

LEAST SQUARES APPROXIMATIONS TO
FIRST ORDER ELLIPTIC SYSTEMS

by

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1. Notations, the Analytic Problem

In the following u, v, \dots will denote pairs $(u^1, u^2), (v^1, v^2)$ of functions defined in a bounded domain $\Omega \subseteq \mathbb{R}^2$ with boundary $\partial\Omega$ sufficiently smooth. If both components are in $L_2(\Omega)$ resp. the Sobolev-spaces $W_2^K(\Omega)$ we will write $u \in H_0$ resp. $u \in H_K$. We will also use the notation

$$(1) \quad \begin{aligned} (u, v) &= (u^1, v^1)_{L_2(\Omega)} + (u^2, v^2)_{L_2(\Omega)} \\ (u, v)_K &= (u^1, v^1)_{W_2^K(\Omega)} + (u^2, v^2)_{W_2^K(\Omega)} \end{aligned}$$

Summary: Linear boundary value problems for elliptic systems in the sense of Petrovski are considered. Using linear finite elements a least squares method is discussed. The concept of nearly zero boundary conditions - i.e. the boundary condition is imposed in the nodes on the boundary exactly - gives quasi-optimal error estimates in the L_2 - and W_2^1 -norms.

$$(2) \quad \|u\| = (u, u)^{1/2}, \quad \|u\|_K = (u, u)_K^{1/2}$$

As a model problem we will consider the elliptic system

$$(3) \quad \begin{aligned} L u &= f & \text{in } \Omega, \\ \text{i.e.} \quad L^1 u &= u_x^1 - u_y^2 + a^{11}u^1 + a^{12}u_y^2 = f^1 \\ L^2 u &= u_y^1 + u_x^2 + a^{21}u^1 + a^{22}u^2 = f^2 \end{aligned} \quad \text{in } \Omega.$$

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We remark that any elliptic system in the sense of Petrovski is - up to a coordinate-transformation-equivalent to (3), see

HAACK-HELLWIG [1]. In addition to (3) we impose the boundary condition

$$(4) \quad u^1 = u^1 \cos \sigma + u^2 \sin \sigma = 0 \quad \text{on } \partial\Omega.$$

Essential for the solvability of the boundary value problem

is the index of $\sigma = \sigma(s) - s$ arc-length on $\partial\Omega$ - defined by

$$(5) \quad \text{ind } (\sigma) = \frac{1}{2\pi} \oint_{\partial\Omega} d\sigma.$$

In case of $n = \text{ind } (\sigma) \geq 0$ then (3), (4) is always solvable, the number of linear independent solutions of the homogeneous problem ($f = 0$) is $2n + 1$. In case of $n < 0$ (3), (4) possesses a solution only if f fulfills $2|n|-1$ linear independent integral relations.

In order to characterize in case of $n \geq 0$ a special solution to following way is possible: Let $\mathfrak{B}_n(\partial\Omega)$ be a subspace of $L_2(\partial\Omega)$ of dimension $2n + 1$ similar to the space of trigonometric functions of order n . This means there is up to a factor exactly one element in $\mathfrak{B}_n(\partial\Omega)$ with zeros in $2n$ prescribed points on $\partial\Omega$, not necessarily distinct. Then there is one and only one solution of (3), (4) such that

$$(6) \quad \oint_{\partial\Omega} p_i u^2 ds = r_i \quad (i = 1, \dots, 2n+1).$$

Here

$$(6') \quad \langle p_i, u^2 \rangle = r_i \quad (i = 1, \dots, 2n+1).$$

For $n = \text{ind } (\sigma) \geq 0$ - only this case will be considered here - let H_1^σ be the set of pairs $u = (u^1, u^2)$ with $u \in H_1$ and denotes the orthogonal complement of the boundary condition (4),

the set $\{p_i\}$ forms a basis of $\mathfrak{B}_n(\partial\Omega)$ and $\{r_i\}$ are fixed real numbers.

In case of a non-negative index n there are especially $2n + 1$ solutions u_j of (3), (4) with $f = 0$ according to

$$(8) \quad \oint_{\partial\Omega} p_i u_j^2 ds = \delta_{ij}.$$

For negative indices $n = \text{ind } (\sigma)$ on the other hand there exists one and only one solution of (3) such that

$$(9) \quad u^1 \in \mathfrak{B}_2^{|n|-1}.$$

In the following we will impose the conditions (6) for non-negative indices of σ and use the weakened form (9) for negative indices.

For functions $p, q \in L_2(\partial\Omega)$ we will use

$$\langle p, q \rangle = \int_{\partial\Omega} p q ds,$$

(10)

$$|p| = \langle p, p \rangle^{1/2}.$$

The conditions (6) may be rewritten:

$$(7) \quad \langle p_i, u^2 \rangle = r_i \quad (i = 1, \dots, 2n+1).$$

fulfilling (4) and $H_1^{\sigma, r}$ the subspace fulfilling (6).

For the sake of simplicity the coefficients a_{ij} in (3) as well as σ in (4) are assumed sufficiently smooth. Without loss of generality the elements of $\mathbb{P}_n(\partial\Omega)$ can be chosen sufficiently smooth and therefore $\mathbb{P}_n(\partial\Omega)$ can be extended to a space $\mathbb{P}_n(\bar{\Omega})$ of functions defined in $\bar{\Omega}$.

For the modified boundary value problems (3), (4), (6) in case of a non-negative index resp. (3), (9) in case of a negative index shift-theorems of the type

$$(11) \quad \|u\|_{k+1} \leq c \left\{ \|f\|_k + \sum \|r_i\| \right\}$$

are valid. c depends besides of σ , a_{ik} on $\partial\Omega$ and k .

The statements of these sections are consequences of the theory on elliptic systems developed by VÉKÚA [2].

2. Finite Element Spaces

Let Γ_h be a subdivision of Ω into generalized triangles Δ , i.e. Δ is a triangle if $\bar{\Delta}$ and $\partial\Omega$ have in common at most one point and otherwise one of the sides of Δ may be curved.

We will only consider regular subdivisions: For $\kappa > 1$ fixed there are to any $\Delta \in \Gamma_h$ two circles with radii $\kappa^{-1} h$ and κh contained in Δ resp. containing Δ .

The finite element spaces $S_h = S_h(\Gamma_h)$ we will work with consist of pairs $x = (x^1, x^2)$ of functions having the properties

$$1) \quad x \in H_1, \text{ i.e. } x^1 \in W_2^1(\Omega) ,$$

ii) x^1 restricted to any $\Delta \in \Gamma_h$ is linear.

Let $\{P_\nu\}$ be the set of nodes of Γ_h on $\partial\Omega$. The subspace $S_h^\sigma \subseteq S_h$ consists of those elements with

$$(12) \quad x^1(x)|_{P_\nu} = 0 \quad .$$

Finally $S_h^{\sigma, r} \subseteq S_h^\sigma$ consists of elements with

$$(13) \quad < p_i, x^2(x) > = r_i \quad (i = 1, \dots, 2n+1) .$$

Then the following approximation property holds:

Lemma: Let $n = \text{ind } (\sigma) \geq 0$. There is a linear projection-operator $Q_h = Q_h^{\sigma, r} : H_1^{\sigma, r} \rightarrow S_h^{\sigma, r}$

with

$$(14) \quad \|u - Q_h u\|_1 \leq c h \|u\|_2$$

$$\text{or} \quad u \in H_1^{\sigma, r} = H_1^{\sigma, r} \cap H_2$$

First let I_h denote the linear interpolation. For $u \in H_2^\sigma$
we have obviously $I_h u \in S_h^\sigma$ and

$$(15) \quad \|u - I_h u\|_1 \leq c h^{2-1} \|u\|_2 \quad (1=0,1)$$

Further we get

$$(16) \quad |r_i - \langle p_i, 1^2(I_h u) \rangle| \leq |p_i| \|1^2(u - I_h u)\|$$

and because of

$$(17) \quad |z|^2 \leq \|z\|_{L_2(\Omega)} \|z\|_{W_2^1(\Omega)}$$

for any $z \in W_2^1(\Omega)$ therefore

$$(18) \quad |r_i - \langle p_i, 1^2(I_h u) \rangle| \leq c h^{2/2} \|u\|_2$$

In order to get Q_h we add a proper combination of the interpolated homogeneous solutions of (8)

$$(19) \quad Q_h u = I_h u + \sum_{i=1}^{2n+1} \alpha_i I_h u_i$$

The conditions on $\{\alpha_i\}$ are

$$(20) \quad \sum_{i=1}^{2n+1} \langle p_i, 1^2(I_h u_i) \rangle > \alpha_i = \langle p_j, 1^2(u - I_h u) \rangle \quad (j = 1, \dots, 2n+1)$$

Because of

$$(21) \quad |1^2(u_i - I_h u_i)| \leq c h^{3/2}$$

and (8) the inverse of the matrix of the linear equations (20)
is bounded away from zero for h small enough. Using the bound

(18) for the right hand side of (20) we get

$$(22) \quad |\alpha_i| \leq c h^{3/2} / \|u\|_2$$

Therefore

$$(23) \quad \|u - Q_h u\|_1 \leq \|u - I_h u\|_1 + \max \{ \|I_h u_i\|_1 \} \sum_{i=1}^{2n+1} |\alpha_i|$$

which gives (14).

3. Least Squares Method, Error Estimates

For $u \in H_1$ and $v \in H_1$ resp. $w \in H_0$ let us define the bilinear functionals

$$a(u, v) = (\mathbb{L}u, \mathbb{L}v)$$

$$(24) \quad = (\mathbb{L}^1 u, \mathbb{L}^1 v) + (\mathbb{L}^2 u, \mathbb{L}^2 v)$$

resp.

$$(25) \quad b(u, w) = (\mathbb{L}u, w)$$

$$= (\mathbb{L}^1 u, w^1) + (\mathbb{L}^2 u, w^2)$$

Obviously we have

$$(26) \quad a(u, v) = b(u, \mathbb{L}v)$$

We get by partial integration - for $w \in H_1$ -

$$b(u, w) = (u, *_{\mathbb{L}w}) +$$

$$+ \int l^1(u) \{ w^1 \sin(\sigma+\gamma) - w^2 \cos(\sigma+\gamma) \} ds$$

$$+ \int l^2(u) \{ w^1 \cos(\sigma+\gamma) + w^2 \sin(\sigma+\gamma) \} ds$$

with γ denoting the angle between the tangent at a point of $\partial\Omega$ and the x-axis, and

$$(33)$$

$$*_{\mathbb{L}w} = g$$

$$*_{\mathbb{L}}^1 w = -w_x^1 - w_y^2 + a^{11} w^1 + a^{21} w^2,$$

$$(28) \quad *_{\mathbb{L}}^2 w = w_x^1 - w_y^2 + a^{12} w^1 + a^{22} w^2$$

being the adjoint of the differential operator \mathbb{L} .

If

$$(29) \quad n = \text{ind } (\sigma) \geq 0$$

then the index of the boundary condition

$$(30) \quad w^1 \cos(\sigma+\gamma) + w^2 \sin(\sigma+\gamma) = 0$$

with respect to the operator $*_{\mathbb{L}} - w^1$ and w^2 have to be interchanged in order to give the Cauchy-Riemann principle part - 1s

$$(31) \quad n* = \text{ind } (\sigma*)$$

$$= \text{ind } (\frac{\pi}{2} - \sigma - \gamma) = -n - 1.$$

According to (9) then (30) has to be modified

$$(32) \quad w^1 \cos(\sigma+\gamma) + w^2 \sin(\sigma+\gamma) \in P_{2n+1} \quad \text{on } \partial\Omega$$

in order that

together with (32) has a unique solution.

For simplicity we will consider the boundary value problem (3), (4), (6) only with $r_i = 0$. Then we have the duality relation

$$(34) \quad (u, \varepsilon) = b(u, w)$$

Further let v be defined by

$$(35) \quad \begin{aligned} L v &= w && \text{in } \Omega, \\ l^1(v) &= 0 && \text{on } \partial\Omega, \\ < p, l^2(v) > &= 0 && \text{for } p \in P_{2n+1}. \end{aligned}$$

Then we have

$$(36) \quad (u, \varepsilon) = a(u, v)$$

Using the shift theorem (11) we get $v \in H_{k+2}$ for $\varepsilon \in H_k$.

In order to approximate the solution u of (3), (4), (6) - with $r_i = 0$ - we use the least squares method: The approximation $u_h \in S_h^{\sigma, 0}$ is defined by

$$(37) \quad a(u_h, \chi) = (f, L\chi) \quad \text{for } \chi \in S_h^{\sigma, 0}$$

Though $a(\cdot, \cdot)$ is positive definite in $H_1^{\sigma, 0}$, it might only be semi-definite in $S_h^{\sigma, 0}$. With $e = u - u_h$ and - using an appropriate approximation U_h on u in $S_h^{\sigma, 0} = e = u - U_h$ and therefore $e = \varepsilon + \tilde{e}$ with $\tilde{e} = U_h - u_h \in S_h^{\sigma, 0}$ we get

$$(38) \quad a(e, \chi) = 0 \quad \text{for } \chi \in S_h^{\sigma, 0}$$

resp.

$$(39) \quad a(\tilde{e}, \chi) = -a(e, \chi) \quad \text{for } \chi \in S_h^{\sigma, 0}$$

By

$$(40) \quad \| \cdot \|^\ast = a(\cdot, \cdot)^{1/2}$$

(35) a semi-norm is defined. Obviously we have for $v \in H_1$

$$(41) \quad \| v \|^\ast \leq c \| v \|_1$$

So we get from (39)

$$(42) \quad \begin{aligned} \| \tilde{e} \|^\ast &\leq \| e \|^\ast \\ &\leq c \| e \|_1 \leq c \| u \|_2 \end{aligned}$$

and consequently

$$(43) \quad \| e \|^\ast \leq 2c \| u \|_2$$

Next we identify $\varepsilon = e$ in (34) and let w resp. v be the solutions of (32), (33) resp. (35). Then we get

Proposition 2: For any $x \in S_h^{\sigma,0}$ nearly zero boundary conditions of the type

$$(44) \quad \begin{aligned} \|e\|^2 &= (e, *_{LW}) \\ &= (\mathbb{I}e, v) - \int \{l_1^1(e) *_{12}(v) - l_2^2(e) *_{11}(v)\} ds \end{aligned}$$

Because of (13) and (32) the last term on the right hand side vanishes. Therefore we get

$$(45) \quad \|e\|^2 = (\mathbb{I}e, Lv) - \int l_1^1(e) *_{12}(v) ds .$$

We will estimate the two terms separately. Using the shift-theorem (11) we find $v \in H_2$ and $\|v\|_2 \leq c\|e\|$. With an appropriate approximation $x \in S_h^{\sigma,0}$ on V we get with (38), (43)

$$(Le, Lv) = a(e, v-x)$$

$$(46) \quad \begin{aligned} &\leq \|e\|' c \|v-x\|_2 \\ &\leq c h^2 \|u\|_2 \|e\| . \end{aligned}$$

In order to find a bound for the second term in (45) we first notice - using (17) -

$$(47) \quad |*_{12}(v)| \leq c\|e\| .$$

Next we make use of conditions (12) which play the role of 'nearly zero boundary conditions' introduced in [2]. With arguments parallel to those we get

Proposition 1: To any $v \in H_1 \cap H_2$ there is a $x \in S_h^{\sigma,0}$ according to

$$(48) \quad |l_1^1(v-x)| \leq c h^2 \|v\|_2 .$$

Proposition 2: For any $x \in S_h^{\sigma,0}$

conditions of the type

$$(49) \quad |l_1^1(x)| \leq c h^{3/2} \|x\|_1$$

hold.

For κ -regular triangulations inverse properties

$$(50) \quad \|x\|_1 \leq c \kappa^{-1} \|x\|$$

hold for $x \in S_h$. Therefore (49) can be replaced by

$$(51) \quad |l_1^1(x)| \leq c h^{1/2} \|x\| .$$

Now let U_h be an approximation on U according to Proposition 1 and put $\phi = U_h - u_h \in S_h^{\sigma,0}$. Then we get

$$\begin{aligned} |l_1^1(e)| &\leq |l_1^1(u-U_h)| + |l_1^1(\phi)| \\ &\leq c h^2 \|u\|_2 + c h^{1/2} \|\phi\| \\ (52) \quad &\leq c h^2 \|u\|_2 + c h^{1/2} \{ \|u-U_h\| + \|u-u_h\| \} \\ &\leq c h^2 \|u\|_2 + c h^{1/2} \|e\| . \end{aligned}$$

The bounds (46) and (47), (52) give - see (45) -

$$(53) \quad \|e\|^2 \leq c n^2 \|u\|_2 \|e\| + c n^{1/2} \|e\|^2$$

and so for n small enough

$$(54) \quad \|e\| \leq c n^2 \|u\|_2$$

Since we now know e to be bounded we have as a consequence the unique solvability of the defining equations (37).

Using (50) we also get the error estimate

$$(55) \quad \|e\|_1 \leq c n \|u\|_2$$

Literature

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