The Finite Element Method with Lagrangian Multipliers*

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Summary. The Dirichlet problem for second order differential equations is chosen as a model problem to show how the finite element method may be implemented to avoid difficulty in fulfilling essential (stable) boundary conditions. The implementation is based on the application of Lagrangian multiplier. The rate of convergence is proved.

1. Introduction

The finite element method has become the most successfull approximation method in engineering. There is a variety of detailed approaches based on the finite element method. See e.g. [22 and 17] many others. The central idea of the finite element method is to use different variational principles together with a Galerkin procedure applied to picewise smooth functions.

The notion of a variational method is used today in engineering mostly in the narrow sense that the solution is a stationary point. In other words, the bilinear form which determines the stationarity is often not positive definite. For variational principles used in the theory of elasticity, see e.g. [25].

The finite element method has been studied recently from a theoretical point of view also. Many theorems about convergence, error estimates, etc., have been proved. See e.g. [5, 6, 8, 9, 12, 14, 16, 26].

It has been shown that it is not computationally easy to handle the Dirichlet (essential) boundary condition if the variational principle requires fullfillment of these conditions. Different methods have been developed to avoid this difficulty. See e.g. [3, 10, 14, 18] and others. There is an obvious and classical technique to deal with restrictions in the variational principles. It is the theory of Lagrangian multipliers, applied to the finite element method. Nevertheless application of this idea to the finite element method has not been so far theoretically studied.

We shall analyze here a model problem, but the approach is quite general and may be applied in other cases too. Some interesting cases will be brought out in subsequent papers. Let our problem be to solve the differential equation

$$-\Delta u + u = f \quad \text{on} \quad \Omega \tag{1.1}$$

with boundary conditions

$$u = g$$
 on Ω . (1.2)

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We shall assume that Ω is a bounded domain and that its boundary Ω^{\bullet} is sufficiently smooth.

The classical technique is to minimize the quadratic functional

$$F(v) = \int_{O} \left[\sum_{i=1}^{n} \left(\frac{\partial v}{\partial x_{i}} \right)^{2} + v^{2} \right] dx - 2 \int_{O} f v \, dx \tag{1.3}$$

over all functions satisfying the prescribed boundary condition on Ω . By the theory of Lagrange multiplier (see e.g. [19]) the solution creates the stationary point for the (not positive definite) functional

$$F(v,\lambda) = \int_{\Omega} \left[\sum_{i=1}^{n} \left(\frac{\partial v}{\partial x_{i}} \right)^{2} + v^{2} \right] dx - 2 \oint_{\Omega} \lambda(v - g) ds - 2 \int_{\Omega} v f dx. \tag{1.4}$$

In the next chapters we will use this principle and show how to use it in the theory of the finite element method. We shall also show that the rate of convergence of this method is the optimal one.

2. The Principle Notions

Throughout the entire paper R_n be the *n*-dimensional Euclidian space, $x \equiv (x_1, \ldots, x_n), \|x\|^2 = \sum_{i=1}^n x_i^2$ and $dx = dx_1 \ldots dx_n$. Let $\Omega \in R_n$ be a bounded region and Ω its boundary. We will assume that Ω is infinitely many times differentiable.

Let $L_2(\Omega)$ be the space of square integrable functions u on Ω such that

$$||u||_{L_{2}(\Omega)}^{2} = \int_{\Omega} |u|^{2} dx < \infty.$$

Some times we shall write $L_2(\Omega)=H^0(\Omega)$. Let $\mathscr{E}(\bar{\Omega})$ be the space of all infinitely many times differentiable functions on Ω and such that all derivatives are continuously extenable on Ω^{\bullet} , analogously $\mathscr{E}(\Omega^{\bullet})$ is the space of all infinitely many times differentiable functions on Ω^{\bullet} . Furthermore let $\mathscr{D}(\bar{\Omega}) \subset \mathscr{E}(\Omega)$ be the subspace of all functions with compact support in Ω .

Let $l \ge 1$, l integral. The Sobolev space $H^l(\Omega)$ (resp. $H^l_0(\Omega)$) will be the closure of $\mathscr{E}(\Omega)$ (resp. $\mathscr{D}(\overline{\Omega})$) in the norm $\|\cdot\|_{H^l(\Omega)}$ where

$$||u||_{H^{l}(\Omega)}^{2} = \sum_{0 \le |\alpha| \le l} ||D^{\alpha}u||_{L_{2}(\Omega)}^{2}$$

and

$$D^{\alpha} = \frac{\partial^{\alpha_1 + \cdots + \alpha_n}}{\partial x_1^{\alpha_1} \cdots \partial x_n^{\alpha_n}}, \quad \alpha = (\alpha_1, \dots, \alpha_n), \quad \alpha_i \ge 0, \quad |\alpha| = \sum_{i=1}^n \alpha_i.$$

Now let $\psi_i(x) \ge 0$, $i = 1, 2, ..., \nu$, $x \in R_n$ be functions infinitely times differentiable and such that

$$\sum_{i=1}^{y} \psi_i(x) = 1 \quad \text{for} \quad x \in \Omega^*.$$

Further let there be a system of local coordinates $x_i^{[s]}$, $i=1,\ldots,n$ $s=1,\ldots,\nu$ and n-1 dimensional domains $J_s \subset R_{m-1}$, and a system of functions φ_s , $s=1,\ldots,\nu$

defined on J_s which are infinitely times differentiable, and which create one to one mappings χ_s of J_s onto $\chi_s(J_s) = \Omega_s^*$ where

$$Q_{s}^{\bullet} = E\left[\left(x_{1}^{[s]}, \ldots, x_{n-1}^{[s]}, \varphi(x_{1}^{[s]}, \ldots, x_{n-1}^{[s]}) \middle| (x_{1}^{[s]}, \ldots, x_{n-1}^{[s]}) \in J_{s}\right]$$

and such that

$$E[x \in \Omega^{\bullet}|\psi_s(x) > 0] \in E[(x_1^{[s]}, \ldots, x_{n-1}^{[s]}, \varphi(x_1^{[s]}, \ldots, x_{n-1}^{[s]})|(x_1^{[s]}, \ldots, x_{n-1}^{[s]}) \in (J_s)_H]$$

where H > 0 and

$$(J_s)_H = E[x \in J_s | d(x, J_s^{\bullet}) > H]$$

with $d(x, J_s^*)$ as the distance of x to J_s^* . it is easy to see that such a system of functions ψ_i , ψ_i and domains J_s actually exist. Let f be defined on Ω^* . Then the function $f_s = \psi_s f = 0$ everywhere outside of Ω_s^* and $f_s(\chi_s(x))$ is defined on J_s and has compact support.

Let us introduce the Sobolev spaces on Ω^{\bullet} . Let $l \geq 0$ an integer. The Sobolev space $H^{l}(\Omega^{\bullet})$ with norm $\|\cdot\|_{H^{l}(\Omega)}$ is the space of all functions f defined on Ω^{\bullet} and such that

$$||f||_{H^{1}(\Omega')}^{2} = \sum_{s=1}^{\nu} ||f_{s}||_{H^{1}(J_{s})}^{2} < \infty.$$

We have introduced the Sobolev spaces with integral derivatives. For Ω with $\Omega \cdot \in C^{\infty}$ we may construct Hilbert scales and obtain spaces with fractional derivatives. See e.g. [20 and 21]. Our notation is in agreement with [21]. These interpolated spaces (resp., their norms) are equivalent with Aronszajn-Slobodeckii norm. See e.g. [23]. For $0 < \alpha = [\alpha] + \sigma$, $0 < \sigma < 1$ and $[\alpha]$ integral we may define

$$||u||_{H^{\alpha}(\Omega)}^{2} = ||u||_{H^{\alpha}(\Omega)}^{2} + \sum_{|k|=[\alpha]} ||D^{k}u||_{H^{\sigma}(\Omega)}$$
 (2.1)

where

$$||u||_{H^{\sigma}(\Omega)}^{2} = \int_{\Omega} \int_{\Omega} \frac{(u(t) - u(\tau))^{2}}{(t - \tau)^{n+2\sigma}} dt d\tau.$$
 (2.2)

Analogously we define the norm $\|\cdot\|_{H^{\alpha}(\Omega^{\bullet})}$ which is equivalent with the interpolated norm. For α negative we shall define the space $H^{\alpha}(\Omega^{\bullet})$ as a dual space, namely $H^{\alpha}(\Omega^{\bullet}) = (H^{-\alpha}(\Omega^{\bullet}))^1$. For additional information see [21] p. 35.

Let us mention some theorems which will be useful later.

Theorem 2.1. Let $f \in H^k(\Omega)$, $k > \frac{1}{2}$. Then there exists a trace of the function f on Ω^* and

$$||f||_{H^{k-\frac{1}{2}}(\Omega)} \le C ||f||_{H^{k}(\Omega)}^{1}$$
 (2.3)

where C does not depend on f.

Theorem 2.2. Let $f \in H^k(\Omega)$, $k > \frac{3}{2}$. Then there exists a trace $\partial f/\partial n$ an Ω and

$$\left\| \frac{\partial f}{\partial n} \right\|_{H^{k-\frac{3}{2}}(\Omega')} \le C \|f\|_{H^{k}(\Omega)} \tag{2.4}$$

where C does not depend on f. For the proof see [21] p. 47.

¹ C will be a generic constant with different values on different places.

The theorems hold only for $k > \frac{1}{2}$ (resp. $k > \frac{3}{2}$) and it is possible to show that the theorem is not valid for $k \leq \frac{1}{2}$ (resp. $k \leq \frac{3}{2}$). See e.g. [21] p. 49. Nevertheless let us show that the theorem is true if we restrict the space $H^k(\Omega)$ to a smaller one.

Let $\mathscr{S}(\Omega) \subset H^1(\Omega)$ be the space of all functions which satisfy the equation

$$-\Delta u + u = 0$$

in the weak sense; i.e., $\mathscr{S}(\Omega)$ be such a subspace of functions u that

$$\int\limits_{\Omega} \left[\sum_{i=1}^{n} \frac{\partial u}{\partial x_{i}} \frac{\partial v}{\partial x_{i}} + uv \right] dx = 0$$

for every $v \in H_0^1(\Omega)$. Now we may prove the next theorem.

Theorem 2.3. Let $u \in \mathcal{S}(\Omega)$. Then we have $\partial u/\partial n \in H^{-\frac{1}{2}}(\Omega^{\bullet})$ and

$$\left\|\frac{\partial u}{\partial n}\right\|_{H^{-\frac{1}{2}}(\Omega)} \le C \|u\|_{H^{1}(\Omega)} \tag{2.5}$$

where C does not depend on u.

Proof. Because $\mathscr{S}(\Omega) \cap \mathscr{E}(\overline{\Omega})$ is dense in $\mathscr{S}(\Omega)$ (see e.g. [1]) we have to prove the for any $v \in H^{\frac{1}{4}}(\Omega^{\bullet}) \cap \mathscr{E}(\Omega^{\bullet})$ and $u \in \mathscr{S}(\Omega) \cap \mathscr{E}(\overline{\Omega})$ we have

$$\left| \oint_{\Omega^{\bullet}} \frac{\partial u}{\partial n} v \, ds \right| \leq C(u) \|v\|_{H^{\frac{1}{2}}(\Omega^{\bullet})} \tag{2.6}$$

with

$$C(u) \leq C \|u\|_{H^{\frac{1}{2}}(\Omega)}.$$

By an inverse imbedding theorem (see e.g. [21]) there exists a linear mapping μ of $H^{\frac{1}{2}}(\Omega^{\bullet})$ into $H^{1}(\Omega)$ such that $\mu(u) = u$ on Ω^{\bullet} and $\|\mu(u)\|_{H^{1}(\Omega)} \leq C \|u\|_{H^{\frac{1}{2}}(\Omega^{\bullet})}$ where C does not depend on u. For $u \in \mathscr{S}(\Omega) \cap \mathscr{E}(\overline{\Omega})$ and $v \in \mathscr{E}(\Omega^{\bullet})$ we compute

$$\int_{\Omega} \left[\sum_{i=1}^{n} \frac{\partial u}{\partial x_{i}} \frac{\partial \mu(v)}{\partial x_{i}} + u \mu(v) \right] dx = B(u, v).$$
 (2.7)

Integrating by parts we obtain

$$B(u, v) = \oint_{\mathcal{O}} \frac{\partial u}{\partial n} v ds$$

and we have obviously

$$|B(u, v)| \le ||u||_{H^1(\Omega)} ||\mu(v)||_{H^1(\Omega)} \le C ||u||_{H^1(\Omega)} ||v||_{H^{\frac{1}{2}}(\Omega)}.$$

Therefore we have

$$|B(u,v)| = \left| \oint_{\Omega^*} \frac{\partial u}{\partial n} v ds \right| \le C ||u||_{H^1(\Omega)} ||v||_{H^{\frac{1}{2}}(\Omega^*)}$$

and this in fact proves our theorem.

Let us mention now some theorems about the Dirichlet and Neumann problems.

Theorem 2.4. Let $f \in H^k(\Omega)$, $k \ge 0$, $g \in H^l(\Omega^{\bullet})$, $l \ge \frac{1}{2}$. Then there exists in $H^1(\Omega)$ exactly one function u such that u is a weak solution of the equation $-\Delta u + u = f$

with boundary condition u = g on Ω^* . Furthermore $u \in H^s(\Omega)$ where $s = \min(k+2, l+\frac{1}{2})$ and

$$||u||_{H^{s}(\Omega)} \le C \left[||f||_{H^{k}(\Omega)} + ||g||_{H^{t}(\Omega^{*})} \right]$$
 (2.8)

where C does not depend on f or g.

Theorem 2.5. Let $f \in H^k(\Omega)$, $k \ge 0$, $g \in H^l(\Omega^*)$, $l \ge -\frac{1}{2}$. Then there exists in $H^1(\Omega)$ exactly function u such that u is a weak solution of the equation $-\Delta u + u = f$ with boundary condition $\partial u/\partial n = g$. Furthermore $u \in H^s(\Omega)$ where

$$s = \min(k+2, l+\frac{3}{2})$$

and

$$||u||_{H^{s}(\Omega)} \le C \left[||f||_{H^{k}(\Omega)} + ||g||_{H^{l}(\Omega^{s})} \right]$$
 (2.9)

where C does not depend on f or g.

For the proof of these theorems see [21], p. 203.

We proved Theorem 2.3 for $u \in \mathcal{S}(\Omega)$ but it is possible to generalize it.

Now let $f \in L_2(\Omega)$ and let $\mathscr{S}(\Omega, f)$ be the space of all functions which satisfy the equation

$$-\Delta u + u = f$$

in the weak sense i.e. let $\mathscr{S}(\Omega)$ be such a subspace of function $u \in H^1(\Omega)$ that

$$\int_{\Omega} \left[\sum_{i=1}^{n} \frac{\partial u}{\partial x_{i}} \frac{\partial v}{\partial x_{i}} + uv \right] dx = \int_{\Omega} f v dx$$

for every $v \in H^1_0(\Omega)$. Now the following theorem is valid.

Theorem 2.6. Let $u \in \mathcal{S}(\Omega, f)$. Then we have $\partial u/\partial n \in H^{-\frac{1}{2}}(\Omega)$ and

$$\left\| \frac{\partial u}{\partial n} \right\|_{H^{-\frac{1}{2}}(\Omega^*)} \le C \left[\| u \|_{H^1(\Omega)} + \| f \|_{L_{\bullet}(\Omega)} \right] \tag{2.10}$$

where C does not depend on u or t.

Proof. Let $u \in \mathcal{S}(\Omega, f)$. Then there exists a linear mapping χ of $L_2(\Omega)$ into $H^2(\Omega)$ such that u = v + w where

$$w = \chi(f), \quad ||w||_{H^{2}(\Omega)} \leq C ||f||_{L_{2}(\Omega)}$$

and

$$v \in \mathcal{S}(\Omega), \quad ||v||_{H^1(\Omega)} \leq C [||u||_{H^1(\Omega)} + ||f||_{L_2(\Omega)}].$$

In fact w is a particular solution of the equation $-\Delta w + w = f$ such that $\|w\|_{H^1(\Omega)} \le C \|f\|_{L_1(\Omega)}$ and we may obviously use for our purpose Theorem 2.4. Combining now Theorem 2.3 with Theorem 2.2, we get our result.

Theorem 2.7. Let $g \in H^{-\frac{1}{2}}(\Omega^{\bullet})$ and u be the weak solution of the Neumann problem for the equation $-\Delta u + u = 0$ on Ω , $\partial u/\partial n = g$ on Ω^{\bullet} in $H^{1}(\Omega)$ [see Theorem 2.5]. There exists constants $0 < C_{1} < C_{2} < \infty$ such that

$$C_1 \oint_{\Omega} guds \le ||g||_{H}^{2-\frac{1}{2}}(\Omega) \le C_2 \oint_{\Omega} guds$$
 (2.11)

and

$$||u||_{H^1(\Omega)}^2 = \oint_{\Omega} g u \, ds \tag{2.12}$$

Proof. The statement (2.12) follows immediately from the definition of a weak solution. For $g \in H^0(\Omega^{\bullet})$ we have

$$||g||_{H^{-\frac{1}{2}}(\Omega^{*})} = \sup_{v \in H^{\frac{1}{2}}(\Omega^{*})} \frac{\left| \oint_{\Omega^{*}} g v \, ds \right|}{||v||_{H^{\frac{1}{2}}(\Omega)}}. \tag{2.13}$$

Let w be the solution of the equation $-\Delta w + w = 0$ with w = v on Ω^* . Using Theorem 2.4 and 2.5 and the definition of the weak solution, we have

$$\frac{\left| \oint_{\Omega^*} g v d s \right|}{\|v\|_{H^{\frac{1}{2}}(\Omega^*)}} \leq C \frac{\|w\|_{H^1(\Omega)} \|u\|_{H^1(\Omega)}}{\|w\|_{H^1(\Omega)}}
\leq C \|u\|_{H^1(\Omega)} = C \left[\oint_{\Omega^*} g u d s \right]^{\frac{1}{2}}.$$
(2.14)

Therefore, we have

$$||g||_{H^{-\frac{1}{2}}(\Omega)}^2 \leq C \oint_{\Omega^*} guds.$$

Because of the density of $H^0(\Omega^{\bullet})$ in $H^{-\frac{1}{4}}(\Omega^{\bullet})$, one side of (2.11) is proved.

Let us prove now the other side of the enequality (2.11). We have

$$\left| \oint_{\Omega^{\bullet}} g u ds \right| \leq \|g\|_{H^{-\frac{1}{2}}(\Omega^{\bullet})} \|u\|_{H^{\frac{1}{2}}(\Omega^{\bullet})}$$

by definition of the norm $\|g\|_{H^{-\frac{1}{4}}(\Omega^{\bullet})}$. Using Theorem 2.5 we get

$$\left| \oint_{\Omega^{\bullet}} g u ds \right| \leq C \|g\|_{H^{-\frac{1}{2}}(\Omega^{\bullet})}^{2}.$$

We in fact show that

$$C_1 \left| \oint_{\Omega^*} g u ds \right| \le \|g\|_{H^{-\frac{1}{2}}(\Omega^*)}^2 \le C_2 \left| \oint_{\Omega^*} g u ds \right| \tag{2.15}$$

instead of (2.11). But using (2.12) we see that $\oint guds \ge 0$ and therefore (2.15) is identical to (2.11).

Let us introduce now two theorems which will play a fundamental role in the next section.

Theorem 2.8. Let H_1 and H_2 be two Hilbert spaces with scalar products $(\cdot, \cdot)_{H_1}$ and $(\cdot, \cdot)_{H_1}$ resp. Further let B(u, v) be a bilinear form on $H_1 \times H_2$, $u \in H_1$, $v \in H_2$ such that

$$|B(u,v)| \le C_1 ||u||_{H_1} ||v||_{H_2},$$
 (2.16)

$$\sup_{\substack{u \in H_1 \\ \|u\|_{H_1} \le 1}} |B(u, v)| \ge C_2 \|v\|_{H_2}, \tag{2.17}$$

$$\sup_{\substack{v \in H_3 \\ v \mid_{H_2} \le 1}} |B(u, v)| \ge C_3 ||u||_{H_1}, \tag{2.18}$$

with $C_2 > 0$, $C_3 > 0$, $C_1 < \infty$. Further let $f \in H_2'$ i.e., let f be a linear functional on H_2 . Then there exists exactly one element $u_0 \in H_1$ such that

$$B(u_0, v) = f(v) (2.19)$$

for all $v \in H_2$ and

$$\|u_0\|_{H_1} \le \frac{\|f\|_{H_2^*}}{C_q}$$
 (2.20)

For the proof see [9].

Theorem 2.9. Let the assumptions of the Theorem 2.8 be fulfilled. Further let there be given two linear subspaces (closed) $M_1 \in H_1$ and $M_2 \in H_2$ and for every $v \in M_2$, let

$$\sup_{\substack{u \in M_1 \\ \|u\|_{H_1} \le 1}} |B(u, v)| \ge d_2(M_1, M_2) \|v\|_{H_1}$$
 (2.21)

with $d_2(M_1, M_2) > 0$, and for every $u \in M_1$

$$\sup_{\substack{v \in M_1 \\ \|v\|_{H_1} \le 1}} |B(u, v)| \ge d_3(M_1, M_2) \|u\|_{H_1}$$
(2.22)

with $d_3(M_1, M_2) > 0$. Let $f \in H_2'$ be given, and let u_0 denote the element of H_1 , such that

$$B(u_0, v) = f(v) (2.23)$$

holds for all $v \in H_2$ [such an element exists and is unique by Theorem 2.8].

Let there exist a $\omega \in M_1$ such that

$$\|u_0 - \omega\|_{H_s} \le \vartheta. \tag{2.24}$$

Further, let $\hat{u}_0 \in M_1$ be such that

$$B(\hat{u}_0, v) = f(v)$$
 (2.25)

for all $v \in M_2$. Then,

$$\|u_0 - \hat{u}_0\|_{H_1} \le \left[1 + \frac{C_1}{d_3(M_1, M_2)}\right] \vartheta.$$
 (2.26)

For the proof see [9].

3. The Finite Element Method — Rate of Convergence in $H^1(\Omega)$

Before we discuss the finite element method, we shall introduce some necessary machinery. Let us introduce $\mathscr{S}_{h}^{t,k}(\Omega)$ and $\mathscr{S}_{h}^{t,k}(\Omega^{\bullet})$ as one parameter families of functions for all 0 < h < 1. The linear, finite dimensional system of functions $\mathscr{S}_{h}^{t,k}(\Omega)$ [resp. $\mathscr{S}_{h}^{t,k}(\Omega^{\bullet})$] will be called a (t,k)-regular system for $t \geq k \geq 0$ (resp. $t \geq k \geq -\frac{1}{2}$) if:

- $(1) \qquad \mathscr{S}^{t,k}_h(\Omega) \subset H^k(\Omega) \text{ [resp. } \mathscr{S}^{t,k}_h(\Omega^{\bullet}) \subset H^k(\Omega^{\bullet})],$
- (2) if $f \in H^1(\Omega)$ [resp. $f \in H^1(\Omega^{\bullet})$]

then there exists $g \in \mathcal{G}_{h}^{t,h}(\Omega)$ [resp. $g \in \mathcal{G}_{h}^{t,h}(\Omega^{\bullet})$] such that for any $0 \le s \le h \le l$ [resp. $-\frac{1}{2} \le s \le h \le l$]

$$\|g - f\|_{H^{\bullet}(\Omega)} \le C h^{\mu} \|f\|_{H^{1}(\Omega)} \quad [\text{resp. } \|g - f\|_{H^{\bullet}(\Omega^{\bullet})} \le C h^{\mu} \|f\|_{H^{1}(\Omega^{\bullet})}]$$
 (3.1)

where

$$\mu = \min(l-s, t-s)$$

and C does not depend on s, h or f.

The system $\mathscr{S}_{h}^{t,k}(\Omega^{\bullet})$ will be called strongly (t,k) regular if it is (t,k)-regular and if for every $g \in \mathscr{S}_{h}^{t,k}$, $-\frac{1}{2} \leq q \leq s \leq k$

$$||g||_{H^{q}(\Omega^{\bullet})} \le C h^{-(s-q)} ||g||_{H^{q}(\Omega^{\bullet})}$$

and C does not depend on s, h or f.

We have studied the construction of the systems $\mathcal{S}_h^{t,k}(\Omega)$, in [4] and [11] and $\mathcal{S}_h^{t,k}(\Omega^{\bullet})$ in [11]. Similarily, since we have introduced the norms $H^s(\Omega^{\bullet})$ by transformation to the case of n-1 dimensional domains, we may construct strongly (t,k)-regular system $\mathcal{S}_h^{t,k}(\Omega^{\bullet})$ using principally the results of [4]. Let us remark that we do not introduce the strongly regular systems $\mathcal{S}_h^{t,k}(\Omega)$. These systems are much more difficult to construct (see [7]). Let us mention that these spaces are finite dimensional.

Now we may explain the finite element method for solving the Dirichlet problem

$$-\Delta u + u = f \quad \text{on} \quad \Omega, \tag{3.2}$$

$$u = g$$
 on Ω . (3.3)

with $f \in L_2(\Omega)$ and $g \in H^{\frac{1}{2}}(\Omega^{\bullet})$. Based on the form (1.4) we may introduce a bilinear form

$$B(u,\lambda;v,\mu) = \int_{\Omega} \left(\sum_{i=1}^{n} \frac{\partial u}{\partial x_{i}} \frac{\partial v}{\partial x_{i}} + uv \right) dx - \int_{\Omega^{*}} (\lambda v + u\mu) ds$$
 (3.4)

and two functionals

$$F(u,\lambda) = \int_{\Omega} f u \, dx, \tag{3.5}$$

$$G(u, \lambda) = -\oint_{\Omega} g \lambda ds. \tag{3.6}$$

Then the stationary point (u_0, λ_0) of (1.4) is obviously such that for every (v, μ) , we have

$$B(u_0, \lambda_0; v, \mu) = F(v, \mu) + G(v, \mu).$$
 (3.7)

Let us define precisely the domain of definition of the bilinear form (3.4). Let us have $H = H_1 = H_2 = H^1(\Omega) \times H^{-\frac{1}{2}}(\Omega^{\bullet})$ with the norm $\|u, \lambda\|_H^2 = \|u\|_{H^1(\Omega)}^2 + \|\lambda\|_H^2 - \frac{1}{2}(\Omega^{\bullet})$ and let us show that the bilinear form (3.4) is defined on $H \times H$ and fulfills the conditions (2.16), (2.17) and (2.18) in Theorem 2.8. Let us study first the condition (2.16). Using the imbedding Theorem 2.1, we have

$$|B(u,\lambda;v,\mu)| \leq ||u||_{H^{1}(\Omega)} ||v||_{H^{1}(\Omega)} + C[||\lambda||_{H^{-\frac{1}{2}}(\Omega^{*})} ||v||_{H^{1}(\Omega)} + ||\mu||_{H^{-\frac{1}{2}}(\Omega^{*})} ||u||_{H^{1}(\Omega)}]$$

and therefore we get

$$|B(u, \lambda; v, \mu)| \le C \|(u, \lambda)\|_H \|v, \mu\|_H$$
 (3.8)

and therefore (2.16) is proved. Because of symmetry by proving (2.18) we will prove (2.17) too.

So let $(u, \lambda) \in H$ be given. Denote by $w \in H^1(\Omega)$ the solution of the Neumann problem for the differential equation $-\Delta w + w = 0$ with the boundary condition $\partial w/\partial n = \lambda$. Using Theorem 2.5, the function $w \in H^1(\Omega)$ exists and for all $v \in H^1(\Omega)$ we have

$$\int_{\Omega} \left[\sum_{i=1}^{n} \frac{\partial w}{\partial x_{i}} \frac{\partial v}{\partial x_{i}} + v w \right] dx = \oint_{\Omega} \lambda v ds.$$
 (3.9)

Furthermore we know that

$$||w||_{H^1(\Omega)} \le C ||\lambda||_{H^{-\frac{1}{2}}(\Omega)}.$$
 (3.10)

Using Theorem 2.7, we have

$$C_1 \|\lambda\|_{H^{-\frac{1}{2}}(\Omega^*)}^2 \le \oint_{\Omega^*} w \,\lambda ds \tag{3.11}$$

and $0 < C_1$. Let us take v = u - w and $\mu = -2\lambda$. Then obviously,

$$\|(v,\mu)\|_{H} \le C \|(u,\lambda)\|_{H}.$$
 (3.12)

Let us show now that

$$B(u, \lambda; v, \mu) \ge C \|(u, \lambda)\|_{H}^{2}.$$
 (3.13)

If (3.13) holds them using (3.12), we get (2.18). So we have

$$B(u, \lambda; v, \mu) = \int_{\Omega} \left[\sum_{i=1}^{n} \left(\frac{\partial u}{\partial x_{i}} \right)^{2} + u^{2} \right] dx$$

$$- \int_{\Omega} \left[\sum_{i=1}^{n} \frac{\partial u}{\partial x} \frac{\partial w}{\partial x} + u w \right] dx$$

$$- \oint_{\Omega} (\lambda u + u \mu) ds + \oint_{\Omega} \lambda w ds.$$
(3.14)

Using (3.9), we get

$$B(u,\lambda;v\mu) = \|u\|_{H^1(\Omega)}^2 - \oint_{\Omega} u(2\lambda + \mu) ds + \oint_{\Omega} \lambda w ds.$$
 (3.15)

Because $2\lambda + \mu = 0$, we obtain (3.13) from (3.11). Therefore we proved the following theorem.

Theorem 3.1. Let $H = H_1 = H_2 = H_1^1(\Omega) \times H^{-\frac{1}{2}}(\Omega^{\bullet})$. Then the bilinear form (3.4) satisfies the assumptions of Theorem 2.8.

Now let $f \in L_2(\Omega)$ and $g \in H^{\frac{1}{2}}(\Omega^*)$. Then obviously the functionals (3.5) and (3.6) are continuous. Using Theorem 2.8, we may find $(u_0, \lambda_0) \in H$ such that (3.7) holds for every $(v, \mu) \in H$. Let us investigate the connection between these functions and our original problem (3.2)-(3.3). Let us state it as a theorem.

Theorem 3.2. Let $f \in L_2(\Omega)$, $g \in H^{\frac{1}{2}}(\Omega^{\bullet})$ and let $(u_0, \lambda_0) \in H$ be such that (3.7) holds for all $(v, \mu) \in H$. Further let $w \in H^1(\Omega)$ be the weak solution of the Dirichlet problem (3.2), (3.3). Then $u_0 = w$ and $\lambda_0 = \partial w / \partial n$.

Proof. First let us state that there exists a solution of the Dirichlet problem (3.2) and (3.3) by Theorem 2.4. Further, it is obvious that $w \in \mathcal{S}(\Omega, f)$ and

therefore $\partial w/\partial n \in H^{-\frac{1}{2}}(\Omega^{\bullet})$ by Theorem 2.6. Let us show now that

$$B\left(w,\frac{\partial w}{\partial n};v,\mu\right) = \int_{\Omega} fv \, dx - \oint_{\Omega^*} g\mu \, ds.$$

Because of Theorem 2.5, the function w is the solution of the Neumann problem for the differential equation $-\Delta u + u = f$ and boundary condition $\partial u/\partial n = \partial w/\partial n$. Therefore

$$\int_{\Omega} \left[\sum_{i=1}^{n} \frac{\partial w}{\partial x_{i}} \frac{\partial v}{\partial x_{i}} + wv \right] dx = \oint_{\Omega} v \frac{\partial w}{\partial n} ds + \int_{\Omega} f v dx.$$

Thus we have

$$B\left(w, \frac{\partial w}{\partial n}; v, \mu\right) = \oint_{\Omega} v \frac{\partial w}{\partial n} ds - \oint_{\Omega} \frac{\partial w}{\partial n} v ds$$
$$- \oint_{\Omega} w \mu ds + \int_{\Omega} f v dx = - \oint_{\Omega} g \mu ds + \int_{\Omega} f v dx$$

which is what we wanted to be proved.

To use out bilinear for the finite element approach, let us introduce a subspace $M \subset H$ so that

$$M = \mathscr{S}^{t_1, k_1}_{h_1}(\Omega) \times \mathscr{S}^{t_1, k_2}_{h_2}(\Omega^{\bullet})$$

where $\mathscr{S}_{h_1}^{t_1,h_2}(\Omega)$ is (t_1,k_1) -regular system and $\mathscr{S}_{h_2}^{t_1,h_2}(\Omega^{\bullet})$ is a strongly (t_2,k_2) -regular system. In addition, let $M=M_1=M_2$. Further let us assume that $k_1 \geq 1$, $k_2 \geq \frac{1}{2}$ and furthermore that

$$h_2 \ge Kh_1 \quad K > 0 \tag{3.16}$$

where K is independent of h_1 and will be determined later. Now we may use Theorem 2.9 which yields the finite element method with Langrage multipliers. We shall seek $(\bar{u}_0, \bar{\lambda}_0) \in M$ such that (3.7) is fulfilled for every $(v, \mu) \in M$. This condition gives a finite number of linear conditions for the element $(\bar{u}_0, \bar{\lambda}_0)$ which is determined by a finite number of parameters. We shall show that the system of linear algebraic equations which we have to solve has, under certain assumptions, exactly one solution. Thus will determine the approximate solution.

Let us prove

Theorem 3.3. For K sufficiently large (see (3.16)) we have in the Theorem 2.9 $d_2(M_1, M_2) = d_3(M_1, M_2) > C > 0$ and C does not depend on h_1 .

Proof. The proof will be similar to the proof of Theorem 3.1. Let $(u, \lambda) \in M$ be given. Denote by $w \in H^1(\Omega)$ the solution of the Neumann problem for the differential equation $-\Delta w + w = 0$ and the boundary condition $\partial w / \partial n = \lambda$. Because $\lambda \in \mathscr{S}^{i_1, i_2}_{h_2}(\Omega^*)$ and because the subspace $\mathscr{S}^{i_1, i_2}_{h_2}(\Omega)$ is strongly regular, we have

$$\|\lambda\|_{H^{\frac{1}{2}}(\Omega^{*})} \le \frac{C}{h_{0}} \|\lambda\|_{H^{-\frac{1}{2}}(\Omega^{*})}.$$
 (3.17)

Using Theorem 2.5, we have

$$||w||_{H^{3}(\Omega)} \le \frac{CC_{1}}{h_{2}} ||\lambda||_{H^{-\frac{1}{2}}(\Omega^{*})}$$
 (3.18)

where C_1 is the constant in (2.9). Using the fact that $\mathcal{S}_{h_1}^{t_1,k_1}(\Omega)$ is (t_1, k_1) -regular we may find $z \in \mathcal{S}_{h_1}^{t_1,k_1}(\Omega)$ such that

$$||w - z||_{H^{1}(\Omega)} \leq C ||w||_{H^{2}(\Omega)} h_{1}$$

$$\leq C C_{1} C_{2} \frac{h_{1}}{h_{2}} ||\lambda||_{H^{-\frac{1}{2}}(\Omega')}.$$
(3.19)

We know that

$$\oint_{\Omega} w \lambda ds \ge C_3 \|\lambda\|_{H^{-\frac{1}{2}}(\Omega^*)}^2. \tag{3.20}$$

Therefore

$$\oint_{\Omega'} z \lambda ds = \oint_{\Omega'} w \lambda ds + \oint_{\Omega'} \lambda (z - w) ds$$

$$\ge C_3 \|\lambda\|_{H^{-\frac{1}{2}}(\Omega')}^2 - \eta$$
(3.21)

where

$$\begin{split} 0 < \eta < & \|\lambda\|_{H^{-\frac{1}{2}}(\Omega)} \|z - w\|_{H^{\frac{1}{2}}(\Omega')} \\ \leq & \|\lambda\|_{H^{-\frac{1}{2}}(\Omega')} C_4 \|z - w\|_{H^1(\Omega)} \leq C \, C_1 \, C_2 \, C_4 \, \frac{h_1}{h_2} \, \|\lambda\|_{H^{-\frac{1}{2}}(\Omega)}^2. \end{split}$$

Therefore we have

$$\oint\limits_{\Omega^*} z\lambda ds \geq \|\lambda\|_{H^{-\frac{1}{2}}\langle\Omega\rangle}^2 \Big[C_3 - \frac{C\,C_1C_2C_4}{K}\Big]\,.$$

Taking K sufficiently large we may obtain

$$C_3 - \frac{CC_1C_3C_4}{K} \ge \frac{C_3}{2} > C_0 > 0 \tag{3.22}$$

and hence

$$\oint_{\Omega'} z \lambda ds \ge C \|\lambda\|_{H^{-\frac{1}{2}}(\Omega')}^{2}.$$
(3.23)

Let us take now

$$v = u - z$$
 and $\mu = -2\lambda$.

Obviously $(v, \mu) \in M$ and the rest of the proof is simple repetition of the remainder of the proof of the Theorem 3.1.

The convergence of the proposed method follows now almost immediately from Theorem 2.9 and the basic properties of (t, k)-regular systems

Theorem 3.4. Let $f \in H^k(\Omega)$, $k \ge 0$, $g \in H^l(\Omega^{\bullet})$, $l \ge \frac{1}{2}$ and let u_0 be the solution of the Dirichlet problem due to Theorem 2.4. Further, let $\bar{u}_0(h_1, h_2)$ and $\bar{\lambda}_0(h_1, h_2)$ be the approximate solution of the finite element method with Lagrange multiplier (with $k_1 \ge 1$, $k_2 \ge \frac{1}{2}$) and K in (3.16) be sufficiently large that Theorem 3.3 holds. Then

$$\|\bar{u}_{0}(h_{1}, h_{2}) - u_{0}\|_{H^{1}(\Omega)} + \|\bar{\lambda}_{0}(h_{1}, h_{2}) - \frac{\partial u_{0}}{\partial n}\|_{H^{-\frac{1}{2}}(\Omega^{*})}$$

$$\leq C \left[h^{\mu_{1}} \|f\|_{H^{1}(\Omega)} + h^{\mu_{1}} \|g\|_{H^{1}(\Omega^{*})}\right]$$
(3.24)

where

$$\mu_1 = \min[t_1 - 1, k + 1, t_2 + \frac{1}{2}]$$
 (3.25)

$$\mu_2 = \min[t_1 - 1, l - \frac{1}{2}, t_2 + \frac{1}{2}]$$
 (3.26)

with $h = \max(h_1, h_2)$.

Proof. Using Theorem 2.9 we have to show the existence of $\hat{\omega}=(\hat{u},\,\hat{\lambda})\in M_1$ such that

$$\|u_0 - \hat{u}\|_{H^1(\Omega)} \leq CQ$$

and

$$\left\| \frac{\partial u_0}{\partial n} - \hat{\lambda} \right\|_{H^{-\frac{1}{2}}(\Omega^*)} \le CQ$$

where

$$Q = C \left[h^{\mu_1} \| f \|_{H^{k}(\Omega)} + h^{\mu_2} \| g \|_{H^{l}(\Omega^*)} \right].$$

But this follows immediately from Theorem 2.4, Theorem 2.2 and the basic property of the regularity of the systems $\mathcal{S}_{h_1}^{t_1,k_1}(\Omega)$ and $\mathcal{S}_{h_2}^{t_2,k_2}(\Omega)$.

4. The Rate of Convergence in the Space $L_2(\Omega)$

On the previous section we established the rate of convergence in the space $H^1(\Omega)$ of the approximate solution $\bar{u}_0(h_1,h_2)$ obtained by the finite element method with Lagrange multiplier. In this section we will be interested in the error estimate in the space $L_2(\Omega)$. We shall use—with minor modification—the usual technique in obtaining estimates of this kind.

Let us denote $\varepsilon = \bar{u}_0(h_1, h_2) - u_0$ and $\eta = \bar{\lambda}_0(h_1, h_2) - \partial u_0/\partial n$. By Theorem 3.4 we know that

$$\|\varepsilon\|_{H^1(\Omega)} \leq Q$$
, $\|\eta\|_{H^{-\frac{1}{2}}(\Omega)} \leq Q$

where Q is given by (3.24). Therefore also $\|\varepsilon\|_{L_4(\Omega)} \leq Q$.

Let V be the solution of the Dirichlet problem for the differential equation $-\Delta V + V = \varepsilon$ with homogeneous boundary condition V = 0 on Ω . Using Theorem 2.4 and 2.2 we have

$$||V||_{H^{3}(\Omega)} \leq CQ, \tag{4.1}$$

$$\left\| \frac{\partial V}{\partial n} \right\|_{H^{\frac{1}{2}}(\Omega^*)} \le CQ, \tag{4.2}$$

and for every $(v, \mu) \in H$ we have

$$B\left(V, \frac{\partial V}{\partial n}; v, \mu\right) = \int\limits_{\Omega} \varepsilon v \, dx.$$

On the other hand, we have $(\varepsilon, \eta) \in H$ and therefore

$$B\left(V, \frac{\partial V}{\partial n}; \varepsilon, \eta\right) = \|\varepsilon\|_{L_{1}(\Omega)}^{2}. \tag{4.3}$$

But,

$$B(w, \varrho; \varepsilon, \eta) = 0 \tag{4.4}$$

for every $(w, \varrho) \in M$ because of the definition of $u(h_1, h_2)$ and $\lambda(h_1, h_2)$. It follows from (4.1), (4.2) and the regularity property that there exists $(z, \xi) \in M$ such that

$$\left\| \left(z - V; \frac{\partial V}{\partial n} - \xi \right) \right\|_{H} \le C h^{\mu} \| \varepsilon \|_{L_{\bullet}(\Omega)}$$
(4.5)

where

$$\mu = \min[t_1 - 1, 1, t_2 + \frac{1}{2}].$$
 (4.6)

Because of (4.4), we have using (4.3)

$$\|\varepsilon\|_{L_{\mathbf{a}}(\Omega)}^2 = B\left(V-z, \frac{\partial V}{\partial n} - \xi; \, \varepsilon, \, \eta\right) \leq \left\|\left(z-V, \frac{\partial V}{\partial n} - \xi\right)\right\|_{H} \|\varepsilon, \, \eta\|_{H}.$$

But using (4.1) and (4.2) and (4.5), we obtain

$$\|\varepsilon\|_{L_{\mathbf{a}}(\Omega)}^{2} \leq C h^{\mu} \|\varepsilon\|_{L_{\mathbf{a}}} Q. \tag{4.7}$$

Therefore

$$\|\varepsilon\|_{L_{2}(\Omega)}^{2} \le C [h^{\mu_{1}+\mu} \|f\|_{H^{k}(\Omega)} + h^{\mu_{2}+\mu} \|g\|_{H^{1}(\Omega)}]$$
(4.8)

where μ_1 is given by (3.25), μ_2 by (3.26) and μ by (4.6). We have proved the following theorem.

Theorem 4.1. Under the assumptions of the Theorem 3.4, (4.8) holds with $\varepsilon = (u - \bar{u}_1(h_1, h_2))$.

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