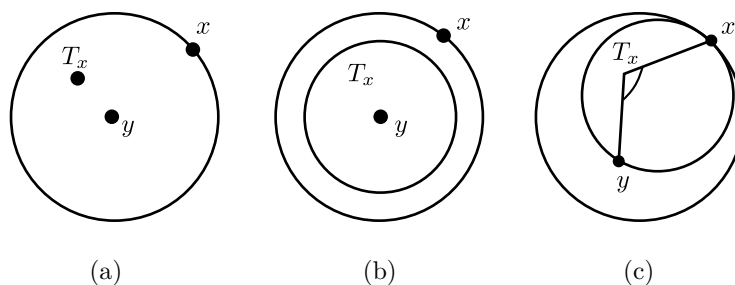


Figure 9.1:

Note that if  $T_1$  and  $T_2$  are contractive then  $T = T_1T_2$  is contractive and  $T$  is strictly contractive if any one of  $T_1$  and  $T_2$  is. However,  $T_1, T_2$  firmly (or quasi-firmly) contractive implies *only*  $T = T_1T_2$  is *contractive*.

**Fixed Points.** We recall that  $x$  is a fixed point of  $T$  iff  $Tx = x$ . Let  $F(T)$  denote the set of all fixed points of  $T$ . If  $T$  is strictly contractive then  $T$  has a unique fixed point and  $F(T)$  is singleton. We have

**THEOREM 3.** *If  $T$  is contractive, then  $F(T)$  is closed and convex.*

*Proof.* Let  $x_n \in F(T), x_n \rightarrow x$ . Then

$$\|x_n - Tx\| \leq \|x_n - x\|$$

169 Taking the limit as  $n \rightarrow \infty$  we get  $\|x - Tx\| = 0$ , i.e.  $x \in F(T)$ . Hence  $F(T)$  is closed.

Let  $x, y \in F(T)$  and  $u = \theta x + (1 - \theta)y$  where  $0 < \theta < 1$ . We have

$$\|x - Tu\| \leq \|x - u\| = (1 - \theta) \|x - y\|, \quad (9.24)$$

$$\|y - Tu\| \leq \|y - u\| = \theta \|x - y\|, \quad (9.25)$$

$$\|x - y\| \leq \|x - Tu\| + \|y - Tu\| \leq \|x - y\|, \quad (9.26)$$

Since  $H$  is strictly convex, we obtain using (9.24) and (9.25),

$$x - Tu = c(y - Tu) \quad (9.27)$$

$$\|x - Tu\|^2 = c(y - Tu, x - Tu), \quad (9.28)$$

$$\|x - y\|^2 = \|x - Tu\|^2 + \|y - Tu\|^2 + 2\|x - Tu\| \|y - Tu\|,$$

by (9.26),

$$\begin{aligned} \|x - y\|^2 &= \|x - Tu + Tu - y\|^2 \\ &= \|x - Tu\|^2 + \|Tu - y\|^2 + 2(x - Tu, Tu - y) \end{aligned}$$

So

$$(x - Tu, Tu - y) = \|x - Tu\| \|y - Tu\| > 0. \quad (9.29)$$

From (9.28) and (9.29) we obtain  $c < 0$ .

Equations (9.26) and (9.27) imply

$$\|y - Tu\| = \frac{1}{1 + |c|} \|x - y\|.$$

This with (9.25) gives  $|c| \geq (1 - \theta)\theta^{-1}$ . Similarly using (9.24), (9.26) and (9.27) we obtain  $|c| \leq (1 - \theta)\theta^{-1}$ . Thus

$$|c| = (1 - \theta)\theta^{-1}.$$

Hence

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$$\theta(x - Tu) = -(1 - \theta)(y - Tu).$$

Therefore

$$Tu = \theta x + (1 - \theta)y = u,$$

that is

$$u \in F(T).$$

□

**REMARK 1.** *Theorem 3 can be proved geometrically.*

Let  $u = \theta x + (1 - \theta)y$ ,  $x, y \in F$ . Since  $x \in F(T)$  and  $T$  is contractive  $Tu$  lies in the closed ball  $C_x$  with  $x$  as centre and  $\|x - u\|$  as radius. Similarly  $Tu$  lies in the closed ball  $C_y$  with  $y$  as centre and  $\|y - u\|$  as radius. But  $C_x \cap C_y = u$ . Hence  $Tu = u$ . Thus  $u \in F(T)$  and  $F(T)$  is convex.

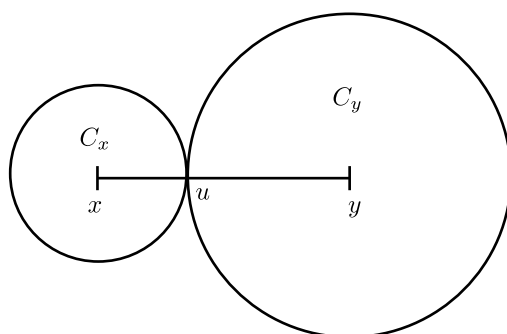


Figure 9.2:

If  $C$  is bounded and  $T$  is contractive, then

$$F(T) \neq \phi.$$

171 In the following we assume  $F(T) \neq \phi$  and study the convergence of the iterative method

$$x^{n+1} = Tx^n.$$

which is known to be strongly convergent to the unique fixed point of  $T$  if  $T$  is strictly contractive. One has

**THEOREM 4.** *If  $T$  is firmly contractive and  $F(T) \neq \phi$  then*

$$x^n \rightarrow \xi \in F(T) \quad \text{as } n \rightarrow \infty$$

*i.e.  $x^n$  converges weakly to a fixed point.*



*Proof.* Let  $y \in F(T)$ . We have

$$\|x^{n+1} - y\|^2 \leq (x^{n+1} - y, x^n - y)$$

But

$$\begin{aligned} \frac{1}{2} \|x^{n+1} - x^n\|^2 &= \frac{1}{2} \|x^{n+1} - y\|^2 + \frac{1}{2} \|x^n - y\|^2 - (x^{n+1} - y, x^n - y) \\ &\leq \frac{1}{2} \|x^n - y\|^2 - \frac{1}{2} \|x^{n+1} - y\|^2. \end{aligned}$$

Therefore

$$\begin{aligned} \|x^{n+1} - y\|^2 + \|x^{n+1} - x^n\|^2 &\leq \|x^n - y\|^2, \\ \|x^{N+1} - y\|^2 + \sum_{n=0}^N \|x^{n+1} - x^n\|^2 &\leq \|x^0 - y\|^2, \end{aligned}$$

which proves that  $\|x^{n+1} - x^n\| \rightarrow 0$  and  $\{x^n\}$  is a bounded sequence. Let  $x^{n'} \rightarrow x$  be a weakly convergent subsequence.

Since

$$(Tx - Ty, Tx - Ty + y - x) \leq 0$$

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choosing  $y = x^{n'-1}$  we obtain

$$(Tx - x^{n'}, Tx - x^{n'} + x^{n'-1} - x) \leq 0$$

As  $n' \rightarrow \infty$ , we get

$$(Tx - x, Tx - x) \leq 0,$$

and hence  $x \in F(T)$ .

As  $\|x^n - y\|^2$  is a decreasing sequence for any  $y \in F(T)$  it converges to some number  $P(y)$ , and we conclude from the following Lemma that the whole sequence  $x^n$  converges.  $\square$

**OPAL'S LEMMA 5.** *Let  $F \subset H$  be a subset of a Hilbert space  $H$  and  $\{x_n\}$  a sequence such that*

$$(i) \quad \|x^n - y\|^2 \rightarrow P(y) \text{ as } n \rightarrow \infty \text{ for any } y \in F$$

(ii) any weakly converging subsequence  $x_{n'} \rightharpoonup z$  is such that  $z$  belongs actually to  $F$ .

Then  $x^n \rightharpoonup \xi \in F$ .

*Proof.* Let  $x_{m'} \rightharpoonup y, x_{n'} \rightharpoonup z$  be two converging subsequences, we have

$$\begin{aligned} \|x_n - y\|^2 &= \|x_n - z + z - y\|^2 \\ &= \|x_n - z\|^2 + 2(x_n - z, z - y) + \|z - y\|^2 \end{aligned}$$

hence taking the limit following  $m$

$$P(y) = P(z) + 2(y - z, z - y) + \|z - y\|^2 = P(z) - \|z - y\|^2$$

173 and taking the limit following  $n'$

$$P(y) = P(z) + 0 + \|z - y\|^2$$

hence  $\|z - y\|^2 = 0 \Rightarrow z = y$ . □

**Exercise 3.** Prove Theorem 4 when  $T$  is quasi firmly contractive.

**THEOREM 6.** Let  $T = QS$  where  $S$  is quasi-firmly contractive and  $Q$  is firmly contractive. Then

$$x^n \rightharpoonup x \in F(T),$$

provided that  $F(T)$  is non-empty.

*Proof.* Let  $y \in F(T)$ . We have

$$\begin{aligned} \|Sx^n - Sy\|^2 &\leq \theta(Sx^n - Sy, x^n - y) + (1 - \theta) \|x^n - y\|^2, \\ (Sx^n - Sy, x^n - y) &= \frac{1}{2} \|Sx^n - Sy\|^2 + \frac{1}{2} \|x^n - y\|^2 \\ &\quad - \frac{1}{2} \|Sx^n - Sy + y - x^n\|^2. \end{aligned}$$

Therefore,

$$\left(1 - \frac{\theta}{2}\right) \|Sx^n - Sy\|^2 + \frac{\theta}{2} \|Sx^n - Sy + y - x^n\|^2 \leq \left(1 - \frac{\theta}{2}\right) \|x^n - y\|^2 \quad (9.30)$$

In the same way,

$$\begin{aligned} \|QSx^n - y\|^2 &\leq (QSx^n - y, Sx^n - Sy) \\ &= \frac{1}{2} \|QSx^n - y\|^2 + \frac{1}{2} \|Sx^n - Sy\|^2 \\ &\quad - \frac{1}{2} \|QSx^n - y + Sy - Sx^n\|^2, \end{aligned}$$

i.e.

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$$\frac{1}{2} \|x^{n+1} - y\|^2 + \frac{1}{2} \|x^{n+1} - y + Sy - Sx^n\|^2 \leq \frac{1}{2} \|Sx^n - Sy\|^2 \quad (9.31)$$

From (9.30) and (9.31) we obtain

$$\|x^{n+1} - y\|^2 + \|x^{n+1} - y + Sy - Sx^n\|^2 + \alpha \|Sx^n - Sy + y - x^n\|^2 \leq \|x^n - y\|^2$$

where

$$\alpha = \frac{\theta}{2 - \theta}$$

This implies

$$\begin{aligned} \|x^{N+1} - y\|^2 + \sum_{n=0}^N (\|x^{n+1} - y + Sy - Sx^n\|^2 + \\ + \alpha \|Sx^n - Sy + y - x^n\|^2) \leq \|x^0 - y\|^2. \end{aligned}$$

Therefore,

$$\begin{aligned} x^{n+1} - Sx^n &\rightarrow y - Sy, \\ Sx^n - x^n &\rightarrow Sy - y, \\ x^{n+1} - x^n &\rightarrow 0. \end{aligned}$$

Let  $x^{n'} \rightarrow x$ ,  $T$  being contractive, we have

$$\|Tx^{n'} - Tx\|^2 \leq \|x^{n'} - x\|^2$$

that is

$$(x^{n'+1} - x^{n'} + x - Tx, x^{n'+1} + x^{n'} - Tx - x) \leq 0$$

and to the limit

$$\begin{aligned} (x - Tx, x + x - Tx - x) &\leq 0 \\ x &= Tx. \end{aligned}$$

- 175 Once again we apply Opial's Lemma to get the convergence of the whole sequence  $x^n$  to a fixed point of  $T$ .  $\square$

**Exercise 4.** Let  $C \subset V$  be a closed convex subset of a Hilbert space  $V$ , then show that the projection map  $P_C : V \rightarrow C$  is firmly contractive.

## 9.4 Application to Unconstrained Problem

We shall apply the previous results to the solution of

$$A(u) = 0.$$

where  $A$  is a monotone operator from  $D(A)$  into  $H$ ; i.e.

$$(Au - Av, u - v) \geq 0 \quad \forall u, v \in D(A).$$

$A$  is said to be *maximal monotone* if  $E \subset H \times H$ ,  $\text{Graph } A \subset E$ ,

$$(x_1 - x_2, y_1 - y_2) \geq 0 \quad \forall \{x_i, y_i\} \in E, i = 1, 2$$

implies

$$\text{Graph } A = E.$$

It is proved in BRÉZIS [4] that

**THEOREM 7.** *A maximal monotone iff*

$$R(I + \lambda A) = H \quad \text{for } \lambda \geq 0.$$

- 176 **EXAMPLE 5.** Let  $A : V \rightarrow V'$  satisfy (9.15) - (9.17) with

$$V \overset{\text{dense}}{\hookrightarrow} H \hookrightarrow V'$$

Then the restriction of  $A$  to

$$D(A) = \{v \in V : Av \in H\}$$

is a maximal monotone operator.

**Exercise 5.** Use Theorem 6 to prove that the operator defined in Example 5 is monotone.

We have

**LEMMA 8.** *If  $A$  is maximal monotone then*

$$T = (I + \lambda A)^{-1}$$

*is firmly contractive.*

*Proof.* Let

$$(I + \lambda A)x = (I + \lambda A)y.$$

Then

$$\lambda(A(x) - A(y)) = -(x - y)$$

Therefore

$$-\|x - y\|^2 = \lambda(A(x) - A(y), x - y) \geq 0.$$

Hence  $x = y$ . This proves  $(I + \lambda A)$  is one-one.

From Theorem 6, we obtain  $R(I + \lambda A) = H$ . Hence  $(I + \lambda A)^{-1}$  is well defined on  $H$ . 177

Let

$$u_i = Tx_i, \quad x_i \in H, \quad i = 1, 2.$$

Then

$$u_i + \lambda Au_i = x_i.$$

We have to prove that

$$\|Tx_1 - Tx_2\|^2 \leq (Tx_1 - Tx_2, x_1 - x_2),$$

i.e.

$$(Tx_1 - Tx_2, Tx_1 - Tx_2 + x_2 - x_1) \leq 0,$$

i.e.

$$(u_1 - u_2, (u_1 - u_2) - (u_1 - u_2) - \lambda(Au_1 - Au_2)) \leq 0,$$

i.e.

$$-\lambda(u_1 - u_2, Au_1 - Au_2) \leq 0,$$

which is true since  $A$  is monotone. □

**COROLLARY 1.** *The algorithm*

$$x^{n+1} = (I + \lambda A)^{-1} x^n \quad (9.34)$$

*converges weakly to a solution of*

$$A(u) = 0 \quad (9.35)$$

*Note that algorithm (9.34) can be written as*

$$\frac{x^{n+1} - x^n}{\lambda} + A(x^{n+1}) = 0 \quad (9.36)$$

**178** *and corresponds to an implicit scheme for*

$$\frac{\partial u}{\partial t} + A(u) = 0. \quad (9.37)$$

*Proof.* Since  $T = (I + \lambda A)^{-1}$  is firmly contractive, algorithm (9.34) converges weakly to a fixed point of  $T$  which is a solution of (9.35).  $\square$

**REMARK 2.** *Algorithm (9.34) is called a proximal point algorithm. Note that computing  $x^{n+1}$  at each step might be as difficult as the original problem except in some special cases.*

**REMARK 3.** *If  $A : V \rightarrow V'$  where  $V$  is a Hilbert space, then it is better to choose  $H = V$ . Let  $J : V' \rightarrow V$  be the Riesz isometry. Then one has to replace  $A$  by  $JA$ . Then algorithm (9.34) is an implicit scheme for*

$$\frac{\partial u}{\partial t} + JA(u) = 0.$$

## 9.5 Application to Problems with Constraint.

We want to solve the problem

$$(A(u), v - u) \geq 0 \quad \forall v \in C. \quad (9.38)$$

If  $u$  is a solution of (9.38) then for any  $\lambda > 0$  we have

$$(u - \lambda A(u) - u, v - u) \leq 0 \quad \forall v \in C$$

179 which implies  $u = P_C S u$ , where

$$S u = u - \lambda A(u).$$

Conversely if  $u$  is a fixed point of  $P_C S$ , then  $u$  is a solution of (9.38). We like to solve (9.38) via the algorithm

$$x^{n+1} = P_C S x^n = P_C(x^n - \lambda A(x^n)). \quad (9.39)$$

Note that if  $J$  is a convex, *l.s.c.*, Gateaux differentiable function and  $A = J'$  then (9.38) is the gradient algorithm with projection for solving

$$\inf_{v \in C} J(v).$$

We will now give some conditions on  $A$  and  $\lambda$  which will ensure the convergence of the algorithm (9.39) to a solution of (9.38).

**THEOREM 9.** *If  $A$  is strongly monotone, i.e.*

$$(A(u) - A(v), u - v) \geq \alpha \|u - v\|^2 \quad \forall u, v \in C \quad (9.40)$$

and Lipschitzian,

$$\|A(u) - A(v)\| \leq c \|u - v\| \quad \forall u, v \in C \quad (9.41)$$

then the algorithm (9.39) converges strongly to the solution (9.38) for all  $0 < \lambda < 2\alpha/c$ .

*Proof.*  $S$  is strictly contractive for  $0 < \lambda < 2\alpha/c$ . Indeed

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$$\begin{aligned} \|S u - S v\|^2 &= \|u - v\|^2 - 2\lambda(A(u) - A(v), u - v) + \\ &+ \lambda^2 \|A(u) - A(v)\|^2 \leq \|u - v\|^2 (1 - 2\lambda\alpha + \lambda^2 c^2) \end{aligned}$$

by (9.40) and (9.41) and

$$1 - 2\lambda\alpha + \lambda^2 c^2 < 1 \quad \text{for } 0 < \lambda < \frac{2\alpha}{c^2}$$

From exercise 4, we know that  $P_C$  is firmly contractive. Therefore  $P_C S$  is strictly contractive for  $0 < \lambda < 2\alpha/c^2$ . Thus the algorithm (9.39) converges strongly to the solution of (9.38).  $\square$

We will now give a condition on  $A$  which will imply weak convergence of the algorithm (9.39).

**THEOREM 10.** *If  $A^{-1}$  is coercive, namely*

$$(A(u) - A(v), u - v) \geq \alpha \|A(u) - A(v)\|^2 \quad \forall u, v \in C \quad (9.42)$$

*then for  $0 < \lambda < 2\alpha$  the algorithm (9.39) converges weakly to a solution of (9.38).*

*Proof.* We claim that  $S$  is quasi firmly contractive for  $\alpha < \lambda < 2\alpha$ . In fact,

$$\begin{aligned} \|Su - Sv\|^2 &\leq \|u - v\|^2 + \left(\frac{\lambda^2}{\alpha} - 2\lambda\right)(A(u) - A(v), u - v) \quad \text{by (9.42)} \\ &= (1 - \theta) \|u - v\|^2 + \theta(Su - Sv, u - v) \end{aligned} \quad (9.43)$$

181 where

$$\theta = 2 - \lambda/\alpha.$$

When  $\alpha < \lambda < 2\alpha$ , we have  $0 < \theta < 1$ .

When  $0 < \lambda < \alpha$  we obtain  $S$  to be firmly contractive. To prove this use (9.43), the Schwarz inequality and the fact that  $\theta \in [1, 2]$  when  $0 < \lambda < \alpha$ . Thus  $S$  is quasi firmly contractive for  $0 < \lambda < 2\alpha$ . Using Theorem 5, we obtain the conclusion of the Theorem.  $\square$

**REMARK 4.** *When  $A$  satisfies (9.42) and  $0 < \lambda < 2\alpha$ , we obtain from the proof of Theorem 5 that*

$$\lambda A(x^n) = x^n - Sx^n \rightarrow x - Sx = \lambda A(x),$$

$$\text{i.e.} \quad A(x^n) \rightarrow A(x), \quad (\text{Strong convergence})$$

whereas

$$x^n \rightarrow x. \quad (\text{Weak convergence})$$

*We also notice that  $x - Sx$  is unique and therefore  $A(x)$  ( $x$ , the solution of (9.38), need not be unique).*



**EXAMPLE 6.** Let

$f : H \rightarrow \mathbb{R}$  be convex, *l.s.c.* differentiable and

$A : V \rightarrow H$  be a linear operator.

We want to solve

$$\inf_{v \in V} f(Av). \quad (9.44)$$

Note that (9.44) is equivalent to

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$$\inf_{y \in C} f(y), \quad (9.45)$$

where  $C = R(A)$ , the range of  $A$ .

Now apply the algorithm (9.39). The projection on  $C$  is easy to compute. In fact,

$$P_C = A(A^*A)^{-1} A^*.$$

The nonlinear Dirichlet problem and the Minimal surface problem are particular cases of the abstract problem.

**EXAMPLE 7.** Let us consider

$$\begin{aligned} \frac{\partial u}{\partial t} + Au &= 0, \\ u(0) &= u^\circ, \end{aligned} \quad (9.46)$$

which has a solution provided  $A$  is maximal monotone.

We like to solve this problem via the algorithm

$$u^{n+1} = F(\lambda) u^n. \quad (9.47)$$

In BREZIS [4] one can find the proof of

**THEOREM 11.** *If  $F(\lambda)$  is a contraction and if*

$$\lim_{\lambda \rightarrow 0} \frac{x - F(\lambda)x}{\lambda} = A(x) \quad \text{exists,} \quad (9.48)$$

then

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$$(F(t/n))^n u^\circ \rightarrow u(t) \quad \text{uniformly.}$$

**Applications:**

1. If  $F(\lambda) = I - \lambda A$ , where  $A$  satisfies (9.40), then  $F(\lambda)$  is a contraction for  $0 < \lambda < 2\alpha$ . The limit in (9.48) exists and hence the algorithm (9.47) converges.
2. Let  $F(\lambda) = (I + \lambda A)^{-1}$ , where  $A$  is maximal monotone. Then by Lemma 7,  $F(\lambda)$  is firmly contractive and hence contractive. Existence of the limit (9.48) is proved in BRÉZIS [4] In this case also algorithm (9.47) converges.

**REMARK 5.** Theorem 11 can be used to prove that  $F(\lambda) = P_C(I - \lambda A)$  gives a sequence converging to the solution of

$$\begin{aligned} \left( \frac{du}{dt} + Au, v - u \right) &\geq 0 \quad \forall v \in C, \\ u(0) &= u^\circ. \end{aligned}$$

**REMARK 6.** If  $A$  is linear, monotone and closed then  $A$  is maximal monotone.

**EXAMPLE 8. The Flow of a Bingham Fluid:** Consider the problem:  
Find  $\sigma \in (L^2(\Omega))^n$ ,  $u \in H_0^1(\Omega)$  such that

$$J'(\sigma) - \nabla u = 0 \tag{9.49}$$

$$(\sigma, \nabla v) = (f, v) \quad \forall v \in H_0^1(\Omega),$$

184 where

$$\begin{aligned} J(\sigma) &= \frac{1}{2} \|\sigma - P_K \sigma\|^2_{(L^2(\Omega))^n}, \\ K &= \{\sigma \in (L^2(\Omega))^n : |\sigma(x)| \leq 1 \text{ a. e. in } \Omega\}. \end{aligned}$$

It is possible to prove that (9.49) is equivalent to the Bingham flow given in Example 4. It can be proved that

$$J'(\sigma) = \sigma - P_K \sigma.$$

Let

$$Z(f) = \{\sigma \in (L^2(\Omega))^2 : (\sigma, \nabla v) = (f, v) \quad \forall v \in H_0^1(\Omega)\}$$

If  $\sigma$  is a solution of (9.49), then  $\sigma$  is also a solution of

$$J(\sigma) = \inf_{\tau \in Z(f)} J(\tau). \tag{9.50}$$

Therefore, we can apply previous results. Note that  $J'$  satisfies (9.42) with  $\alpha = 1$ , so that previous results can be applied. Note also that the projection on  $Z(f)$  is easy to compute:

$$P_{Z(f)}(\sigma) = \sigma + \nabla(-\Delta)^{-1} \operatorname{div} \sigma + \nabla(-\Delta)^{-1} f.$$

**EXAMPLE 9. Transonic Flows.**

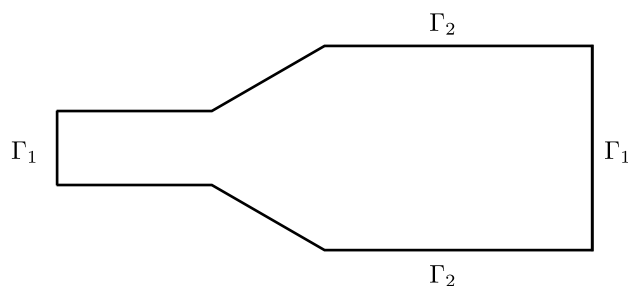


Figure 9.3:

The potential flow at transonic speed in a nozzle is governed by the equations 185

$$\begin{aligned} \operatorname{div}(\rho \vec{q}) &= 0, \\ \operatorname{rot} \vec{q} &= 0, \quad \vec{q} = \nabla \phi, \\ \rho &= \left(1 - \frac{\gamma - 1}{2} M_\infty^2 (1 - q^2)\right)^{1/\gamma - 1}. \end{aligned}$$

(Practical value of  $\gamma = 1.4$ ) where  $q = |\vec{q}|$ .

Finally we solve

$$\operatorname{div}(\rho(|\nabla \phi|) \nabla \phi) = 0,$$

$$\begin{aligned}\phi|_{\Gamma_1} &= \phi_0, \\ \frac{\partial \phi}{\partial n}|_{\Gamma_2} &= 0\end{aligned}$$

Let  $M = q/a$ , where  $a = \frac{\rho^{\gamma-1}}{M_\infty^2}$ .  $M$  is called the *Mach number*.

The equation is elliptic for  $M < 1$  and hyperbolic for  $M > 1$ .

When  $M < 1$  continuous piecewise linear finite element can be used.

For  $M > 1$ , we do not know much (see COURANT- FRIEDRICHS [13]).

The reader can refer to GLOWINSKI-PIRONNEAU [20],[21], RAVIART [36], CIAVALDINI-POGU-TOURNEMINE [12], J. ROUX [40].

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