THE $p$ AND $h$-$p$ VERSIONS OF THE FINITE ELEMENT METHOD:
BASIC PRINCIPLES AND PROPERTIES

by

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In this paper, the fundamental theoretical ideas behind the relatively recent p-version and h-p version are discussed. In the p-version, a fixed mesh is used and p is allowed to increase. The h-p version combines both approaches. The basic properties and characteristics of these newer versions, especially in areas where their behavior is significantly different from that of the h-version, are described and explained.
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BASIC PRINCIPLES AND PROPERTIES

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Abstract

In the classical form of the finite element method, called the \textit{h} version, piecewise polynomials of fixed degree \( p \) are used and the mesh size \( h \) is decreased for accuracy. In this paper, we discuss the fundamental theoretical ideas behind the relatively recent \textit{p version} and \textit{h-p version}. In the \textit{p} version, a fixed mesh is used and \( p \) is allowed to increase. The \textit{h-p} version combines both approaches. We describe and explain the basic properties and characteristics of these newer versions, especially in areas where their behavior is significantly different from that of the \textit{h} version. We include simplified proofs of key concepts and provide computational illustrations of several results.

Subject Classifications: AMS (MOS): 65N30
1. INTRODUCTION

A major tool in computational mechanics today is the finite element method (FEM). Although this method has a long history (see e.g., [Od],[Wl]) the beginning of its computational success can be related to [TC],[Cl]. Originally, the main application of the FEM was in the area of structural mechanics but today, its use has expanded to various other fields such as fluid mechanics, thermal analysis, electrical engineering, etc., with applications to both linear and nonlinear, stationary and transient analysis. Consequently, there is a vast body of published research in the field of finite element theory and applications. In Figure 1.1, we have presented data showing the number of papers related to the FEM. (See also [Ma].)

![Graph showing the number of FEM papers over years](image)

Figure 1.1. Number of FEM papers in MAKEBASE database (1990, 1991 are not completely correct—delayed literature acquisition. Computational fluid mechanics is not included.) (Courtesy: Jaroslav Mackerle, Linkoping Inst. of Technology, Sweden).
Most of the papers represented in Figure 1.1 are related to the classical form of the FEM, called the \textit{h version}, in which polynomials of fixed degree \( p \) are used and the mesh is refined to increase accuracy. There are several commercial finite element programs based on the \( h \) version of the FEM which are currently available (e.g., MSC/NASTRAN, Cosmos/M, Abaqus, Aska, Adina, Ansys, etc.).

In contrast, only a few (less than 100) of the above papers are related to the relatively recent \( p \) and \( h-p \) versions. The \( p \) version of the FEM uses a fixed mesh but increases the polynomial degree \( p \) to increase accuracy. In the \( h-p \) version, the two approaches of mesh refinement and degree enhancement are combined. (The \( h \) and \( p \) versions can be considered as special cases of the \( h-p \) version.) The first theoretical papers on these new versions ([BSKI], [BD]) were published in 1981. Recently, a few commercial and large research programs have become available on these versions, such as MSC/PROBE, Applied Structure, PHLEX and STRIPE. In contrast to the many books devoted to the \( h \) version FEM in engineering and mathematics, the only book addressing the \( p \) and \( h-p \) versions is [SzB].

In this paper, we discuss the basic theoretical ideas of the \( p \) and \( h-p \) versions of the FEM. Rather than an exhaustive summary of published work, our goal here is to describe only the most salient (and practically useful) features which are characteristic of the \( p \) and \( h-p \) versions. We give numerous elaborations and examples of key theoretical results to help explain the underlying reasons for the behavior of these versions, and to bring out the differences that exist in comparison to the classical \( h \) version. We also include informal proofs of several important theorems to make them more accessible and to help introduce the reader to \( p \) version theoretical
techniques of approximation (which can be quite different from standard h version techniques). Throughout this survey, we concentrate on problems related to structural mechanics, as opposed to other applications.

The plan of our paper is as follows. In Section 2, we consider the one-dimensional case, which allows us to present various ideas and results (which hold for higher dimensions as well) in a simplified setting. Perhaps the most important characteristic of the p and h-p versions, the enhanced asymptotic rates of convergence that are possible, is discussed at length in Sections 2.3 - 2.4. We show, by elementary arguments, exactly what leads to the well-documented phenomenon of the "doubling" in the rate of convergence of the p version (over the h version) in the presence of singularities. We also demonstrate (by means of some simple proofs) the following basic difference: the rate for the h version can never be better than algebraic (no matter which mesh is used), while the h-p version can yield exponential rates. In Sections 2.5, 2.6, we discuss, respectively, the error in the \( L_2 \) norm and the problem of "pollution", two areas in which the p version behaves differently from the h version.

The two- and three-dimensional analogs of results in Section 2 are discussed in Section 4. The refined approximation theorems for the p and h-p versions presented in Section 4 are extremely dependent on precise regularity results for the solution (which are more detailed than those required for the h version). Therefore, Section 3 is devoted to a survey of various types of regularity theorems for elliptic problems that are used in the development of such convergence results. In particular, regularity results in terms of countably normed spaces, which are crucial in deriving exponential rates of convergence for the h-p version, are presented in Section 3.3.
An inherent property of the p version is its robustness in most situations, a characteristic that can be put to good use in applications to problems that are susceptible to numerical locking. We explain this phenomenon in Section 5 and survey some recent related results, using the equations of linear elasticity as a model.

The convergence results in Section 2, 4 are mostly in the energy norm. Often, instead of the energy, some other functional of the solution is of interest. In Section 6, we discuss extraction procedures by which these functionals can be recovered from the solution. As a result, the high rates of convergence for the h-p version in the energy are preserved for these quantities of interest as well.

In Section 7, we consider the use of the p version in plate modeling. We show how a hierarchy of plate models may be set up, which approximates the three-dimensional plate with increasing accuracy. The implementation of such a hierarchy of plate models can then be effected in a natural way by using the p version. This may be used to ensure convergence to the solution of the three-dimensional plate, as opposed to that of a two-dimensional plate model which is what happens when a fixed model (e.g., the Reissner-Mindlin plate) is discretized.

Implementational aspects of the p and h-p version are briefly discussed in Section 8.

We have not tried to compose an exhaustive bibliography of publications on the p and h-p version (as could, for instance, be obtained using MAKEBASE). Rather, we have provided only selected representative references, for the interested reader to obtain more information on the topics discussed. In this connection, some representative references on the p and h-p versions
for areas not discussed in this paper are as follows: boundary element methods [PS], [SS], mixed methods [JV], [S1], parabolic problems [BJ], free boundary problems [ZB], shape optimization [Sh], non-linear problems [Ng], solution of large systems [Mn]. In addition, see [Sz] for a survey on computational questions related to the p and h-p version.

2. The h, p AND h-p VERSIONS IN ONE DIMENSION

The one-dimensional setting allows a simple and convenient vehicle for the exposition of our basic ideas and results. The fundamental properties of and differences between the h, p and h-p versions are most clearly explained and understood through the one-dimensional results and numerical illustrations in this section. These also indicate what can be expected in two and three dimensions (discussed in succeeding sections), for which strongly analogous results hold. A detailed analysis of the 1-d case has been made in [GiB], from which most of the results in this section have been cited.

2.1. The model problem and its discretization

Let us consider the model problem

\[(2.1) \quad -u'' = f \quad \text{on} \quad \Omega = (0,1) = I, \quad u(0) = u(1) = 0.\]

We will consider the general case when it is only known that \( f \in H^{k-2}(\Omega) \) \(^1\), \( k \geq 1 \), as well as the particular case when \( f \) is such that the exact solution of (2.1) is

\[1\) Here and elsewhere, \( W^{k,t}(\Omega), H^k(\Omega), H^k_0(\Omega) \) will denote the usual Sobolev spaces.
It will be seen shortly that much more refined convergence results exist for the particular case (2.2) than the general case. The form (2.2) is obviously singular at the endpoint \( x = 0 \). The reason we choose it is that it serves as a model for two (and higher) dimensions where analogous singularities exist at the corners of domains which are curvilinear polygons (see Section 3). The restriction \( \alpha > \frac{1}{2} \) is made to ensure that \( u \in H^1(\Omega) \).

In general, the regularity of the solution \( u \) is determined by the regularity of \( f \). Obviously, if \( f \in H^{k-2}(\Omega) \), then \( u \in H^k(\Omega) \). An analogous but more complicated relation between the regularity of the input data and of the solution occurs in two and three dimensions and will be discussed in Section 3.

As a prelude to using the finite element method, the problem (2.1) is cast into the following variational form. Find \( u \in H^1_0(\Omega) \) such that for any \( v \in H^1_0(\Omega) \),

\[
(2.3) \quad B(u,v) = \int_\Omega u'(x)v'(x)dx = f(v)
\]

with

\[
 f(v) = \int_\Omega f(x)v(x)dx.
\]

The above may now be discretized by constructing a sequence of finite-dimensional (finite element) subspaces \( V^n \subset H^1_0(\Omega) \) of dimension \( N_n = N(V^n) \), \( n = 1,2,\ldots \), called the number of degrees of freedom. These spaces consist of piecewise polynomials on a sequence of meshes \( \mathcal{G}_n \), \( n = 1,2,\ldots \), on \( \Omega \) where

\[
\mathcal{G}_n = \{0 = x_0^n < x_1^n < \cdots < x_m^n = 1\}.
\]
We denote $i^n_j = [x^n_j, x^n_{j-1}]$, $h_{n,j} = x^n_j - x^n_{j-1}$, $j = 1, \ldots, m$, $h_n = \max_{1 \leq j \leq m} h_{n,j}$, and call $i^n_j$ elements.\(^2\)

Let us now consider three basic types of meshes, which are analogous to those used in the two and three dimensional settings.

1. **The quasiuniform mesh.** Here $\{i^n\}$ is such that we have for all $n$
\[
\frac{h_n}{h_{n,j}} \leq \tau,
\]
where $\tau > 0$ will be called the quasiuniformity constant. A special case is the uniform mesh, where $\tau = 1$ and
\[
(2.5) \quad x^n_i = \frac{i}{m}, \quad i = 0, \ldots, m.
\]

2. **The radical mesh** with power $\beta > 0$ (with respect to $x = 0$). This is defined by choosing
\[
(2.6) \quad x^n_i = \left(\frac{i}{m}\right)^\beta, \quad i = 1, \ldots, m.
\]
(the uniform mesh is obviously a radical mesh with $\beta = 1$).

The radical mesh is a special case of a "graded mesh" with "grading function" $g(y) = y^\beta$, defined by
\[
(2.7) \quad x^n_i = g\left(\frac{i}{m}\right), \quad i = 1, \ldots, m.
\]

**Remark 2.1:** A generalization of graded meshes uses the notion of density $\psi (\psi = g'$ in the one-dimensional case). It is naturally applicable to two dimensions. See [Hu].

\(^2\)We consider here only a denumerable family of meshes. Of course, we could consider a one (or more) parametric family of meshes as well.
3. The geometric mesh with ratio $q$, $0 < q < 1$ (with respect to $x = 0$). Here,

\begin{align}
(2.7a) & \quad x^n_i = q^{n-1}, \quad i = 1, \ldots, m_n \\
(2.7b) & \quad x^n_0 = 0.
\end{align}

A more general definition will be presented in Section 4 for the 2-d case.

Let $p_n = (p_{n,1}, p_{n,2}, \ldots, p_{n,n})$ be an integer, be called the degree vector. By $V^n(F^n, P_n) \subset H^1_0(\Omega)$, we denote the set of all functions $v$ such that the restriction $v\big|_{I^n_j}$ is a polynomial of degree $p_{n,j}$, i.e. it lies in $P_{n,j}$. If $p_{n,j} = p_n$, then we speak about the uniform degree vector $P_n$, for which we write $p_n$. $V^n = V^n(F^n, P_n)$ will be called the space of elements of degree $p_n$.

Having constructed a finite-dimensional subspace as above, we now define the finite element solution as $u_n \in V^n$ satisfying (2.3) for all $v \in V^n$.

Then we immediately obtain, with $e_n = u - u_n$,

\begin{equation}
(2.8) \quad \|e_n\|_E \leq \inf_{v \in V^n} \|u - v\|_E,
\end{equation}

where $\|u\|_E = (B(u, u))^{1/2}$ is the energy norm, which is equivalent to the $H^1(\Omega)$ norm.

In this paper, we will be most interested in the energy norm measure of the error in the finite element solution. Any change of the norm can (and usually does) change the results and conclusions following from them. In practice, many different accuracy assessments are used.
2.2 Extension processes

In order to increase the accuracy of the finite element solution, we define an extension process which entails the construction of a sequence of subspaces \( V^n \) (of particular character). The classical extension process uses \( V^n(J^n,p_n) \), where \( J^n \), \( n = 1, 2, \ldots \), is a sequence of meshes with \( h_n \) decreasing (\( h_n \to 0 \) as \( n \to \infty \)) and \( p_n = p \), usually \( p = 1, 2 \). Here the meshes \( J^n \) could be a sequence of uniform, quasiuniform or non-quasiuniform meshes. This choice of \( V^n \) is called the \( h \) version of the FEM.

The \( p \) version of the FEM on the other hand uses \( V^n(J^n,p_n) \), where \( J^n = J \) is a fixed mesh and the degrees \( p_{n,j} \to \infty \) as \( n \to \infty \). The increase of \( p_{n,j} \) can be uniform (\( p_n = p \)), or selective (i.e., \( p_{n,j} \) are different for different \( j \)). The mesh \( J \) can be essentially uniform or it may be strongly refined in certain areas.

Finally, the \( h-p \) version of the FEM consists of \( V^n(J^n,p_n) \) with \( J^n \) in general being refined and \( p_n \) increasing with \( n \). Obviously, the \( h-p \) version is a generalization of the \( h \) and \( p \) versions.

It is useful to distinguish between hierarchic and non-hierarchic extensions. Hierarchic ones are such that \( V^{n+1} \supset V^n \). This assures that the error decreases monotonically, i.e., \( \| \varepsilon_{n+1} \|_E \leq \| \varepsilon_n \|_E \). If the spaces \( V^n \) are not hierarchic then, in general, this monotone behavior does not hold. The basis functions for the \( p \) version are in general hierarchic.

In practice, of course, we always use finite meshes with elements of finite degree so that for concrete computations, the distinction between different extensions is blurred. Nevertheless, this distinction is important not only theoretically, but also for practical purposes, since while performing an actual computation it is usually assumed that the results are in
the asymptotic range of the extension process. Since, in practice, one or more steps of the extension process are usually carried out, it is here that the behavior and asymptotic properties of various extensions play an important role.

2.3. Approximation theorems for quasiuniform meshes

The rate at which the finite element approximations converge to the true solution will depend upon two factors: the definition of each of the subspaces \( V^n \) (i.e., quasiuniform or non-quasiuniform meshes, choice of \( P_n \), etc.) as well as the type of extension procedure (h, p or h-p). We now present various approximation theorems which bound the right hand side of (2.8), in terms of the number of degrees of freedom \( N_n = \dim V^n \). If various methods are to be compared, they should be assessed in terms of the total work required to obtain a given accuracy. Although \( N_n \) is one measure of the work required, it does not give the complete picture, since various other factors like amount of user time, special solution techniques, etc., are involved (see [BE], [BEM], [CM]). With these qualifications in mind, we first discuss the case of the h-p version using quasiuniform meshes. We will omit the subscript \( n \) where no confusion would occur.

**Theorem 2.1.** Let \( p_n = p_\ast = p \), and \( J^n \) be quasiuniform. Then for \( k \geq 1 \)

\[
\| e \|_{E} \leq C h^{(\mu - 1)p - (k - 1)} \| u \|_{H^k(\Omega)}
\]

where \( \mu = \min(k, p+1) \) and \( C \) is independent of \( n \) (i.e., of \( h, p \)) and \( u \), but depends on \( k \) and the quasiuniformity constant \( \tau \).

From this, noting that \( N = N(V^n) \approx C h^{-1} p \), we obtain the following corollaries.
Corollary 2.1. For the h version (i.e., $p_n = \bar{p}$) we get

\[ (2.10) \quad \|e\|_E \leq C h^{(\mu-1)} \|u\|_{H^k(\Omega)} \leq C N^{-(\mu-1)}. \]

Corollary 2.2. For the p version (i.e., $q_n = \bar{q}$) we get

\[ (2.11) \quad \|e\|_E \leq C p^{-(k-1)} \|u\|_{H^k(\Omega)} \leq C N^{-(k-1)}. \]

The estimate (2.9) is a refined version of the estimate (2.10) which is the one found in many books on the FEM. It shows explicitly how the constant $C$ in (2.10) behaves with respect to $p$. The two-dimensional version of this estimate (see also Theorem 4.1) was proven in [BS1] while that of the estimate (2.11) was proven in [BS2].

Notice that (2.9) shows that if the h version is used in a case that $k < p + 1$, increasing $p$ further will decrease the error by decreasing the "constant" $C p^{-(k-1)}$, though it will not affect the exponent of $h$. This exponent of $N$ (or $h$) is often called the asymptotic order of convergence, and, for example, for (2.10) will often be written as

\[ \|e\|_E = O(N^{-(\mu-1)}). \]

The above estimates are optimal in the following sense. For any $N$, there is a function $u \in H^k(\Omega)$, possibly depending on $N$, such that

\[ (2.12) \quad \|e\|_E \geq C N^{-(k-1)} \|u\|_{H^k(\Omega)} \]

with $C > 0$ independent of $u$, $N$, the mesh $\mathcal{J}$, and element degrees $p$. (2.12) follows from the theory of $n$-widths (see [P1]).

If $k \geq p + 1$, then for the h version there is a $u \in H^k(\Omega)$ such that
Let us remark that the error bounds (2.12) and (2.13) must be distinguished from asymptotical expansions of the type

\[ \|e\|_E \geq C(p)h^p \|u\|_{H^k(\Omega)} \geq C N^{-p} \]

with \( C \) independent of \( N \).

In practice, we are usually interested in the convergence of the FEM for a specific \( u \), namely the solution of the problem of interest. For this specific \( u \), the rate of convergence could be higher than the rate for the "worst" case (which is the one addressed in (2.9) and (2.11)). For example, function \( u \) defined in (2.2) belongs maximally to \( H^{\alpha-1/2-\epsilon}(I) \) but the convergence rate in (2.9) can be too pessimistic since for particular relations between \( h \) and \( p \), better asymptotic rates of the form (2.14) hold (see e.g. Theorem 2.2). In general, we must therefore distinguish between the convergence rate for the worst case (which may sometimes only be achievable for a sequence of exact solutions dependent on \( N \)) and the convergence rate for a particular case, especially when the solution under consideration is "reasonable". Consequently, when the accuracy of the FEM is experimentally studied by numerical determination of the exponent \( k-1 \) of \( N \) in (2.14), it should be kept in mind that this only gives an upper bound for the true exponent in (2.12), (2.13) and cannot always serve as an indicator of the performance of the method for all cases.

The estimate (2.13) shows that the maximum possible rate of convergence using the \( h \) version is always bounded by the fixed degree \( p \) of polynomials.
used. In contrast, (2.11) shows that the p version gives arbitrarily high rates of convergence, depending only on the smoothness of the function \( u \).

The above estimates are based only on the information that \( u \in H^k(\Omega) \). If more information is available, then better estimates can be obtained. For example, if \( u \) is analytic on \([0,1] \), then the error for the p version decays exponentially.

Let us now assume that the solution has a singularity at a point, or specifically, that \( u \) is given by (2.2). Then we have

**Theorem 2.2.** Let the solution \( u \) be given by (2.2). Then

\[
\|e\|_E \leq C \min\{h^{\alpha-1/2}, h^{\min(\alpha-1/2,p-\alpha+1/2)} p^{-2(\alpha-1/2)}\}.
\]

This gives immediately

**Corollary 2.3.** For the h version

\[
(2.15) \quad \|e\|_E \leq Ch^{\alpha-1/2} = O(N^{-(\alpha-1/2)}).
\]

**Corollary 2.4.** For the p version

\[
(2.16) \quad \|e\|_E \leq Cp^{-2(\alpha-1/2)} = O(N^{-2(\alpha-1/2)}).
\]

The estimate (2.15) follows from (2.10) when Besov spaces \( B^\alpha_{2,\infty} \) are used instead of fractional Sobolev spaces \( H^k(\Omega) \). It is sharp when quasiuniform meshes are used, in the sense that there exists \( \tilde{C} > 0 \) such that

\[
(2.17) \quad \|e\|_E \geq \tilde{C}h^{\alpha-1/2}.
\]

Also, (2.16) is sharp as well. Theorem 2.2 is once more a special version of the theorem for the two-dimensional case proven in [BS1], [BD] (where \( \alpha - \frac{1}{2} \))
has to be replaced by $\alpha$).

We see that the $p$ version gives twice the rate of convergence as the $h$ version (on quasiuniform meshes) when the solution has an $x^\alpha$ type of singularity at an endpoint of the element. If the singularity is located inside the element, then this doubling in the rate does not occur. The above result about doubling of the rate was first proven for the two-dimensional case in [BSK], where a slightly weaker result, showing the doubled rate of convergence reduced by an amount $O(N^e)$ ($e > 0$ arbitrary), was obtained. The dependence on $e$ was removed in [BS2].

Let us outline the reason for obtaining the rate of convergence $O(N^{-2(\alpha-1/2)})$ in Corollary 2.4. Suppose, for simplicity, we have just one element and the domain is $\Omega = (-1,1)$ instead of $\Omega = (0,1)$, with the solution (2.2) correspondingly scaled. Then if we were to use the standard estimate (2.11) with $u$ being imbedded in the Besov space like we did for the $h$ version, we would get an identical rate $O(N^{-(\alpha-1/2)})$ as in (2.15). It turns out, however, that using various orthogonality properties of Jacobi polynomials, one can show that [GiB]

\begin{equation}
\|u - u_p\|_E \leq C_p^{-k}\|u'\|_{V_k^k(\Omega)}
\end{equation}

where

\begin{equation}
\|v\|^2_{V_k^k(\Omega)} = \int_{-1}^{1} (1-t^2)^k (v(t))^2 \, dt
\end{equation}

so that the usual Sobolev norm in (2.11) is replaced by the weighted Sobolev norm in (2.19). Now it may be readily verified that although the solution (2.2) satisfies $u' \in H^k(\Omega)$ for $k < \alpha - \frac{1}{2}$ only, due to the weight $(1-t^2)^k$
in (2.19), we have instead, \( u' \in V_k^k(\Omega) \) for \( k < 2(\alpha - \frac{1}{2}) \), leading to the doubled rate up to an arbitrary \( \varepsilon > 0 \).

In order to obtain a rate of \( O(p^{-k}) \) with \( k \) actually equal to \( 2(\alpha - \frac{1}{2}) \) without the \( \varepsilon \), a more detailed analysis is required. Specifically, we expand \( u' \) in terms of the Legendre polynomials \( \{P_n\} \)

\[
(2.20) \quad u'(x) = \sum_{n=1}^{\infty} a_n P_n(x)
\]

and choose \( u_p \) such that \( u_p(\pm 1) = 0 \) and \( u'_p \) is the sum of the first \( (p-1) \) terms of (2.20). Then we have, by the orthogonality property of \( \{P_n\} \),

\[
(2.21) \quad \|u - u_p\|_E^2 = \sum_{n=1}^{\infty} \frac{a_n^2}{2n+1}.
\]

Here, \( a_n \) may be explicitly computed by the formula

\[
a_n = \frac{2n+1}{2} \int_{-1}^{1} u' P_n(t) dt = \frac{2n+1}{2} \int_{-1}^{1} (\alpha(1+t)^{\alpha-1} - 2^{\alpha-2}) \frac{1}{2n!(1-t^2)^n} dt
\]

using the Rodriguez form of Legendre polynomials. Once again, the weight \( (1-t^2)^n \) leads to an "absorption" of the singularity at the end of \( \Omega \).

Integrating by parts and using some properties of the gamma function leads to the following estimate, proved in [GiB]

\[
(2.22) \quad a_n = (-1)^{n-1} \frac{C(\alpha)}{n^{2\alpha-1}} \left( 1 + \frac{C_1(\alpha)}{n} + \frac{C_2(\alpha)}{n^2} + \cdots \right)
\]

with

\[
C(\alpha) = \frac{2^{\alpha}|\Gamma(\alpha)|^2}{\pi} \text{sin}(\alpha-1)
\]

(2.21), (2.22) lead to the following theorem, which has been stated in terms
of the original domain \( \Omega = (0,1) \) ([GIB]):

**Theorem 2.3.** For the \( p \) version in the case of a single element, we have

\[
\|e\|_E = C_0(\alpha) \frac{1}{p^{2\alpha-1}} (1 + O(\frac{1}{p}))
\]

as \( p \to \infty \), where

\[
C_0(\alpha) = \frac{\alpha(\Gamma(\alpha))^2 |\sin \pi \alpha|}{\pi \sqrt{2\alpha-1}}.
\]

Let us show a numerical example. Consider the cases \( \alpha = 0.7 \) and \( \alpha = 3.5 \). In Table 2.1 we show computed values of \( \|e\|_E = E_p \) and \( R_A^A = \frac{E_A}{E_p} \) where \( E_p \) is the leading term in the right hand side of (2.23). The numerical results are seen to be in close agreement with (2.23).

**TABLE 2.1. The error of the finite element solution.**

<table>
<thead>
<tr>
<th>( p )</th>
<th>( \alpha = 0.7 )</th>
<th>( \alpha = 3.5 )</th>
</tr>
</thead>
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<td></td>
<td>( E_p )</td>
<td>( R_A^A )</td>
</tr>
<tr>
<td>1</td>
<td>4.743E-1</td>
<td>0.9877</td>
</tr>
<tr>
<td>2</td>
<td>3.637E-1</td>
<td>0.0067</td>
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<tr>
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<td>1.626E-1</td>
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</table>
2.4. Approximation theorems for non-quasiuniform meshes

Let us now consider the case when non-quasiuniform meshes are used and the solution is given by (2.2). We will show that it is possible to get an exponential rate of convergence (in $N$) using the $h$-$p$ version of the FEM. First, however, let us discuss the $h$ version. Let us show that the best possible rate is $O(N^{-p})$ (no matter what mesh is used) i.e., the convergence is always algebraic. As before, $N$ is the number of degrees of freedom and $p$ is the degree of polynomials used in the $h$ version.

For simplicity, we take $p = 1$. Let $w(x) = u'(x) = ax^{\alpha - 1} - 1$ on $\Omega = (0,1)$. Then $w(x)$ is monotonic and $|w'(x)| \geq a|\alpha - 1| = \kappa$. Now over $I_j = (x_{j-1}, x_j)$ we have

$$\int_{I_j} |e'|^2 dx \geq \inf_{c \in \mathbb{R}} \int_{I_j} |w(x) - c|^2 dx = \frac{2}{3} \left( \frac{h_j}{2} \right)^2 \kappa^2.$$ 

Hence

$$\|e\|_E^2 \geq C \sum_{j=1}^{m_n} h_j^3.$$ 

(Here $m_n$ is the number of elements). Since the $h_j$'s in the above are arbitrary, we estimate the above sum as follows. Let

$$J = \{ j | h_j \geq \frac{1}{2m_n} \}, \quad J^* = \{ j | h_j < \frac{1}{2m_n} \}.$$ 

Then obviously

$$\sum_{j \in J^*} h_j < m_n \frac{1}{2m_n} \leq \frac{1}{2}$$

so that since $I = (0,1)$ we have

$$\sum_{j \in J} h_j \geq \frac{1}{2}.$$
Therefore
\[
\sum_{j=1}^{m} h_{j}^3 \geq \sum_{j \in J} h_{j}^3 \geq \left( \frac{1}{\frac{2m}{n}} \right)^2 \sum_{j \in J} h_{j} \geq \frac{C}{(m_n)^2}.
\]

Because \( N = m_n - 1 \) we get
\[
\|e\|_E \geq C N^{-1}
\]
no matter what kind of mesh is used. A similar result holds for \( p > 1 \).

The following theorem shows that we also have \( \|e\|_E \leq C N^{-P} \) provided the mesh is properly chosen (i.e., the above bound is attained).

**Theorem 2.4.** For the \( h \) version, using polynomials of fixed degree (i.e., \( p_n = p \)), the radical mesh using \( m \) elements with
\[
\beta_0 = \frac{p+1/2}{a-1/2}
\]
is the optimal choice and the following estimate holds with \( h = \frac{1}{m} \)
\[
\lim_{h \to 0} \frac{\|e\|_E}{h^P} = C(\alpha, p) \beta_0^{p+1/2}
\]
where
\[
C(\alpha, p) = \frac{a \Gamma(\alpha) \sin \pi \alpha \Gamma(p-\alpha+1)}{\sqrt{\pi} q^P \sqrt{2p+1} \Gamma(p+1/2)}.
\]

The notion of optimality in the above theorem is understood in the following sense. Consider the sequences of meshes defined by grading functions \( g \) that satisfy
1) \( g \in C^0(0,1) \cap C^1(0,1) \)
2) \( g(0) = 0, g(1) = 1 \) and \( g \) is strictly increasing.

Then the error of the FEM with the radical mesh (\( g = \beta_0 \)) attains the
minimum among all grading meshes satisfying (i), (ii).

It is better to (mildly) overrefine the mesh than underrefine it. This can be seen from the detailed results proven in [GiB]. Mild overrefinement is advisable in two and three dimensions as well.

We see that proper refinement leads to the highest possible rate in the h version, algebraic in $N$, which is the same as would be achieved if the solution were smooth. We also see that for fixed $p$ (uniform) the best mesh to use (for the h version) is a radical one, with the refinement strength (i.e., $\beta$) being increased linearly with $p$.

Let us now show that an exponential rate of convergence is possible using the h-p version. For simplicity, we will take a uniform degree vector. Let us take the mesh to be the geometric one defined by (2.7) and let

$$x_{i-1/2} = \frac{1}{2} (x_{i-1} + x_i).$$

As before, let $w(x) = u'(x)$, with $u(x)$ given by (2.2), and let us expand $w(x)$ in terms of its Taylor expansion about the point $x_{i-1/2}$ so that

$$w(x) = w_{1,p}(x) + r_{1,p}(x) = \sum_{k=0}^{p-1} a_k (x-x_{i-1/2})^k + \sum_{k=p}^{\alpha} a_k (x-x_{i-1/2})^k$$

where $w_{1,p}(x)$ is the Taylor polynomial of degree $p-1$ and $r_{1,p}(x)$ is the remainder. Using (2.7a) (2.7b) we see that

$$x_{i-1/2} = \frac{x_{i-1} + x_{i-1} - q^{-1}}{2} = q^{-1}(\frac{1+q}{2})$$

and

$$x_{i-1} - x_{i-1} = q^{-1}(\frac{1-q}{2}).$$

Now for $i = 2, 3, \ldots, m$, the function $x^{q^{-1}}$ is analytic in the closed disk
with center \( x_{i-1/2} \) and radius \( q^{(m-1)(1+q/2 - q^2)} = x_{i-1/2} - x_{i-2} \) for \( i \geq 3 \) and is bounded by \( Cq^{(m-1)(\alpha-1)} \) there. Hence, applying the Cauchy theorem,

\[
|a_k^{(1)}| \leq C q^{(m-1)(\alpha-1)} (1-q^k) k \left(\frac{1+q^2}{2} - q^2\right)^k \leq C q^{(m-1)(\alpha-1)} \rho^k
\]

where \( |\rho| = |2q+1|^{-1} < 1 \), and \( C \) is independent of \( m, i, k \). Hence

\[
|w(x) - w_{1,p}(x)| \leq C q^{(m-1)(\alpha-1)} \rho^p.
\]

Using this,

\[
E_1 = \inf_{r \in \mathcal{P}} \int_{I_1} |r - w|^2 \, dx \leq \int_{I_1} |w_{1,p} - w|^2 \, dx \\
\leq Cq^{2p} q^{(m-1)(2(\alpha-1)+1)}.
\]

Therefore

\[
\sum_{i=2}^{m} E_i \leq C \rho^{2p}.
\]

Next, we note that

\[
\inf_{r \in \mathcal{P}} \int_{I_1} |r - w|^2 \, dx \leq \int_{I_1} |w(x)|^2 \, dx = C q^{(m-1)(2\alpha-1)} \leq C \rho_1^m
\]

for some \( 0 < \rho_1 < 1 \). Choosing \( p = \kappa m \) and noting that \( N = m(p+1) \) we see using (2.31), (2.32) that

\[
\sum_{i=1}^{m} E_i \leq C e^{-\gamma N^{1/2}}
\]

for some \( \gamma > 0 \) and \( C \) depending on \( q, \alpha \).

Now let \( r_i \) denote the minimizer of \( E_i \) in (2.30). Then obviously
\[ \int_{I_1} (r_1 - \omega) \, dx = 0. \]

Hence defining \( R(x) \), \( 0 < x < 1 \), \( R(x) \big|_{I_1} = r_1(x) \) we can define

\[ \omega(x) = \int_0^x R(t) \, dt \]

and get \( \omega(1) = 0 \). Hence \( \omega \in V^n(\mathcal{J}, p) \) such that

\[ \| u - \omega \|_{H^1(1)} \leq C e^{-\gamma N^{1/2}} \]

and the exponential rate is proven. See [GiB] for a more precise analysis.

We have used \( p \) uniform. We could also analyze the case for \( p \) non-uniform. In fact, the optimal rate is obtained by choosing a non-uniform degree vector as described in the following theorems.

**Theorem 2.5.** Let \( [a] \) denote the integral part of \( a \). Then the optimal mesh for the h-p version is the geometric one with ratio \( q_0 = (\sqrt{2} - 1)^2 = e^{-1.7626} \approx 0.17 \) and linear degree vector \( p = \{ p_j \} \), \( p_j = [s_0] + 1 \) with \( s_0 = 2a - 1 \). With this choice,

\[ \| e \|_E \leq C(\alpha) q_0^{\sqrt{(a-1/2)N}} = C(\alpha) e^{-1.7626\sqrt{(\alpha-1/2)N}}. \]

Let us be more precise about the notion of optimality. We have the following theorem.

**Theorem 2.6.** For any mesh \( \mathcal{J} \), and any degree vector \( p \), we have

\[ \| e \|_E \leq C(\alpha) \frac{q_0^{\sqrt{(\alpha-1/2)N}}}{\sqrt{N^{(\alpha-1/2)}}}. \]

21
Comparing the lower bound given in Theorem 2.6 and the upper one in Theorem 2.5, we see that the geometrical mesh with linear degree vector gives the maximal possible exponential rate. This is the justification for calling the mesh optimal.

For any given geometrical mesh with ratio \( q \) there is an optimal \( p \), which is linear, with slope

\[
(2.35) \quad s_0(q) = (\alpha - \frac{1}{2}) \frac{\log q}{\log r}, \quad r = \frac{1 - \sqrt{q}}{1 + \sqrt{q}}.
\]

For this choice, the error is

\[
(2.36) \quad \|e\|_E \leq C(\alpha, q) e^{-\sqrt{(\alpha - 1/2)N} \sqrt{2 \log q \log r}}.
\]

This shows that for small \( q \) (i.e. strong refinement) the optimal degree vector is less uniform.

In Theorems 2.5 and 2.6, we have addressed the case when the degree vector was non-uniform. For uniform degree vector, the situation is slightly different.

**Theorem 2.7.** Let \( q \) be the ratio of the geometric mesh with \( m \) elements and \( p = s_0m \) where \( s_0 \) is as in (2.35). Then

\[
\|e\|_E \leq C(\alpha, q) e^{-\sqrt{(\alpha - 1/2)N} \sqrt{\log q \log r}} N^{-\sigma/2}
\]

where

\[
\sigma = \min(2\alpha - 1, \alpha).
\]

Substituting the optimal \( q = q_0 = (\sqrt{2} - 1)^2 \) in the above, we get

\[
(2.37) \quad \|e\|_E \leq C(\alpha) e^{-1.2464\sqrt{(\alpha - 1/2)N} \sqrt{\log q \log r}} N^{-\sigma/2}.
\]

Hence we see that the exponential rate is smaller in this case.
In practice (in higher dimensions) one often uses a geometric mesh with a fixed number of elements \( m \) and determines \( p \) more or less adaptively. For every \( m \), the convergence will have two parts. In the pre-asymptotic portion, the rate is exponential, while in the asymptotic region, one obtains the algebraic rate of the \( p \) version. If the result of a computation lies in the algebraic part, then this indicates the mesh was not sufficiently refined.

If a uniform degree vector \( p \) is to be used, then it turns out that the geometric mesh does not give the best rate of convergence, i.e. (2.37) is not optimal. We have seen that for fixed uniform degree vector (i.e. "\( h \) version), the radical (and not the geometric) mesh is optimal. If we consider a sequence of uniform degree vectors and perform the \( h \) version for each of them with the corresponding optimal radical mesh, then the envelope of the resulting error curves leads to an \( h-p \) version which gives better results than the one in Theorem 2.7. We have the following theorem.

Theorem 2.8. There is an \( h-p \) version with uniform degree vector \( p \) which has the following estimate

\[
\|e\|_E \leq \frac{C(\alpha)}{N^{(\alpha-1/2)/2}} e^{-4/e\sqrt{(\alpha-1/2)/N}} = \frac{C(\alpha)}{N^{(\alpha-1/2)/2}} e^{-1.4715\sqrt{(\alpha-1/2)/N}}.
\]

As \( N \to \infty \), the meshes (inside \((0,1)\)) tend to the geometric mesh with ratio \( q = e^{-4/e^2} \approx 0.5820 \) and the relation between the degree \( p \) and the number of elements \( m \) tends to a linear one,

\[
p \approx \frac{4}{e^2} \left( \alpha - \frac{1}{2} \right) m.
\]

Various numerical example illustrating these results may be found in 23.
2.5. The rate of convergence in the $L_2$ norm

So far, we have discussed the $h$, $p$ and $h$-$p$ versions in terms of their convergence in the energy norm. An important difference between the $h$ version on the one hand and the $p$ (and $h$-$p$) version on the other that is apparent from the above is that the $p$ (and $h$-$p$) version leads to different asymptotic rates of convergence, depending upon the nature of the solution (for example, the "doubling" effect for $x^a$ singularities). In contrast, the only way that the $h$ version could lead to a rate better than $O(h^p)$ for a function $f(x)$ is if $f$ were a polynomial of degree $\leq p$. This difference leads to an interesting effect in the error in the $L_2$ and other lower-order norms.

Let us discuss the error in the $L_2$ norm for the $p$ version. Suppose we first consider the case where $u \in H^k(\Omega)$. Then, by the usual duality argument, we may show

\begin{equation}
\|e\|_{L_2(\Omega)} \leq C p^{-1} \|e\|_E
\end{equation}

so that using Corollary 2.2,

\begin{equation}
\|e\|_{L_2(\Omega)} \leq C p^{-k} \|u\|_{H^k(\Omega)}
\end{equation}

In fact, the usual duality gives the following theorem.

**Theorem 2.9.** Let $p_n = p$ and $\{f^n\}$ be quasiuniform. Then for $k \geq 1$,

\begin{equation}
\|e\|_{L_2(\Omega)} \leq C h_\mu p^{-k} \|u\|_{H^k(\Omega)}
\end{equation}
where $\mu = \min(k,p+1)$.

This result is optimal in the sense that

$$\sup_{u \in H^k(\Omega)} \inf_{v \in V^n} \|u-v\|_{L^2(\Omega)} \geq Ch^{\mu-k}.$$

If we now consider the special singular solution of the form (2.2) (scaled on to $\Omega = (-1,1)$), then Theorem 2.9 predicts a rate of $O(h^{\min(\alpha+1/2,p+1)})$ for the $h$ version which is the best $L_2$ error possible with a quasiuniform mesh, and so is again optimal. For the $p$ version (using a single element), we may now use (2.39) with (2.23) to obtain an $O(p^{-2\alpha})$ rate. The optimal $L_2$ approximation estimate is essentially obtained by taking $\alpha = \alpha + 1$ in Theorem 2.3, which gives

$$\inf_{v \in V^n} \|u-v\|_{L^2(\Omega)} = C \mu^{\alpha+1/2} (1+O(p^{-1})).$$

Hence the estimate through duality does not give the optimal rate (2.41).

An improved estimate of $O(p^{-2\alpha})$ may be obtained using the fact that

$$\int_{-1}^{\infty} P_\alpha(t) dt = \frac{1}{2p+1} (P_{p+1}(x) - P_{p-1}(x))$$

Then by (2.20), the definition of $u_p$, and (2.42),

$$u - u_p = \sum_{r=p}^{\infty} \frac{a_r}{2r+1} \int_{-1}^{\infty} P_r(t) dt
= - \frac{a_p}{2p+1} P_{p-1} - \frac{a_{p+1}}{2p+3} P_p + \sum_{r=p}^{\infty} \left( \frac{a_r}{2r+1} - \frac{a_{r+2}}{2r+5} \right) P_{r+1}.$$
Define
\[ b_r = \frac{a_r}{2r+1} - \frac{a_{r+2}}{2r+5}. \]

Then, using (2.22), we obtain
\[ (2.44) \quad b_r = (-1)^{r-1} \frac{\bar{C}(\alpha)}{r^{2\alpha+1}} \left(1 + \frac{\bar{C}_1(\alpha)}{r} + \frac{\bar{C}_2(\alpha)}{r^2} + \cdots \right). \]

Now by (2.43),
\[ (2.45) \quad \| u - u_p \|_{L^2(\Omega)} = \frac{a_p^2}{(2p+1)^2} \frac{2}{2p-1} + \frac{a_{p+1}^2}{(2p+3)^2} \frac{2}{2p+1} + \sum_{r=p}^{\infty} b_r \frac{2}{2r+3}. \]

Using (2.44), we that
\[ \sum_{r=p}^{\infty} b_r \frac{2}{2r+3} = \frac{C(\alpha)}{p^{2\alpha+2}} (1 + O(p^{-1})). \]

Also, (2.22) may be used to estimate the first two terms in (2.45), yielding
\[ (2.46) \quad \| u - u_p \|_{L^2(\Omega)} = \frac{C(\alpha)}{p^{2\alpha+1/2}} (1 + O(p^{-1})). \]

This estimate, which is sharp, improves upon the rate through duality, but still does not give the optimal rate of (2.41).

Hence we observe the following interesting difference when the projection of the function (2.2) is taken in the energy norm (equivalent to the \( H^1(\Omega) \) norm). For the \( h \) version, the convergence in the \( L^2(\Omega) \) norm is
\[ \min(\alpha+1, p+1) \quad O(h^{-\alpha-1/2}) \] which is optimal. For the \( p \) version, this \( L^2 \) convergence rate is higher, \( O(p^{-2\alpha-1}) \). However, the best \( L^2 \) approximation is \( O(p^{-2\alpha-1}) \), which is even higher. As a result, the energy norm projection does not give the optimal rate in the \( L^2(\Omega) \) norm for (2.2) (see [Erm]). Let us mention that in two dimensions, for non-convex domains, the \( h \) version also fails to give the optimal \( L^2 \) rate (see Section 4.2).
2.6. The pollution problem

Let us now discuss another difference between the h and p versions, the phenomenon of pollution. This is said to occur when unsmoothness of the solution $u$ in one area influences the accuracy of $u_h$ in a different area, i.e. the accuracy of $u_h$ is not governed by the local behavior of the solution $u$. Consider once again the approximation of $u$ given by (2.2) (scaled to $(-1,1)$). Then $u$ is singular at $x = -1$ but analytic everywhere else. Suppose $u_h$ is the solution using the h version (say with $p = 1$). Then it is easily proven that $u_h$ is simply the linear interpolant of $u$ (an analogous result holds for higher $p$). This obviously implies that the error $|(u-u_h)(0)|$, for instance, will only depend upon the behavior of $u$ over some $I_\delta = (-\delta,\delta)$, $\delta < 1$. In fact, this error will converge asymptotically at the optimal rate of $O(h^2)$ for any $\alpha > 1/2$. Also, for fixed $\delta$, $\|u - u_h\|_{H^r(I_\delta)}$ for $r = 0,1$ will similarly converge at the optimal rate $O(h^{2-r})$, even if $1/2 < \alpha < 3/2$. We say, therefore, that the h version is free from pollution. The h-p version also turns out to be free from pollution, so that $|(u-u_h)(0)|$, $\|u - u_h\|_{L^2(I_\delta)}$, $\|u - u_h\|_{H^1(I_\delta)}$ all converge at the optimal exponential asymptotic rate determined only by the analytic behavior of $u$ over $I_\delta'$, $\delta'$ close to $\delta$, and not by the singular behavior at $x = -1$.

The p version, however, will exhibit pollution. Suppose we use a single element, then with $u_h = u_p$, $\|u - u\|_{H^r(I_\delta)}$ will be no better than the rates given in Sections 2.3, 2.5 for $\|u - u_p\|_{H^r(I)}$, even though $u$ is...
analytic over $I_\delta$. More precisely, let us estimate $|(u-u_p)(0)|$. Noting that

$$P_{2m}(0) = (-1)^m \frac{(2m-1)!!}{(2m)!!} = \frac{(-1)^m}{\sqrt{\pi}} \frac{1}{\sqrt{m}} \left(1 + \frac{C}{m} + \ldots\right)$$

we get for $p = 2k+1$ (using (2.22), (2.43)),

$$(u-u_p)(0) = \frac{C}{p^{2\alpha+1/2}} (1 + O(p^{-1})) + \sum_{m=k}^{\infty} \frac{(-1)^{m+1}b_{2m+1}}{\sqrt{\pi} (m+1)^{1/2}} (1 + O(m^{-1})).$$

Using (2.44) and the fact that the terms in the sum are of alternating sign, this gives

$$|(u-u_p)(0)| = \frac{C}{p^{2\alpha+1/2}} (1 + O(p^{-1})).$$

The above estimate clearly shows the pollution effect, since in the absence of pollution, the rate should be exponential because $u$ is analytic over $I_\delta$. The reason this pollution occurs should be clear: unlike the h and h-p versions, there is no "buffer" of elements now to isolate the singularity from the rest of the domain. For more results on the pollution problem in 1-d, in the context of the h-version as well, see [W2].

3. THE MODEL PROBLEM AND ITS REGULARITY

In Section 2 we have analyzed a simple one-dimensional problem. In this section we will address the basic properties of the analogous problem in two and three dimensions. We will concentrate on model problems which are typical in the field of computational (structural) mechanics. Mathematically, these problems are described by elliptic partial differential equations with piecewise analytic input data. Hence this class of practical problems is relatively narrow and well-defined, leading to several important common
properties which can be exploited in their numerical treatment.

3.1. The model problem

Let us consider \( \Omega \subset \mathbb{R}^2 \) to be a curvilinear polygon with piecewise analytic boundary \( \partial \Omega = \bigcup_{i=1}^{M} \Gamma_i \), where \( \Gamma_i \) is an open arc and \( \bar{\Gamma}_i \) is analytic with endpoints \( A_1, A_{i+1} \). \( A_1, A_2, \ldots, A_M \) \((A_{M+1} = A_1)\) will denote the vertices with internal angles \( \omega_1, \ldots, \omega_M \) \((0 < \omega_i \leq 2\pi)\). (If \( \omega_i = 2\pi \) then we have a slit domain, where the boundary is taken to be two-sided.) \( \partial \Omega \) will consist of two parts, \( \Gamma_D = \bigcup_{i \in \mathcal{D}} \bar{\Gamma}_i \) being the Dirichlet boundary and \( \Gamma_N = \partial \Omega - \Gamma_D \) the Neumann boundary. It should be noted that the boundary could be smooth everywhere, with \( \omega_1 = \pi \) and the vertices \( A_i \) corresponding to points in \( \bar{\Gamma}_D \cap \bar{\Gamma}_N \). We restrict ourselves to simply connected domains in order to simplify the exposition. An example of such a domain and the notation used is shown in Figure 3.1.

In three dimensions, \( \Omega \subset \mathbb{R}^3 \) will be a bounded domain with piecewise analytic boundary \( \partial \Omega \) consisting of faces \( \Gamma_i \), \( i = 1, 2, \ldots, M \) which are curved polygons in \( \mathbb{R}^3 \), joined by edges \( \gamma_i \), \( i = 1, \ldots, n_e \) (curves in \( \mathbb{R}^3 \)) and vertices \( A_i \), \( i = 1, \ldots, n_v \). \( \Gamma_D \) and \( \Gamma_N \) will now be the union of some faces.

Analogously to (2.1) we will now consider a second order elliptic system of \( m \) differential equations for the unknown vector \( \mathbf{u} = (u_1, \ldots, u_m) \) which, when cast into the variational form \( B(\mathbf{u}, \mathbf{v}) \) (analogous to (2.3)), is given by

\[
(3.1) \quad B(\mathbf{u}, \mathbf{v}) = \int_{\Omega} \left( \sum_{i,j=1}^{m} \sum_{k,l=1}^{t} a_{ijkl} \frac{\partial u_i}{\partial x_k} \frac{\partial v_j}{\partial x_l} \right) dx,
\]

where \( t = 2 \) or \( 3 \) and \( a_{ijkl} = a_{jilk} \) are analytical functions on \( \Omega \).
satisfying an ellipticity condition that ensures

(3.2) \[ B(\overline{\mathbf{u}}, \overline{\mathbf{v}}) = \gamma |\overline{\mathbf{u}}|^2_{H^1(\Omega)}, \quad \gamma > 0. \]

Let

\[ H^1_D(\Omega) = \{ \overline{\mathbf{v}} \in H^1(\Omega), \quad \overline{\mathbf{v}} = 0 \text{ on } \Gamma_D \}. \]

Figure 3.1. The scheme of the domain.

Further, let \( F \) be a continuous functional, defined on \( H^1_D(\Omega) \), of the form

(3.3) \[ F(\overline{\mathbf{v}}) = \int_\Omega \overline{\mathbf{r}} \cdot \mathbf{v} \, dx + \int_{\Gamma_N} \overline{\mathbf{r}} \cdot \mathbf{v} \, ds, \]

where \( \overline{\mathbf{r}} \) is analytic on \( \overline{\Omega} \) and \( \overline{\mathbf{r}} \) is analytic on every \( \overline{\Gamma}_1, \Gamma_1 \subset \Gamma_N \). Moreover, let \( \overline{\mathbf{g}} \in C^0(\Gamma_D) \) and be analytic on each \( \overline{\Gamma}_1 \subset \Gamma_D \). Then our model problem consists of finding \( \overline{\mathbf{u}} \in H^1(\Omega) \) satisfying

(3.4) \[ B(\overline{\mathbf{u}}, \overline{\mathbf{v}}) = F(\overline{\mathbf{v}}), \quad \forall \overline{\mathbf{v}} \in H^1_D(\Omega) \]
For \( \Gamma_D \neq \emptyset \) condition (3.2) ensures a unique solution for (3.4), (3.5). For \( \Gamma_D = \emptyset \), we need an additional solvability condition on \( \Gamma \) and \( \Omega \).

We define the energy norm \( \| \mathbf{u} \|_E = (B(\mathbf{u}, \mathbf{u}))^{1/2} \), equivalent to the \( H^1(\Omega) \) norm. By restricting our formulation to the above class of problems with piecewise analytic input data of the type assumed, we ensure that the solution will possess certain special features.

Let us now describe two concrete examples of (3.4), (3.5) in which we shall be interested. The first problem is a scalar one \((m = 1)\) on \( \Omega \subset \mathbb{R}^t \) with

\[
B(u,v) = \int_{\Omega} \left( \sum_{i=1}^{t} k_i \frac{\partial u}{\partial x_i} \frac{\partial v}{\partial x_i} \right) dx, \quad k_i = k_i(x) > 0.
\]

For \( k_i \neq 1 \), we get the Poisson equation. Most of our cited results are in two dimensions but the three-dimensional problem will be mentioned as well.

In Section 5 we will be interested in the 2-d case that \( k_i \) are very different in magnitude. Obviously, in this case the bilinear form becomes degenerate as \( k_i/k_2 \to 0 \). Such degeneracies occur in various practical problems and their numerical treatment will be discussed in the context of the robustness of the method (the locking problem).

Our second example is the typical model problem of two-dimensional elasticity. We will restrict ourselves here to the case of isotropic homogeneous materials (plane strain) for which \( m = 2, t = 2 \) and

\[
B(\mathbf{u}, \mathbf{v}) = \frac{E}{1+\nu} \int_{\Omega} \left[ \sum_{i,j=1}^{t} \epsilon_{i,j}(\mathbf{u}) \epsilon_{i,j}(\mathbf{v}) + \frac{\nu}{1-2\nu} (\text{div } \mathbf{u} \text{ div } \mathbf{v}) \right] dx.
\]
where \( E > 0 \) is the Young's modulus of elasticity, \( 0 \leq \nu < \frac{1}{2} \) is the Poisson ratio and

\[
(3.8) \quad \varepsilon_{ij}(\tilde{u}) = \frac{1}{2} \left( \frac{\partial u_i}{\partial x_j} + \frac{\partial u_j}{\partial x_i} \right).
\]

For \( \nu \to \frac{1}{2} \), the bilinear form becomes degenerate as in the first example. Nevertheless, there is a large difference in these two cases. In the first one, the solution becomes less smooth as the form degenerates and in the limiting case the problem is not elliptic. In the second example, the degeneracy does not influence the smoothness of the solution and the limiting case is still an elliptic problem. See e.g. [AF1] for a general analysis of the degeneracy with the smoothness being preserved.

In three-dimensional elasticity the formulation is analogous to (3.7), (3.8).

### 3.2. Regularity of the two-dimensional problem

Let us now analyze the regularity of the solution of our model problem for \( \Omega \subset \mathbb{R}^2 \). First assume that \( \partial \Omega \) is smooth and that either \( \Gamma_N = \partial \Omega \) or \( \Gamma_D = \partial \Omega \). Then

\[
(3.9) \quad \| \tilde{u} \|_{H^k(\Omega)} \leq C \left( \| \tilde{\Phi} \|_{H^{k-2}(\Omega)} + \| \tilde{g} \|_{H^{k-1/2}(\Gamma_D)} + \| \tilde{h} \|_{H^{k-3/2}(\Gamma_N)} \right).
\]

If the data (i.e., \( \partial \Omega, \tilde{g}, \tilde{h} \)) are not sufficiently smooth (at vertices) or both parts \( \Gamma_N \) and \( \Gamma_D \) are present, then (3.9) does not hold. (Note that the analyticity of \( \tilde{g}, \tilde{h} \) is only on \( \bar{\Gamma}_1 \).) However, under our assumptions, \( \tilde{u} \) will be analytic on \( \tilde{\Omega} - \bigcup_{1=1}^{M} A_1 \) and in the neighborhood of a vertex \( A_1 \), for the scalar equation (\( m = 1 \)), \( u \) will have the form (see [CS], [Gr2], [Gr4], [Da])

32
\begin{equation}
(3.10) \quad u = \sum_{j=1}^{J} \sum_{s=0}^{S} \sum_{t=0}^{T} c_{jst} \psi_{jst}(\varphi) r^{j+s+t} \ln r + u_0 = \sum_{j=1}^{J} u_j^{(1)} + u_0
\end{equation}

Here, \((r, \varphi)\) are local polar coordinates at the vertex \(A_i\) and \(\psi_{jst}\) are analytic function in \(\varphi\). The decomposition \((3.10)\) for a proper choice of \(J, S, T\) is such that the remainder \(u_0\) is smoother than the terms in the sum and belongs to a weighted space with the weight \(w(r) \to \infty\) as \(r \to 0\).

In the simpler case of the Laplace operator on a straight-sided polygonal domain and homogeneous boundary conditions, one has \(S = T = 0\), i.e.,

\begin{equation}
(3.10)' \quad u = \sum_{j=1}^{J} c_j \psi_j(\varphi) r^{\alpha_j} + u_0
\end{equation}

except for the case of \(\alpha_j\) being an integer, when \(S = 1\). The \(\psi_j\) and \(\alpha_j\) are explicitly known,

\[
\alpha_j = \begin{cases} 
j\pi/\omega & \text{for Dirichlet or Neumann boundary conditions} \\
(2j-1)\pi/2\omega & \text{for mixed boundary conditions}
\end{cases}
\]

\[
\psi_j(\varphi) = \begin{cases} 
\sin \alpha_j \varphi & \text{for the Dirichlet or mixed problem} \\
\cos \alpha_j \varphi & \text{for the Neumann problem.}
\end{cases}
\]

The mixed problem here means Dirichlet conditions for \(\varphi = 0\) and Neumann conditions for \(\varphi = \omega\).

The number \(T\) in \((3.10)\) is different from 0 in the following cases:

1) The coefficients of the differential operator are not constant;

2) the operator contains lower order terms;
iii) the arcs $\Gamma_1$ are curvilinear.

In the case of the Laplacian, the singular term due to the curvature of $\Gamma_1$ (with $j = 1$, $t = 1$) may be less regular than the second term ($j = 2$, $t = 0$). This happens if

$$\alpha_1 + 1 < \alpha_2 \text{ i.e. } \omega < \pi$$

for all three boundary conditions.

The coefficients $c_{jst}$ have physical meaning and often are called stress intensity factors. For details, see [CS]. We see here that the solution is more complicated, but the singularity of the solution is analogous to (2.2), the singularity in the one-dimensional setting.

Taking the correct number of terms in (3.10), we can write

$$u = \sum_{i=1}^{M} \sum_{j} u_{ij} x_i + u_0,$$

where $x_i$ is a smooth cut-off function in the neighborhood of $A_1$. Then we have

$$\|u_0\|_{H^k(\Omega)} \leq C \left( \|f\|_{H^{k-2}(\Omega)} + \sum_{\Gamma_1 \subset \Gamma_D} \|g\|_{H^{k-1/2}(\Gamma_1)} + \sum_{\Gamma_1 \subset \Gamma_N} \|h\|_{H^{k-3/2}(\Gamma_1)} \right)$$

i.e., the smoothness of $u_0$ is analogous to that when the domain and data are smooth as in (3.9). Note that although the exponents $\alpha_j$ in (3.10) depend continuously on the angles $\omega_1$, the value of $S$ does not. As a result, the constant $C$ in (3.12) can be very large for angles $\omega_1$ close to the exceptional values for which $S = 1$ instead of $0$. This problem can be avoided, however, by adjusting the expansion of $u$. 
So far, we only discussed the scalar equation, which was represented by the example of Laplace's equation. The case of a system (like the elasticity problem) is more complicated but very analogous. See, for example, [Sa]. The decomposition (3.10) still holds with \( \psi_{jst}(\varphi) \) being a vector. In contrast to the previous case, \( \alpha_j \) may now be complex (and there are more exceptional cases than before).

3.3. Countably normed spaces

The regularity of the solution in our setting of piecewise analytic input data can be very advantageously described in the form of countably normed spaces. (For details, see [BG2],[BG3].) For \( \delta < 1 \), let \( S_\delta = \{ x \in \Omega, r < \delta \} \) and let \( \Phi_\beta = r^\beta \), for \( 0 < \beta < 1 \). By \( H_\beta^{m,\ell}(S_\delta) \), \( m \geq \ell > 0 \) integers, we denote the completion of the set of all infinitely differentiable functions under the norm

\[
\| u \|_{H_\beta^{m,\ell}(S_\delta)}^2 = \| u \|_{H_\delta^{\ell-1}(S_\delta)}^2 + \sum_{k=\ell}^{m} \| \Phi_\beta u \|_{L_2(S_\delta)}^2
\]

where

\[
|D^m u|^2 = \sum_{|\alpha| = m} |D^\alpha u|^2
\]

and

\[
\| u \|_{H_\beta^{m,0}(S_\delta)}^2 = \sum_{k=0}^{m} \| \Phi_\beta u \|_{L_2(S_\delta)}^2.
\]

Further, let us introduce \( H_\beta^{\ell}(S_\delta) \), \( \ell \geq 0 \) consisting of those \( u \in H_\beta^{\ell,\ell}(S_\delta) \) for which

\[
\| \Phi_\beta u \|_{L_2(S_\delta)}^2 \leq C d^{k-1}(k-\ell)!,
\]

35
Finally, introduce the space $\mathcal{E}^2_\beta(S_\delta)$,

$$\mathcal{E}^2_\beta(S_\delta) = \{ u \in H^2_\beta(S_\delta) \mid |Du|^2(x) \leq C_\delta k! (\beta + k - 1)^{-1} \}.$$  

We then have

$$B^2_\beta(S_\delta) < \mathcal{E}^2_\beta(S_\delta) < B^2_{\beta + \epsilon}(S_\delta), \quad \epsilon > 0 \text{ arbitrary.}$$

Under the assumptions made above, we have the following theorem.

**Theorem 3.1** [BG2] [BG3]. Let $u$ be the solution corresponding to a second order (scalar) operator.

1) Let $\Omega$ be a polygon (i.e., $\Gamma_i$ are straight lines). Then the solution $u$ is analytic in $\Omega - \bigcup_{i=1}^m \Lambda_i$, and in the neighborhood of any $\Lambda_i$, we have $u \in B^2_\beta(S_\delta)$, $\forall \delta < \delta_0$ where $0 < \beta < 1$ depends on the angle $\omega_i$.

ii) For curvilinear polygons, $u$ is analytic on $\Omega - U\Lambda_i$ and $u \in \mathcal{E}^2_\beta(S_\delta)$, $0 < \beta < 1$.

The characterization given in Theorem 3.1 is very useful in practice because it allows us to construct the $h$-$p$ version so that its convergence rate is exponential, without knowing the decomposition form (3.10) exactly. For analogous results for problem (3.7) we refer to [GoB].

### 3.4. Regularity of the three-dimensional problem

In three dimensions, the situation is more complicated. The solution is again analytic on $\Omega \setminus (\text{edges and vertices})$. There is a singular behavior in the neighborhood of the edges (called edge singularities) where the solution is smooth along the edge and singular in the direction perpendicular to the
edge. In addition, near the vertices, the solution has an essentially radial vertex singularity, as well as transitional edge-vertex singularities that appear in cones around the edges. The stress intensity factors, i.e analogues to \( c_{jst} \), are now functions along the edges which are smooth but have singular behavior in the neighborhood of the endpoints (the vertices). These stress intensity factors may be bounded or unbounded in dependence on the geometry. In addition, there are stress intensity factors for the vertex singularity.

As an illustration, let us present some regularity results from [Gr1] [Gr2], (which were used in [D2] for the following model problem: Given \( f \in C^\infty(\bar{\Omega}) \), find \( u \) such that

\[
\begin{align*}
(3.13) \quad -\Delta u &= f \quad \text{in} \ \Omega \\
 u &= 0 \quad \text{on} \ \Gamma,
\end{align*}
\]

where \( \Omega \) is a polyhedral domain in \( \mathbb{R}^3 \).

First we describe the singularity close to \( \gamma \), an open segment of an edge \( \gamma_0 \) of \( \Omega \) assumed to be along the \( x_3 \)-axis. Let the dihedral angle at \( \gamma_0 \) be \( \omega \), then this singularity is characterized only by the values \( \alpha_j = j\pi/\omega \). Then, in a neighborhood \( G \) of \( \gamma \) (such that \( \bar{G} \) does not intersect any corners or edges other than \( \gamma_0 \)) we have, analogously to (3.10) (for \( f \in C^\infty(\bar{\Omega}) \)),

\[
(3.14) \quad u = \sum_{j=1}^{J} \sum_{s=0}^{S} \sum_{t=0}^{T} c_{jst}(x_3)\psi_{jst}(\phi) r^{\alpha_j t} \ln^{s}(r) + u_0,
\]

where \( (r,\phi,x_3) \) are local cylindrical coordinates, \( c_{jst} \in C^\infty(\gamma) \) and \( \psi_{jst} \) are analytic in \( \phi \). As before, \( S = 0 \) for the problem (3.13) except for integral values of \( \alpha_j \), for which \( S = 1 \).
Next, we consider the behavior of \( u \) in the neighborhood of a vertex. Let \( A_0 \), the vertex, be situated at the origin and be an endpoint of the edge \( \gamma_0 \). Let \( K \) be the intersection of \( \Omega \) and a solid circular cone with vertex \( A_0 \) and axis along \( \gamma_0 \), which does not intersect any edge having an end point at \( A_0 \) other than \( \gamma_0 \). Let \( K_0 = K \cap B_0 \) where \( B_0 \) is a small ball of radius \( \rho_0 \) about the origin. Also, let \( \tilde{K} \) be the infinitely extended cone coinciding with \( \Omega \) in a neighborhood of the origin and let \( R = K \cap S \), where \( S \) is the unit sphere. Let \( 0 < \lambda_1 \leq \lambda_2 \leq \cdots \) be the eigenvalues of the Laplace-Beltrami operator on \( R \) with homogeneous Dirichlet boundary data on \( \partial R \), let \( \beta_1 = \sqrt{\lambda_1 + \frac{1}{4}} \), and let \( \tilde{\alpha}_j \) be the coefficients (corresponding to the singular coefficients \( \alpha_j \) above) for this operator for the singularity in a neighborhood of a vertex with angle \( \omega \). Then, with \( f \in C^\infty(\Omega) \), we have on \( K_0 \)

\[
(3.15) \quad u = \sum_{j=1}^{J} \sum_{s=0}^{S} \sum_{t=0}^{T} \left( \sum_{i=0}^{I} c_{jst} \rho^{-1/2+\beta_1} + f_{jst}(\rho) \right) \tilde{\psi}_{jst}(\phi) \theta^j \ln^5(\theta) \\
+ \sum_{i=0}^{I} d_i \rho^{-1/2+\beta_1} \tilde{g}_i(\theta,\phi) + \tilde{u}_0,
\]

where \( (\rho,\theta,\phi) \) are the local spherical coordinates, with origin at \( A_0 \), and with \( \theta = 0 \) corresponding to edge \( \gamma_0 \). The functions \( f_{jst}, \tilde{\psi}_{jst}, \tilde{g}_i \) and \( \tilde{u}_0 \) are all smooth, their exact regularity being stated in Lemma 4.3 of [D2]. Here, the first term is a transitional edge-vertex singularity corresponding to edge \( \gamma_0 \) and vertex \( A_0 \), while the second term is the vertex singularity of radial type which will be present throughout \( R \).

Equations (3.14), (3.15) then characterize the two types of singular
behavior (near edges and vertices) for three-dimensional problems. For the singular behavior of the elasticity problem we refer to [Gr3], [Pe].

We see that in two dimensions (for the Laplace equation), the singularity depends only on \( \omega_j \), the interior angles of the domain and upon the boundary conditions. In three dimensions, the situation is similar for edge singularities, but not for the singularities near vertices, which depend upon the Laplace-Beltrami eigenvalues. This makes the a priori determination of the exact form of the singularities more difficult. In general, for any vertex \( A_\ell \), \( \ell = 1, 2, \ldots, n_v \), we now define \( S_\ell \) to be the portion of the unit sphere subtended by the infinite cone which coincides with \( \Omega \) in a neighborhood of \( A_\ell \). Then, in addition to the dihedral angles \( \omega_j \), \( j = 1, 2, \ldots, n_e \), the regularity also depends upon the coefficients \( \beta_{\ell, i} = \sqrt{\lambda_{\ell, i} + \frac{1}{4}} \), where \( 0 < \lambda_{\ell, 1} \leq \lambda_{\ell, 2} \leq \ldots \) are the eigenvalues of the Laplace-Beltrami operator on \( S_\ell \) with homogeneous boundary conditions on \( \partial S_\ell \). These eigenvalues are directly related to the Steklov eigenvalue problem on \( K_0 \). See [BPA].

Let us mention that we can also describe the regularity of the solution using weighted spaces, as was done in Theorem 3.1 for the two-dimensional case. Unlike the isotropic weighted spaces used in the two-dimensional case, however, the natural spaces in three dimensions to describe the different types of singular behavior in the neighborhood of edges and vertices are anisotropic. We will not elaborate on this here. For regularity and error estimates in terms of countably normed spaces in 3-d, we refer to [BG4].

39
4. FINITE ELEMENT SPACES AND BASIC APPROXIMATION RESULTS

We will proceed here analogously to the one-dimensional case. Assume for now that $\beta = 0$. Then as in Sections 2 and 3 we will consider $V^n \subset H_D^1(\Omega)$ to be the finite element space of dimension $N_n$ (number of degrees of freedom). The finite element solution $\hat{u}_n \in V^n$ is then such that

$$B(\hat{u}_n, \hat{v}) = F(\hat{v}) \quad \forall \hat{v} \in V^n.$$

4.1. The finite element spaces

We will now restrict ourselves mainly to the two-dimensional problem and define the meshes and elements analogously as before.

Let $S = (-1, 1)^2$ and $T = \{\xi, \eta | 0 < \eta < (\xi + 1)\sqrt{3}, -1 < \xi < 0, \ 0 < \eta < (1-\xi)\sqrt{3}, 0 < \xi < 1\}$ be the standard square and standard triangle respectively. The mesh $T^n$ is now a partition of $\Omega$ into the open curvilinear triangles or quadrilaterals denoted by $T_j^{(n)}$, $j = 1, 2, \ldots, m_n$ such that to every $T_j^{(n)}$, we may associate an invertible mapping $F_j^{(n)}$ given by

$$x = X_j^{(n)}(\xi, \eta), \ y = Y_j^{(n)}(\xi, \eta),$$

which is a one-to-one mapping of $S$ onto $T_j^{(n)}$ (if $T_j^{(n)}$ is a quadrilateral) or $T$ onto $T_j^{(n)}$ (if $T_j^{(n)}$ is a triangle.) The functions $X_j^{(n)}$ and $Y_j^{(n)}$ are analytic functions in $\xi, \eta$ and have certain uniform bounds of the form

$$|D^\alpha X_j^{(n)}|, \ |D^\alpha Y_j^{(n)}| \leq C_1 h_n, j, \ |\alpha| = 1, 2 \quad (4.1a)$$

$$C_2 h_n^2 < |J_j| \leq C_3 h_n^2 \quad (4.1b)$$

where $J_j$ is the Jacobian and the constants $C_i$, $i = 1, 2, 3$, are independent
of the partition. Here, $h_{n,j}$ is the maximal side length of $\tau_j^{(n)}$. There are some additional assumptions about uniform analytic extendability of functions $X_j^{(n)}(\xi,\eta)$ and $Y_j^{(n)}(\xi,\eta)$ outside $S$ or $T$, which we will not discuss here (see [BG1]). In the case when $\tau_j^{(n)}$ are straight sided triangles, the mapping $F_j^{(n)}$ is linear and the uniformity condition is satisfied if the minimal (and maximal) angle is bounded from below (and above) by constants independent of the partition under consideration. Further, we assume that $\tau_j^{(n)} \cap \tau_k^{(n)}$ is either empty, a single common vertex or a single common side and if $s_{j,k} = \tau_j^{(n)} \cap \tau_k^{(n)}$ is the common side, then $x \in s_{j,k}$ is mapped by the mappings $(F_j^{(n)})^{-1}$ and $(F_k^{(n)})^{-1}$ into the same relative position on the side of the standard square or triangle. The above assumptions about the partitioning are analogous of the ones for the partitioning of $(0,1)$ into $I^n_j$ in Section 2.

Now we will define the quasiuniform, radical and geometric meshes.

1) **The quasiuniform mesh.** Here $\{\mathcal{G}^n\}$ is such that

$$\frac{h_n}{h_{n,j}} \leq \tau,$$

where $h_n = \max h_{n,j}$ and $\tau > 0$ is uniform for all meshes $\mathcal{G}^n$ under consideration.

2) **The radical mesh** with power $\beta \geq 1$ (with respect to $A_i$). If $A_i \notin \tau_j^{(n)}$, then with $\gamma = 1 - \frac{1}{\beta}$

$$C_1 h_n \left( \min_{x \in \tau_j^{(n)}} r^\gamma(x) \right) \leq h_{n,j} \leq C_2 h_n \left( \max_{x \in \tau_j^{(n)}} r^\gamma(x) \right)$$

and if $A_i \in \tau_j^{(n)}$, then
Here \( r(x) \) is the distance of \( x \) from \( A_I \). This definition, which holds in a neighborhood of \( A_I \), is the two-dimensional version of the definition of the radical mesh introduced in Section 2. It was employed in [BKPI].

3) The geometric mesh with ratio \( q, 0 < q < 1 \) (in the neighborhood of \( A_I \)). Let us number the elements of \( \mathcal{J}^n \) by a double index \( \tau_{j,k}^{(n)} \) with \( j = 1, \ldots, p(k), \ p(k) \leq \rho_0 \) and \( k = 1, 2, \ldots, n+1 \). Then if \( d_{n,j,k} \) denotes the distance between \( \tau_{j,k}^{(n)} \) and \( A_I \) and \( A_I \notin \tau_{n,j,k}^{(n)} \),

\[
C_1 h_n \left( \max_{x \in \tau_{j,k}^{(n)}} r(x) \right) \leq h_{n,j} \leq C_2 h_n \left( \max_{x \in \tau_{j,k}^{(n)}} r(x) \right).
\]

\[
C_1 q^{n+2-k} \leq d_{n,j,k} \leq C_2 q^{n+1-k}
\]

\[
\kappa_1 d_{n,j,k} \leq h_{n,j,k} \leq \kappa_2 d_{n,j,k}
\]

\( j = 1, \ldots, p(k), \ k = 2, \ldots, n+1 \).

If \( A_I \in \tau_{j,k}^{(n)} \) then \( k = 1 \) and

\[
\kappa_3 q^n \leq h_{n,j,k} \leq \kappa_4 q^n, \ j = 1, \ldots, \rho(1),
\]

where constants \( C_1 \) and \( \kappa_1 \) are uniform for all meshes under consideration.

The geometric mesh will be denoted by \( \{\mathcal{G}_q^n\} \). For more about these meshes we refer to [BG1].

As in one dimension, the radical and geometric meshes are refined in neighborhoods of the vertices \( A_I \) where the singularities of the solution are located.

Now let us define

\[
\tilde{V}^n(\Omega) = \{ u \in H^1(\Omega) | \text{For } x \in \tau_{j,k}^{(n)}, \ u \circ F_{j}^{(n)} \in U_{p_j}(K) \}.
\]
where \( K = T \), \( U_p = P_p \) if \( \tau_j^{(n)} \) is a curvilinear triangle and \( K = S \),

\[ U_p = Q_p \text{ or } Q_p' \] if and \( \tau_j^{(n)} \) is a curvilinear quadrilateral. Here,

\[ P_p(K) \] denotes the set of all polynomials of total degree \( \leq p_j \) on \( K = T \)

or \( K = S \). Further, \( Q_p(S) \) is the set of all polynomials of degree \( \leq p_j \)
in each variable and \( Q_p'(S) = P_p(S) \oplus \{ \xi_j, \eta_j \} \). The elements \( Q_p \) are

sometimes called product elements, while \( Q_p' \) are called serendipity or trunc elements, see e.g. [Ci], [SzB]. They are used for example, in the programs

MSC/PROBE, FIESTA, PEGASYS and STRIPE.

The theoretical question of under what conditions the space \( Q_p(S) \) or \( Q_p'(S) \) is preferable has not yet been resolved sufficiently. For the analysis of some aspects of this problem, we refer to [BE].

We now set \( V^n = \tilde{V}^n \cap H_D^1(\Omega) \). In general, the condition \( u \in H_D^1(\Omega) \) in the
definition of \( V^n(\Omega) \) restricts the functions \( u(F_j^{(n)}(x)) \) to belong to a
proper subset of \( P_p(T) \) (or \( Q_p(S), Q_p'(S) \)). This is necessary to enforce
continuity of \( u \in H_D^1(\Omega) \), for example, when \( p_j \) is different over adjacent
elements.

As in Section 2, we now deal with the space \( V^n(J^n, p_n) < H_D^1(\Omega) \) where
\( p_n \) is the degree vector. We note that for quasiuniform meshes and uniform
\( p_n = p_n \), we have

\[ N_n \approx h_n^{-2} p_n^2 \]

which also holds for the case of radical meshes and uniform \( p_n \). For the
geometric mesh and uniform \( p_n \) we have

\[ N_n \approx n p_n^2 \]
Analogously as in Section 2, we have in our model problem (3.1), $B(\tilde{u}, \tilde{v}) = \|\tilde{u}\|^2_E \approx \|\tilde{v}\|^2_{H^1(\Omega)}$ and hence with $\tilde{e}_n = \tilde{u} - \tilde{v}_n$ as before,

$$\|\tilde{e}_n\|_E \leq \inf_{\tilde{v} \in V^n(J^n, P_n)} \|\tilde{u} - \tilde{v}\|_E.$$

4.2. Approximation theorems

Let us now formulate theorems analogous to the ones in Section 2 from which the error bounds for the finite element method immediately follow.

Let us remark that when comparing the performance of different elements like $Q_p$ and $Q'_p$ elements (for example), the rate of convergence with respect to $p$ or $h$ in the theorems below may be the same, but the performance may be different (see [BE]). First, for the h-p version with quasiuniform meshes, we have the following theorem.

Theorem 4.1 [BS1]. Let us consider $p_n = p$ and a family of quasiuniform meshes $J^n$. Then

$$\inf_{\tilde{v} \in V^n} \|\tilde{u} - \tilde{v}\|_{H^1(\Omega)} \leq C h^{-\mu+1} p^{-(k-1)} \|\tilde{u}\|_{H^k(\Omega)},$$

where $\mu = \min(p+1, k)$ and $C$ is a constant independent of $h, p$ and $\tilde{u}$.

From the above theorem, it is clear that for the p version, we have a rate of $O(p^{-(k-1)})$. Let us outline how this rate is established—this will allow us to explain some key techniques used in p version approximation theory.

Suppose we consider a mesh of triangles, then for the p version, there will only be a finite number of elements, this number being independent
of \( p \), i.e. \( n \). Consider the triangular element \( \tau_j^{(n)} \) (with \( h_n, j < 1 \)). The first step is to extend the function \( \tilde{u} \) defined on \( \tau_j^{(n)} \) to an open unit square \( R \) satisfying \( R > \tau_j^{(n)} \), such that the norm-preserving property
\[
\|\tilde{u}\|_{H^k(R)} \leq C\|\tilde{u}\|_{H^k(\tau_j^{(n)})}
\]
holds.

Next, we use the Tchebysheff transformation, which transforms algebraic polynomials into periodic trigonometric polynomials. \( \tilde{u} \) gets transformed into \( \tilde{\tilde{u}} \) on the mapped square \( \tilde{R} \) such that the norm is once again preserved, i.e.
\[
\|\tilde{\tilde{u}}\|_{H^k(\tilde{R})} \leq C\|\tilde{u}\|_{H^k(R)}.
\]
The function \( \tilde{\tilde{u}} \), of course, is now periodic and therefore can be approximated by Fourier expansions. For each \( p = 1, 2, \ldots \) we therefore obtain a trigonometric polynomial \( \tilde{\omega}_p \) satisfying
\[
\|\tilde{\tilde{u}} - \tilde{\omega}_p\|_{H^s(\tilde{R})} \leq C_p^{-(k-s)}\|\tilde{\tilde{u}}\|_{H^k(\tilde{R})}, \quad 0 \leq s \leq k.
\]
Transforming the functions \( \tilde{\tilde{u}}, \tilde{\omega}_p \) back, we get a \( \tilde{\omega}_p^{(j)} \) for each \( p, j \) such that
\[
|\tilde{\omega}_p^{(j)}|_{H^s(\tau_j^{(n)})} \leq C_p^{-(k-s)}\|\tilde{\tilde{u}}\|_{H^k(\tau_j^{(n)})}.
\]
Note that in the above, we used the fact that \( \tau_j^{(n)} \) is a proper subset of \( R \), i.e. the distance of \( \tau_j^{(n)} \) from the boundary of \( R \) is strictly positive.

The above construction yields a separate polynomial for every triangle. The next step is to adjust these polynomials so that they will be continuous on \( \Omega \). As shown in [BS2], for the case \( k > 3/2 \), we may first adjust \( \tilde{\omega}_p^{(j)} \) so that it coincides at the vertices with the function \( \tilde{u} \) and estimate (4.4) is preserved. Then the jump (discontinuity) \( \tilde{\psi} \) on a side \( \gamma \) of \( \tau_j^{(n)} \) is a polynomial which vanishes at the ends of \( \gamma \) and satisfies (by (4.4) and the
trace theorem),

\[ (4.6) \quad \| \tilde{u} \|_{H^{1/2}(\gamma)} \leq C_p^{-(k-1)} \| u \|_{H^k(\tau_j)} \]

The final step is to extend \( \tilde{u} \) to a polynomial \( \tilde{z}^{(j)}_p \) over \( \tau_j^{(n)} \) which coincides with \( \tilde{u} \) on \( \gamma \) and vanishes on the other two sides of \( \gamma \) and which has a suitably bounded \( H^1(\tau_j^{(n)}) \) norm. In the original proof from [BSK], this was accomplished in a non-optimal way, giving an estimate

\[ (4.7) \quad \| \tilde{z}^{(j)}_p \|_{H^1(\tau_j^{(n)})} \leq C \| \tilde{u} \|_{H^1(\gamma)} \]

The right side of (4.7) was bounded above by \( C_p^{-(k-2)} \| u \|_{H^k(\tau_j)} \), so that subtracting \( \tilde{z}^{(j)}_p \) from \( \tilde{u}^{(j)}_p \) for each \( j \) gave a polynomial \( \tilde{w}^p \in V^n \) satisfying

\[ (4.8) \quad \| \tilde{u} - \tilde{w}^p \|_{H^1(\Omega)} \leq C_p^{-(k-2)} \| \tilde{u} \|_{H^k(\Omega)} \]

Hence, we obtained the desired estimate, except that a factor \( p^{-1} \) was lost. This loss can be reduced to \( p^{-c} \), for \( c > 0 \) arbitrary, by a useful interpolation argument, which we illustrate below.

We obviously have

\[ (4.9) \quad \inf_{\tilde{w} \in V^n} \| \tilde{u} - \tilde{w} \|_{H^1(\Omega)} = \| \tilde{u} - P_p \tilde{u} \|_{H^1(\Omega)} \leq C \| \tilde{u} \|_{H^1(\Omega)} \]

where \( P_p \tilde{u} \) is the \( H^1 \) projection of \( \tilde{u} \) onto \( V^n \). Also, from (4.8), we get

\[ (4.10) \quad \| \tilde{u} - P_p \tilde{u} \|_{H^1(\Omega)} \leq C_p^{-(k-2)} \| u \|_{H^k(\Omega)} \]
Now given $k$, we choose $K \geq 4$ and using (4.9), (4.10) and the interpolation theorem for Banach spaces, obtain

$$
\| \hat{u} - P_p \hat{u} \|_{H^1(\Omega)} \leq C \left( P^{-(K-2)} \right)^{\mu} \| \hat{u} \|_{(H^1(\Omega), H^K(\Omega))}
$$

where $\mu = \frac{k-1}{K-1}$. Noting that the norm on the right is just $\| u \|_{H^K(\Omega)}$, this gives

$$
\| \hat{u} - P_p \hat{u} \|_{H^1(\Omega)} \leq C P \left( \frac{K-2}{K-1} \right)^{-(k-1)} \| \hat{u} \|_{H^K(\Omega)}.
$$

Selecting $K$ large enough, we get $\frac{K-2}{K-1} = 1 - \frac{\epsilon}{k-1}$, which leads to the desired result, with the constant $C$ in (4.3) depending on $k$ and $\epsilon$.

Let us state the above interpolation argument as a theorem.

**Theorem 4.2.** Suppose $P_p : H^K(\Omega) \rightarrow V^n$ is a projection that satisfies

$$
\| u - P_p u \|_{H^K(\Omega)} \leq C \| u \|_{H^K(\Omega)}
$$

$$
\| u - P_p u \|_{H^s(\Omega)} \leq C(s) p^{-(s-\ell)} \| u \|_{H^s(\Omega)} \quad \text{for arbitrary } s > k, \ s \geq \ell.
$$

Then given $\epsilon > 0$, there exists $C = C(\epsilon, s)$ such that

$$
\| u - P_p u \|_{H^K(\Omega)} \leq C \left( P^{-(s-k)+\epsilon} \right) \| u \|_{H^s(\Omega)}
$$

The above interpolation argument has been used frequently in the $p$ version literature, e.g. in [BSK], [D1], [JV], [S1]. Note that such an argument does not work for the $h$ version, since the rate of convergence is bounded by the polynomial degree.

To obtain the optimal order (4.3), without the $\epsilon$, the extensions $z^{(j)}_p$ in (4.7) must be constructed in an optimal way. This was accomplished in
[BS1], using the theorems 4.3 and 4.4 below, which have found applications in other contexts as well (see [BS3], for example). Using these theorems together with (4.6) gives the desired optimal estimate for \( \frac{1}{p} \), leading to (4.3).

**Theorem 4.3.** [BS1], [BCM]. Let \( T \) be the standard triangle with sides \( \gamma_i \), \( i = 1, 2, 3 \), and \( f \) be continuous on \( \partial T \) with \( f|_{\gamma_i} \in \mathcal{P}_p(\gamma_i) \), \( i = 1, 2, 3 \). Then there exists \( U \in \mathcal{P}_p(T) \) such that \( U = f \) on \( \partial T \) and

\[
\|U\|_{H^1(T)} \leq C \|f\|_{H^{1/2}(\partial T)},
\]

where the constant \( C \) is independent of \( p \) and \( f \).

**Theorem 4.4.** [BS1]. Let \( S \) be the standard square with sides \( \gamma_i \), \( i = 1, 2, 3, 4 \). Let \( f \) be a continuous function on \( \partial S \) with \( f|_{\gamma_i} \in \mathcal{P}_p(\gamma_i) \), \( i = 1, 2, 3, 4 \). Then there exists \( U \in \mathcal{Q}_p(S) \) such that \( U = f \) on \( \partial S \) and

\[
\|U\|_{H^1(T)} \leq C \|f\|_{H^{1/2}(\partial S)},
\]

where the constant \( C \) is independent of \( p \) and \( f \). (See also [BCM]).

Theorems 4.3 and 4.4 are the analogs of the classical extension theorems. Here the extension is constrained to polynomials of the same degree as the traces. We mention that when using the space \( \mathcal{Q}'_p(S) \) instead of \( \mathcal{Q}_p(S) \), the constant \( C \) in Theorem 4.4 is mesh dependent ([Md]).

Let us now consider singular functions.

**Theorem 4.5.** [BS1]. Let \( \tilde{u} = r^\alpha \log^5 r \tilde{\psi}(\theta) \) where by \( r^\alpha \), \( \alpha \) complex, we denote \( \text{Re} r^\alpha \) or \( \text{Im} r^\alpha \), and let the meshes be quasiuniform with uniform \( p \). Then
\[
\inf_{\vec{v} \in V^n(\mathcal{G}^n,p)} \|\vec{u} - \vec{v}\|_{H^1(\Omega)} \leq C g(h,p,s) \min\left[ h^{\min(\hat{\alpha},p-\hat{\alpha})}, \frac{1}{p^{2\hat{\alpha}}} \right].
\]

where \( g(h,p,s) = \max(\log h^S, \log p^S) \), \( \hat{\alpha} = \text{Re} \alpha \) and \( C \) is a constant independent of \( h,p \).

See also [BD]. The next theorem is analogous to Theorem 2.5 and gives the rate of convergence of the \( h \) version with radical meshes.

**Theorem 4.6.** Let \( \vec{u} = r^{\alpha} \log r \psi(\theta) \). Then for the radical mesh with power \( \beta \) and uniform \( p \),

\[
\inf_{\vec{v} \in V^n} \|\vec{u} - \vec{v}\|_{H^1(\Omega)} \leq C(\hat{\alpha},\beta,s) h^p
\]

provided that \( 1 > \gamma = 1 - \frac{1}{\beta} > 1 - \frac{\hat{\alpha}}{p} > 0 \).

The proof is a simple generalization of the results in [BKP].

Finally, let us state the following theorem which shows that the \( h-p \) version with geometric meshes leads to exponential convergence.

**Theorem 4.7 [BG].** Let \( u \in B^2_\beta(\Omega), \ 0 < \beta < 1 \). Let \( \{\mathcal{G}^n\} \) be geometric meshes and \( \mu n \leq p \leq \nu n, \ p \geq 1, \ 0 < \mu, \nu < \infty \). Then we have

\[
(4.11) \quad \inf_{\vec{v} \in V^n} \|\vec{u} - \vec{v}\|_{H^1(\Omega)} \leq C e^{-\gamma \sqrt{n}}
\]

with \( \gamma > 0 \).

See also [BD]. The proof of Theorem 4.7 uses a combination of the approaches used in the outline of the proof of the convergence of the \( p \) version in this section and the proof in the one dimensional setting presented in Section 2 for the \( h-p \) version exponential rate.

It is more advantageous to use a linear degree vector than a constant
The use of the geometric mesh is advantageous when the effort of mesh generation is considered for complex geometries, since the refinement is restricted to a small subset. It is also possible to use a sequence of radical meshes as in one dimension and by means of the envelope obtain an exponential rate, but the mesh generation effort is now much more complex.

So far we mentioned only the problems where the approximation in the spaces $H^1(\Omega)$ was relevant. There are problems where the norms in $H^m(\Omega)$ have to be used. Various analogous results can be proven here. We refer to [Go] [S2] for details. An analysis of the error in the $L_2(\Omega)$ and negative norms has been presented in [JS]. Let us briefly describe a result of interest from this reference, related to the duality results presented in Section 2.

Consider the L-shaped domain in Figure 4.1. The solution of Laplace's equation on $\Omega$ is known to have a singularity behaving like $r^{2/3}$ by the results of the previous section. If the h version is used to approximate the solution of Laplace's equation (using a uniform mesh), the error will satisfy
The usual duality argument does not lead to a full extra power in $h$ for the $L_2(\Omega)$ norm over the $H^1(\Omega)$ norm. This is because of the non-convexity of the domain: the solution of the dual problem

$$-\Delta w = e \quad \text{on } \Omega$$

will satisfy

$$\|w\|_{H^k(\Omega)} \leq C\|e\|_{L_2(\Omega)}$$

with $k = 4/3$ only and not $k = 2$, as required. We note that the estimates in (4.12) cannot be improved, see [Wi].

For the $p$ version, however, we obtain the full extra power of convergence, i.e.

$$\|e\|_{H^1(\Omega)} \approx Cp^{-4/3}, \quad \|e\|_{L_2(\Omega)} \approx Cp^{-1}\|e\|_{H^1(\Omega)} \approx Cp^{-7/3}.$$  

This is because the solution of (4.13) can now be decomposed into a singular and smooth part as in equation (3.10). The smooth part will satisfy (4.14) with $k = 2$. The singular part will not; however, since it behaves like $r^{2/3}$, it is approximated at double the usual rate by the $p$ version, and the duality argument easily follows. The same argument works for other non-convex domains (see [JS]).

So far, we have assumed that $g = 0$ on $\Gamma_D$. In the inhomogeneous case, $u$ is approximated by first finding a suitable approximation $\tilde{g}_n$ to $g$ in the trace space $W^h = W^h(\Gamma_D) = \{w|_{\Gamma_D}, w \in V^h\}$ where $V^h$ is the finite element space when $\Gamma_D = \emptyset$. This introduces a new component, depending upon the difference $\tilde{g} - \tilde{g}_n$, into the error. There are various interesting
results on the effect of this component on the total error, depending upon how $g_n$ is defined. We refer to [BS3] for example.

Let us mention here that the spectral element method, which has many features in common with the p version, also leads to various approximation results. We will not survey them here but refer instead to [MdP] and the references therein.

In the three-dimensional case the solution is essentially similar but more complex. Let us mention the following theorem for the p version.

**Theorem 4.8** [D2]. Consider a fixed mesh $\mathcal{G}$ and a sequence of degrees $p_n \geq 0$. Let $f \in C^\infty(\overline{\Omega})$. If $u$ is the solution of (3.13) and $\omega_j, \beta_{\ell,i}, j = 1, \ldots, n_e, \ell = 1, \ldots, n_w, i = 1, 2, \ldots$, are as described in Section 3, then for arbitrary $c > 0$

\[
\inf_{v \in V_n} \| u - v \|_{H^1(\Omega)} \leq C p^{-2\mu+c}
\]

where

\[
\mu = \min\{\frac{\pi}{\omega_j}, \beta_{\ell,i}\}
\]

and $C$ depends on $u$ and $c$ but is independent of $p$.

The above theorem is based on the regularity results (3.14), (3.15) for the solution $u$.

Theorem 4.5 indicates what was observed computationally for the model problems we considered, namely that in most practical cases the edge singularity (and not the vertex one) governs the rate of convergence for the p version.

In three dimensions the proper mesh selection is complex because of the (anisotropic) singular behavior of the solution. Here, in the neighborhood of edges, long (needle) elements have to be used, because while the solution is
smooth along the edge, it is unsmooth in the direction perpendicular to the edge. In the neighborhood of the vertices, the singularity is isotropic and elements of the usual type are used. Hence, the class of meshes (radical and geometric) introduced in two dimensions has to be enlarged. In addition, the mesh generation has to be more sophisticated, to take into account the transition between various meshes. For the approximation properties of "needle" elements which violate the usual minimal angle condition, we refer to [BA], [Kr].

As we saw in Theorem 4.7, for a proper sequence of meshes and degree distributions, we have an exponential rate of convergence (with the h-p version). In the three dimensional case, we get a result which is quite analogous to Theorem 4.7, with the exponent \( \frac{3}{N} \) in (4.11) now being replaced by \( \frac{3}{N} \). For more see [BG5].

4.3. The effect of numerical quadrature

So far, we have assumed throughout that exact integration is used in the calculations of stiffness matrices and load vectors. In practice, of course, these are generally computed using numerical integration. Unlike the h version, where a rule of fixed precision may be used as \( h \to 0 \), in the p version, the precision of the rule must also increase as \( p \to \infty \). In [BnS], the optimal convergence of the p version using numerical integration in \( \mathbb{R}^t \), \( t = 1, 2, 3 \) was established under a set of sufficient conditions on the quadrature rule. Let us describe these results for the model problem (3.6), where \( \partial \Omega = \Gamma_D \), \( g = 0 \), a fixed mesh \( \mathcal{J}^n = \mathcal{J} \) is used and \( p_j = p \) in (4.2).

Suppose \( K \) is a reference element, \( K = S \) or \( T \). Then a quadrature rule on \( K \) induces a corresponding rule on each element \( \tau_j \). We then use a family of rules \( \{ R_j \} \) satisfying the following two assumptions.
(A) The weights are positive and the quadrature points lie within $K$.

(B) $R_p$ is exact for all $v \in \mathcal{U}_m(K)$ with $m \geq 2p$.

Condition (B) can be weakened, see [BnS].

Newton-Cotes rules are unsuitable for the $p$ version, since they have low degree of precision and violate the positivity of weights when $p$ is large. The most advantageous rules are Gaussian rules, which obviously satisfy (A).

In [BnS], it is shown that in $1 - d$, the $p$ point Gauss-Legendre and $p - 1$ point Gauss-Lobatto rules are the minimal required when polynomials of degree $p$ are used. For $Q_p$ type elements in higher dimensions, the corresponding tensor product rules using $p + 1$ points in each direction are minimally required. For other rules, see [BnS].

Let $A = \text{diag}(k_1)$, where $k_\ell = k_\ell(x)$ are the coefficients in $(3.6)$. Also, let $DF_j^{-1}(x)$ denote the Jacobian matrix of the inverse of $F_j$ and define the matrix $B_j = B_j(x)$ by

$$B_j = [b_{k\ell}^j] = (DF_j^{-1})A(DF_j^{-1})^T.$$ 

Moreover, let $\beta_{s,q} = \max_{j,k,\ell} \|b_{k\ell}^j\|_{W,s,q(K)}$. Then the following theorem is proved in [BnS].

**Theorem 4.9.** Let $\Omega \subset \mathbb{R}^t$. Let $f$ in $(3.3)$ be in $H^s(\Omega)$, $s > t/2$ and let $u$, the solution of the model problem $(3.6)$, lie in $H^k(\Omega)$, $k > 1$. Let $b_{k\ell}^j \in H^d(K)$ for each $j, k, \ell$, with $d > t/2$. Let $\tilde{e} = \tilde{e}_p$ be the error when the finite element solution $\tilde{u}_p$ is computed using rules $\{R_p\}$ which satisfy (A), (B). Then for any $0 \leq q < m-p$, $r = \min(p,m-p-q)$,

$$\|\tilde{e}\|_{H^1(\Omega)} \leq C\left((m-p)^{-\frac{s-t}{2}}\|f\|_{H^s(\Omega)} + q^{-\frac{d-t}{2}}\beta_{d,2}\|u\|_{H^1(\Omega)}\right)$$

$$54$$
The above estimate contains three terms. The first term, which is
\[ -(s^{-\frac{t}{2}}) \quad O((m-p)^2) \] arises from the approximation of \( f \), the second, which is
\[ -(d^{-\frac{t}{2}}) \quad O(q) \] arises from approximating the mappings \( F_j \) and coefficients \( k_i \),
and the final term, \( O(r^{-\frac{t}{2}}) \), arises from the solution \( u \).

For the case that \( f, F_j \) and \( k_i(x) \) are sufficiently smooth, the
above theorem shows that if \( m = 2p \), then we once again recover the
\( O(p^{-\frac{t}{2}}) \) convergence predicted by Theorem 4.1. A similar result holds for
the case that the solution is singular, for which it may be shown that the
optimal rate of \( O(p^{-2\alpha}) \) may be recovered.

If, however, \( d \) is small, then "overintegration" (i.e. taking \( m > 2p \))
may be needed to preserve the error bound. See [BnS], [Ki] for additional
theorems and various computational results.

4.4. Numerical examples

The theorems above are asymptotic in nature. To verify that they hold
for a practical range of parameters as well, we consider a computational
example involving the elasticity problem (3.7) over the L-shaped domain shown
in Figure 4.1.

We specify tractions on \( \partial \Omega \) (with \( \Gamma_D = \phi, \ \Gamma_T = 0 \) and sides OA, OE
being traction free) such that the true solution \( \mathbf{u} = (u_1, u_2) \) is given by

\[ u_1 = \frac{1}{2G} r^\alpha [(\kappa-Q(\alpha+1))\cos \alpha \theta - \alpha \cos(\alpha-2)\theta] \quad \text{(Mode 1)} \]
\[ u_2 = \frac{1}{2G} r^{\alpha} \left[(\kappa \cdot Q(\alpha+1))\sin \alpha \theta + \alpha \sin(\alpha - 2)\theta \right] \quad \text{(Mode 2)} \]

where \( \alpha = 0.5444837 \) and \( Q = 0.5430756 \). Here, \( G \) is the modulus of rigidity and \( \kappa = 3 - 4\nu \), with \( \nu = 0.3 \). Equation (4.16) shows that the solution has an \( r^{\alpha} \) type singularity at the origin 0, of the same type as in the decomposition (3.10) - (3.11). It may also be mentioned that the solution lies in \( H^{1+\alpha-c}(\Omega) \) for arbitrary \( c > 0 \) and in the Besov space \( B^{1+\alpha}_{2,\infty}(\Omega) \).

We first consider the FEM approximation of this problem using a uniform mesh with square elements, as shown in Figure 4.2.

Since we have a quasiuniform mesh, by Theorem 4.5, for \( p \geq 1 \) we obtain the estimate

\[ \|e\|_E \leq C \min \left[ h^{\alpha}, \frac{h^{\min(\alpha,p-\alpha)}}{p^{2\alpha}} \right]. \]

This leads to the predicted asymptotic rates \( h^{\alpha} \) and \( p^{-2\alpha} \) for the \( h \) and \( p \) versions, respectively.

![Figure 4.2. The uniform mesh.](image)
Figures 4.3 and 4.4 in which the relative error in the energy norm \( \| e \|_{E,R} \) is plotted respectively as a function of \( h \) and \( p \), show that the error is in the asymptotic range even for moderate \( h \) and \( p \). We reach the same conclusions as in the one-dimensional case.

![Figure 4.3. The error for the h version.](image)

![Figure 4.4. The error for the p version.](image)
Let us now consider the case of strongly refined (geometric) meshes. Such a mesh is shown in Figure 4.5 for $n = 2$ layers. This is a geometric mesh with ratio 0.15 (which leads to nearly optimal convergence).

Figure 4.6 shows the relative error vs. the number of degrees of freedom for various $n$ (with $p$ being increased in each case to increase $N$) in $\log ||e||_{E,R}$ vs $N^{1/3}$ scale. In Figure 4.7, $\log ||e||_{E,R}$ for various combinations of $n$ and $p$ has been plotted. It is observed that $||e||_{E,R} \propto Ce^{-\gamma N^{1/3}}$, which is consistent with Theorem 4.7.

Remark 4.1. The meshing required by the h-p version can be expensive to implement. Often (especially in the context of commercial p version codes), it is enough to select a fixed "good" mesh, i.e. one that is sufficiently refined near potential singularities, so that the error decreases largely in the pre-asymptotic exponential region. This region is clearly seen in Fig. 4.6 for any fixed $n$ — when $p$ is increased further, the error goes into the asymptotic region, which is algebraic. In general, for small required
accuracies, the mesh should be over-refined (low p), while for high accuracies, it is under-refined (high p). The optimal design of the mesh would require either a fully adaptive approach (see e.g. [DO]) or an approach based on an expert system (see e.g. [BR]).

We presented here only examples of two dimensional problems. For some computational results of three dimensional problems, we refer, for example, to [AB], [CM]. The basic features are completely similar to the two dimensional case.

![Figure 4.6. The error $|\text{e}|_{E,R}$ as a function of $n$ and $p$.](image)

![Figure 4.7. The error $|\text{e}|_E$ for selected combinations of $(n,p)$.](image)
4.5. The pollution problem

In Section 2, we discussed briefly the "pollution effect" in the p version. By "pollution", we mean the phenomenon by which a singularity in the solution at one point influences the behavior of the approximation over the entire element. With large elements, the accuracy can thereby be unfavorably affected even at points that are distant from the singularity. This effect is even more pronounced when it is the derivatives of the solution (like stresses) that are of interest. Some mathematical analysis of pollution may be found in [WI], however, a complete theory is not yet available.

Let us show the effect of this pollution by an example from [BO]. Consider the plane strain problem \( (E = 10^6, \nu = 0.3) \) of a plane loaded by a concentrated load as shown in Figure 4.8. The boundary conditions are traction conditions such that the solution is the well known Bousinesgue solution on a half plane. The solution then has infinite energy but it has finite energy on any subdomain, particularly in the shaded area shown in Figure 4.8.

In the case of the h version, the energy error in the shaded area will be essentially the same as the error of the best approximation on the domain (let us call it pollution free). This property of the h version is caused by the filtering out of the singularity through a layer of elements. The same effect of filtering occurs in the h-p version, as long as there are a sufficient number of layers of elements between the singularity and the shaded domain. In Figure 4.9, we show a sequence of four meshes (on half of the panel because of symmetry) used for the computation of the error in the shaded area. Figure 4.10 shows the relative error. We have also depicted the
pollution-free error which is the error of the best approximation of the solution on the shaded area. We see clearly the effect of the pollution, particularly when the p version is used on an insufficiently refined mesh, as well as the effect on this pollution when the layers are increased.

![Diagram of a panel loaded by a concentrated force.](image)

**Figure 4.8.** The scheme of a panel loaded by a concentrated force.

![Scheme of refined meshes.](image)

**Figure 4.9.** The scheme of refined meshes.
5. LOCKING AND ROBUSTNESS

Suppose we consider (3.7) when the Poisson ratio $\nu$ is very close to 0.5 and use the h version with piecewise linear elements to approximate this problem. Then it is well-known that the finite element solution for practical choices of the parameter $h$ is generally not very accurate. This is due to a phenomenon called locking (Poisson ratio locking in the particular case mentioned), which is said to occur when the accuracy of numerical schemes for the approximation of certain parameter-dependent problems deteriorates when the parameter is close to a limiting value. Another example of locking occurs when the h version is used on (3.6) when the parameter $d = k_1/k_2$ is close to 0 (highly anisotropic materials). Other examples of locking include various plate and shell models when the thickness $t$ is close to 0 (see, for example, [AF2], [Pt], [BrF]).
For such problems, one needs methods that are robust, i.e. which work equally well, essentially independent of the parameter value. In this connection, it has been demonstrated for various problems that the p version is often robust in situations where low-order standard h version schemes are not. For example, in [Vo], it has been shown that in the case of Poisson ratio locking, the p version is essentially free from locking for the displacements, while in [Li], the Timoshenko beam model (for which the standard h version shows locking, see [Ar]) was shown to be free from locking when the p and h-p versions are used. Analogous results have been shown for shell models, see [Pt].

Let us point out that other robust, locking-free methods exist. For example, in [SV], it was established that for the case of Poisson ratio locking, using the h version with high order elements \( p \geq 4 \) also avoids locking. Another strategy to overcome locking is to use mixed methods, as done for the Timoshenko beam problem in [Ar] and for various other problems in e.g. [AF2] [Pt] [BrF]. However, the reformulations of the variational form necessary for mixed methods may not be feasible in the context of available commercial codes, so that it may be more desirable to have a method (like the p version) which is robust in terms of the standard formulation of the problem.

In [BS4], we have presented a general theory of locking and robustness. Let us present some definitions and results from there, adapted to the context of the model problem (3.7), when the material is nearly incompressible.

Let us denote the bilinear form (3.7) by \( B_{\nu}(...,\cdot,\cdot) \) and let the exact solution be \( \hat{u}_n^{\nu} \) when the Poisson ratio is \( \nu \) \( (0 \leq \nu < 0.5) \). Then the corresponding finite element solution \( \hat{u}_n^{\nu} \) will satisfy
Let us denote the energy norm \( (B_{\nu}(\tilde{w},\tilde{v}))^{1/2} \) by \( \|\tilde{w}\|_{E,\nu} \). We will restrict our attention to the case when the exact solutions \( \tilde{u}^\nu \) belong to the closed unit balls \( H^k_{B,\nu} \subset H^k(\Omega) \) (\( k \geq 2 \)) defined in terms of the norms

\[
\|\tilde{w}\|^2_{k,\nu} = \|\tilde{w}\|^2_{H^k(\Omega)} + (1-2\nu)^{-2}\|\text{div}\;\tilde{w}\|^2_{H^{k-1}(\Omega)}.
\]

As shown in [BS5], \( H^k_{B,\nu} \) correspond to precisely the correct weighted spaces that characterize the regularity of the solution when the data is appropriately bounded.

We assume that the sequence \( \{V^n\} \) is such that

\[
C_1F_0(N) \leq \sup_{\tilde{w} \in H^k_{B,\nu}} \inf_{\tilde{v} \in V^n} \|\tilde{w} - \tilde{v}\|_{1,\Omega} \leq C_2F_0(N)
\]

where \( H^k_{B,\nu} \) is the unit ball in the \( H^k(\Omega) \) norm. Here, \( C_1, C_2 \) are independent of \( N \) and \( F_0(N) \to 0 \) as \( N \to \infty \). Equations (5.1), (5.2) imply that for any \( 0 < \nu_0 < 0.5 \), we can find constants \( C_1(\nu_0), C_2(\nu_0) \), such that the following holds uniformly for all \( 0 \leq \nu < \nu_0 \) (\( N = N_n \)),

\[
C_1(\nu_0)F_0(N) \leq \sup_{\tilde{u}^\nu \in H^k_{B,\nu}} \|\tilde{u}^\nu - \tilde{u}^\nu_n\|_E \leq C_2(\nu_0)F_0(N).
\]

If (5.3) can be made to hold uniformly for all \( 0 \leq \nu < 0.5 \), we say that \( \{V^n\} \) is free from locking. More precisely, we have the following definitions.

**Definition 5.1.** Let \( f \) be an increasing function with \( \lim_{N \to \infty} f(N) = \infty \). The extension process \( \{V^n\} \) shows locking of order \( f(N) \) for the family of problems (5.1), \( \nu \in [0,0.5) \), with respect to the solution sets \( \{H^k_{B,\nu}\} \) in
the energy norm iff

$$0 < \lim_{N \to \infty} \sup_{\nu \in [0, 0.5]} \left( \sup_{u \in H^1_{\text{div}, \nu}} \left( \sup_{u' \in H^1_{\text{div}, \nu}} \| u' - u \|_{L^2(E)} \right) \right) (F_0(N)f(N))^{-1} - M < \infty.$$  \hspace{1cm} (5.4)

For the case that $M$ (depending on $f$) is bounded (respectively, infinite), we say that the order of locking is at most (respectively, at least) $f(N)$. If (5.4) holds with $f(N) \equiv 1$, we say that $\{V^n\}$ is free from locking.

**Definition 5.2.** Let $g$ be a decreasing function, $g(N) \to 0$ as $N \to \infty$. The extension process $\{V^n\}$ is robust with uniform order $g(N)$ for the family of problems (5.1), $\nu \in [0, 0.5)$, with respect to the solution sets $\{H^1_{\text{div}, \nu}\}$ in the energy norm iff

$$\lim_{N \to \infty} \sup_{\nu \in [0, 0.5]} \left( \sup_{u \in H^1_{\text{div}, \nu}} \left( \sup_{u' \in H^1_{\text{div}, \nu}} \| u' - u \|_{L^2(E)} \right) \right) (g(N))^{-1} = K < \infty.$$  \hspace{1cm} (5.5)

We see from the above definitions that if $f(N)$ is such that $f(N)F_0(N) = g(N) \to 0$ as $N \to \infty$, then $\{V^n\}$ shows locking or order $f(N)$ iff it is robust with maximum uniform order $g(N)$. Note that in the above definitions, we could choose other error measures instead of the energy norm (e.g. the $H^1(\Omega)$ norm), see [BS4].

As $\nu \to 0.5$, we see that the incompressibility constraint,

$$Cw = \nabla \cdot \vec{w} = 0$$  \hspace{1cm} (5.6)

gets imposed on the exact solution. The reason locking occurs is that (5.6) gets imposed on the finite element solution as well. For there to be no locking in the energy or $H^1(\Omega)$ norm, the following necessary condition must be satisfied.
Condition (5.7) is sufficient as well, under certain general conditions (see [BS4]) which are shown in [BS5] to be satisfied for our model problem. This reduces the whole question of locking to one of approximability alone. Let us give some sample results from [BS5] in this connection.

Let us take $\Omega$ to be a square domain. Let $\mathcal{J}_1^h$ and $\mathcal{J}_2^h$ be the uniform triangular and uniform square meshes, respectively (Figure 5.1). Define, for $i = 1, 2, 3,$

$$V_{p,h}^i = \{ u \in C^0(\Omega), \ u|_S \in R^i_p(S), \ \forall \ S \in \mathcal{J}_i^h \}$$

where $R^1_p(S) = \mathcal{P}_p(S)$, $R^2_p(S) = \mathcal{Q}_p(S)$, $R^3_p(S) = \mathcal{Q}_p'(S)$ and $\mathcal{J}_3^h = \mathcal{J}_2^h$.

![Figure 5.1. Uniform triangular and square mesh.](image-url)
Theorem 5.1. Let the extension process be the h version for the problem (3.7), (5.1) using subspaces $V^n = (V_{p,h}^i)^2$ with $p$ fixed. Let the solution sets be $H_{k,v}^P$ and let $k \geq p + 1$. Then the following is true for locking in the energy norm and also in the $H^1$ norm.

<table>
<thead>
<tr>
<th>Space $V^n$</th>
<th>Degree $p$</th>
<th>Order of locking $\ell$ $f(N) \approx C N^\ell$</th>
<th>Robustness order $r$ $g(N) \approx C N^{-r}$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$(V_{p,h}^1)^2$</td>
<td>$1 \leq p \leq 3$</td>
<td>$\ell = 1/2$</td>
<td>$r = (p-1)/2$</td>
</tr>
<tr>
<td></td>
<td>$p \geq 4$</td>
<td>$\ell = 0$</td>
<td>$r = p/2$</td>
</tr>
<tr>
<td>$(V_{p,h}^2)^2$</td>
<td>$p \geq 1$</td>
<td>$\ell = 1/2$</td>
<td>$r = (p-1)/2$</td>
</tr>
<tr>
<td>$(V_{p,h}^3)^2$</td>
<td>$1 \leq p \leq 2$</td>
<td>$\ell = 1/2$</td>
<td>$r = (p-1)/2$</td>
</tr>
<tr>
<td></td>
<td>$p \geq 3$</td>
<td>$\ell = 1$</td>
<td>$r = (p-2)/2$</td>
</tr>
</tbody>
</table>

The above shows that locking cannot be avoided for $p \leq 3$ for the h version. Also, using rectangular elements leads to locking for all $p$. The result for $p \geq 4$ for $(V_{p,h}^1)^2$ was first proved in [SV] and holds for more general meshes.

In contrast, as mentioned earlier, it was shown in [Vo] that the $p$ version (using straight sided triangles) leads to uniform robustness of optimal order for the displacements, so that locking is eliminated. However, the use of curvilinear elements is indispensable for the $p$ version in most cases. In this connection, we have the following computational results.

Let us consider the $p$ version for a single curvilinear element of type $Q_p^r$ for problem (3.7). We assume that this element is the image of the standard square under a curvilinear mapping $F$. Below, we show the relative errors in the energy norm for the cases that the mapped element is a square, a
trapezoid, a quadrilateral with one side a parabola, and a quadrilateral with one side a circle. We do this for $v = 0.3$ (Figure 5.2) and $v = 0.5 - 10^{-10}$ (Figure 5.3). It is observed that the rate of convergence now behaves like $(p-s)^{(k-1)}$, with $s \geq 0$ depending on $F$ ($s = 0$ for affine $F$). There is no change in the apparent asymptotic rate so that in this sense, there is no observed locking.

We have treated various other examples of locking in [BS4] using the above definitions and principles. Like the above case, the reason locking occurs in general is due to the imposition of a constraint similar to (5.6). Such constraints usually can be interpreted as inter-element continuity requirements, with (5.7) breaking down when the constrained finite element space is too small. (For example, (5.6) can be interpreted as imposing $C^{(1)}$ continuity on a space related to $V^1$, see [BS5].) Unlike low-order h version schemes, the p version remains robust under such continuity constraints, which is why it is more resistant to locking.

As another example, consider problem (3.6) with $k_1/k_2$ close to zero. In this case, the constraint is $\frac{\partial w}{\partial x_2} \neq 0$. Suppose the h version is used, with piecewise polynomials of degree $\leq p$ on either of the meshes in Figure 5.1.

![Fig. 5.2. Error behavior for curved elements for $v = 0.3$.](image-url)
Then, since the meshes are aligned with the $x_1 - x_2$ directions, there is no locking (see [BS4]) since there are a number of functions satisfying $\frac{\partial v}{\partial x_2} = 0$.

However, if these meshes are not aligned with $x_1 - x_2$ direction (for example, as in Figure 5.4), then the only functions satisfying $\frac{\partial v}{\partial x_2} = 0$ will be actual (rather than piecewise) polynomials and (5.7) will not hold. There
will be complete locking in this case, with the method being non-robust. On
the other hand, if we use the p version with just one element, then the
constraint $\frac{\partial v}{\partial x_2} = 0$ has no effect and there is no locking, i.e. the method is
robust. See [BS4] for a detailed analysis. Also, see [SBS] for an analysis
of locking in the Reissner-Mindlin plate.

6. EXTRACTION AND POST-PROCESSING TECHNIQUES

Usually in computational practice, the solution $u$ of the differential
equation considered is only a tool to get the primary quantity of interest.
For example, the goal of the computation may be to find the stresses at a
point, the stress intensity factors, or the resultant (e.g. the moments) in
shell and plate theory.

The results in Sections 2, 4 indicate the high asymptotic rates of
convergence, including exponential convergence, that can be obtained by the
h-p version, when the error is measured in the energy norm. The question that
we ask here is how these rates can be obtained for quantities of interest
other than the energy. For instance, if high-order derivatives of the h-p
solution were computed directly, then severe oscillations could result and
these rates would be lost.

Mathematically, we are interested in evaluating the value of a certain
function of the solution. There are various forms in which this function may
be expressed, using properties of the equation under consideration. Choosing
one such form, we then substitute the finite element solution instead of the
exact one to get an approximation of the desired value. If the exact solution
were substituted, then the computed value would be the same, no matter which
form of the function were used. However, when the finite element method is
used, the accuracy of the computed value can be very different. In this section, we indicate how the convergence rate in the energy can be achieved for other quantities as well.

We illustrate the main idea by considering the simple one-dimensional problem,

\[-u'' + u = f \quad \text{on } \Omega = (0,1) = I\]

\[u(0) = u(1) = 0\]

where \( f \) is assumed to be in \( L^2(\Omega) \) so that \( u \in H^2(\Omega) \). Suppose the quantity we are interested in is the value of \( u'(0) \).

Problem (6.1) can be cast in the variational form

\[(6.2) \quad B(u,v) = \int_0^1 (u'v' + uv) \, dx = \int_0^1 fv \, dx\]

with \( u \in H^1_D(\Omega) = H^1(\Omega) \), \( v \in H^1_D(\Omega) \).

Let \( \psi \in H^1(\Omega) \) satisfy \( \psi(0) = 1, \, \psi(1) = 0 \). Then

\[B(u,\psi) = \int_0^1 (-u'' + u)\psi \, dx - u'(0) = \int_0^1 f\psi \, dx - u'(0)\]

so that

\[(6.3) \quad u'(0) = \int_0^1 f\psi \, dx - B(u,\psi).\]

Hence, we have two possible ways of obtaining the desired value \( u'(0) \): we could either use differentiation, or equation (6.3). Function \( \psi \) is called an extraction function. Let us compare the accuracy of these two approaches.

Suppose that \( \{V^n\} \) is the sequence of subspaces defining the extension process \( (V^n \subset H^1_D) \) and \( \{u^n\} \) are the corresponding finite element
solutions. Define

\[(6.4) \quad F = F(u) = \int_0^1 f_\psi dx - B(u, \psi) (= u'(0))\]

\[(6.5) \quad F_n = F(u_n) = \int_0^1 f_\psi dx - B(u_n, \psi).\]

Then

\[(6.6) \quad |F - F_n| = B(u_n - u, \psi).\]

Now let \( z \in H^1_D(\Omega) \) be such that

\[(6.7) \quad B(v, z) = B(v, \psi), \quad \forall \; v \in H^1_D(\Omega).\]

Then for any \( \omega^{(n)} \in V^n \) we have

\[|F - F_n| = |B(u_n - u, z)| = B(u_n - u, z - \omega^{(n)}) \leq C\|u_n - u\|_{H^1(\Omega)} \|z - \omega^{(n)}\|_{H^1(\Omega)}.\]

Hence,

\[|F - F_n| \leq C\|u_n - u\|_{H^1(\Omega)} \inf_{\omega \in V^n} \|z - \omega\|_{H^1(\Omega)}\]

and the theory of Section 2 is applicable.

The smoothness of the function \( z \) can be easily established. Usually,

\[\inf_{\omega \in V^n} \|z - \omega\|_{H^1(\Omega)} \approx \|u_n - u\|_{H^1(\Omega)}.\]

We see therefore that the error \( |F - F_n| \) is, in general, of the order \( \|u_n - u\|_{H^1(\Omega)}^{2} \), i.e. the order of the error in the energy. In contrast, if we defined \( F_n \) by direct differentiation, i.e. \( F_n = u'_n(0) \), then we would only get \( \|u_n - u\|_{H^1(\Omega)} \) (modulo a logarithmic term) as the rate of convergence of \( |F - F_n| \).

We therefore see that by methods of the type above, called extraction methods, more accurate values can be obtained when \( \psi \) is properly selected.
(ψ smooth). (In the case above, we could, for instance, take ψ = 1 - x). Let us note that the extracted value is not changed when we replace ψ by ψ - ω, ω ∈ V^N, so that, for example, in the h version, ψ can be restricted only to the element containing x = 0. (See [BM] for more)

Similar ideas can be used for other values of interest. One important case is the extraction of stress intensity factors (see Section 3). Let us elaborate briefly on this case. Assume that we are solving the problem

-Δu = 0 on Ω

where Ω is the L-shaped domain shown in Figure 4.1. Suppose on ∂Ω we have the boundary condition

\[ \frac{∂u}{∂n} = g \text{ on } ABCDE \]
\[ u = 0 \text{ on } OE, OA. \]

Then using (3.10)', we have

\[ u = Fr^{2/3} \cos(\frac{2θ}{3}) + v \]

with

\[ |v| |r^{-2/3}|, \quad |\text{grad } v| |r^{1/3}| = o(1). \]

Here, F is the stress intensity factor. Let us define

\[ ψ = (r^{-2/3} \cos(\frac{2θ}{3}))χ(r) \]

where \( χ = 1 \) for \( 0 ≤ r < a, χ = 0 \) for \( r > b, a < b < 1, \ χ \in C^∞(Ω) \).

Obviously, \( w = Δψ \in C^∞(Ω) \) and

\[ w = 0 \text{ for } r < a \text{ and } r > b. \]

Now denoting \( Ω_ε = Ω - \bar{R}_ε \) where \( \bar{R}_ε = \left\{(r,θ) \mid r ≤ ε \right\} \) we have
\[\int_{\Omega} u \Delta \psi dx = \int_{\Omega} u w dx = \int_{\partial \Omega} \frac{\partial \psi}{\partial n} ud s - \int_{\partial \Omega} \frac{\partial u}{\partial n} \psi ds = - \frac{4}{3} F \int_{0}^{\pi/2} \cos^2 \frac{2\theta}{3} d\theta + o(1).\]

Hence, taking the limit as \( c \to 0 \),

\[F = -\left( \int_{\Omega} u w dx \right) \frac{1}{\pi}.\]

Proceeding as before, we define

\[F_n = -\left( \int_{\Omega} u_n w dx \right) \frac{1}{\pi}\]

and get

\[|F - F_n| \leq C\|u - u_n\|_{H^1(\Omega)} \|z - z_n\|_{H^1(\Omega)}\]

where \( z \) is the solution of the problem

\[-\Delta z = w\]

\[\frac{\partial z}{\partial n} = 0 \quad \text{on } ABCDE\]

\[z = 0 \quad \text{on } OE, OA.\]

Because \( w \in C^\infty(\Omega) \), and all internal angles except at 0 are \( \frac{\pi}{2} \), \( z \) is \( C^\infty \) except in a neighborhood of the origin. In this case (when \( F \neq 0 \)) we have

\[\|z - z_n\|_{H^1(\Omega)} \leq C\|u - u_n\|_{H^1(\Omega)}\]

and hence

\[|F - F_n| \approx \|u - u_n\|_{H^1(\Omega)}^2\]

(or better), i.e. \( |F - F_n| \) is once more of the order of the error in the energy. This is typical for this kind of extraction.

The stress intensity factor extraction method obviously requires a
priori knowledge of the form of the singularities. In two dimensions, a numerical algorithm for obtaining the singular form for the general elasticity problem may be found in [Pa] (see also [Sa]).

As a numerical example, we consider the L-shaped domain shown in Fig. 4.1. We consider the elasticity problem described in Section 4 with prescribed traction conditions on \( \partial \Omega \) such that the exact solution is given by

\[
u = u_1 + 2u_2
\]

where \( u_1, u_2 \) are given by (4.16). The mesh shown in Fig. 4.5 with two layers is used. In Table 6.1, we report results using the p-extension for the given mesh, showing the number of degrees of freedom \( N \), the normalized strain

<table>
<thead>
<tr>
<th>P</th>
<th>N</th>
<th>STRAIN ENERGY</th>
<th>STRESS INTENSITY FACTOR-MODE 1</th>
<th>STRESS INTENSITY FACTOR-MODE 2</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>41</td>
<td>6.42072796</td>
<td>0.95268</td>
<td>2.29075</td>
</tr>
<tr>
<td>2</td>
<td>119</td>
<td>6.74137580</td>
<td>1.02177</td>
<td>2.08422</td>
</tr>
<tr>
<td>3</td>
<td>209</td>
<td>6.77029847</td>
<td>1.00250</td>
<td>2.02239</td>
</tr>
<tr>
<td>4</td>
<td>335</td>
<td>6.77575144</td>
<td>1.00073</td>
<td>2.00437</td>
</tr>
<tr>
<td>5</td>
<td>497</td>
<td>6.77683967</td>
<td>0.99991</td>
<td>2.00097</td>
</tr>
<tr>
<td>6</td>
<td>695</td>
<td>6.77719530</td>
<td>0.99985</td>
<td>2.00022</td>
</tr>
<tr>
<td>7</td>
<td>929</td>
<td>6.77736281</td>
<td>0.99987</td>
<td>2.00005</td>
</tr>
<tr>
<td>8</td>
<td>1199</td>
<td>6.77749228</td>
<td>0.99990</td>
<td>2.00001</td>
</tr>
<tr>
<td>(\infty)</td>
<td>(\infty)</td>
<td>6.77776914</td>
<td>1.00000</td>
<td>2.00000</td>
</tr>
</tbody>
</table>

well as the (normalized) values of the first two stress intensity factors. In addition, we present the exact values \((p = \infty)\) of these stress intensity

75
factors. In Fig. 6.1, we show the relative error in the stress intensity factors as well as the relative error in the strain energy of the solution (i.e. square of the error in the energy norm). We see that the errors in the stress intensity factors are in fact of the order of those in the strain energy. Note the typical form in the first phase, where the rate is exponential, and in the second phase, where it is algebraic. For more, we refer to [SzB]. For the extraction of stress intensity factors in 3-d, we refer to [AB], [BPA].

![Fig. 6.1. Error of the stress intensity factors and strain energy.](image-url)
7. PLATE MODELLING

A very common problem in structural mechanics is the problem of solving partial differential equations on a thin domain. A typical problem of this type is the plate problem, which we elaborate on below.

Let \( \Omega = \{ (x = x_1, x_2, x_3) \in \mathbb{R}^3, (x_1, x_2) \in \omega \subset \mathbb{R}^2, |x_3| < \frac{d}{2} \} \) where \( \omega \) is a bounded domain with piecewise smooth boundary. Denote \( S = \{ x \in \mathbb{R}^3 \mid (x_1, x_2) \in \partial \omega, |x_3| < \frac{d}{2} \} \) and \( R_\pm = \{ x \in \mathbb{R}^3 \mid (x_1, x_2) \in \omega, x_3 = \pm \frac{d}{2} \} \). We assume that \( d \ll \text{diam}(\omega) \).

We are interested in the three-dimensional elasticity problem which we mentioned in Section 3. Denoting

\[
\bar{u} = (u_1, u_2, u_3), H_D^1(\Omega) = \{ \bar{u} \in (H^1(\Omega))^3 \mid u_3 = 0 \text{ on } S \}
\]

and

\[
F(\vec{v}) = \frac{1}{2} \int_\omega g(x_1, x_2)(v_3(x_1, x_2, -\frac{d}{2}) + v_3(x_1, x_2, \frac{d}{2}))dx_1dx_2,
\]

the solution \( \bar{u} \in H_D^1(\Omega) \) of the plate problem satisfies (3.4) for every \( \vec{v} \in H_D^1(\Omega) \), with \( B(\bar{u}, \vec{v}) \) given by (3.7) (with \( t = 3 \)). This problem is the so-called plate problem with soft simple support. Using other constrained spaces, we can describe other physical conditions (for example, using the condition \( \bar{u} = 0 \) instead of \( u_3 = 0 \) on \( S \) gives the clamped boundary condition).

In order to numerically approximate this problem, a common technique is to first replace the three-dimensional model by a two-dimensional one, which leads to a considerable reduction in the number of degrees of freedom and makes the problem more numerically tractable. The simplest model that could be considered is that of the Kirchoff plate. However, this has the disadvantage of requiring \( C(1) \) elements.
Two-dimensional models that only require $C(0)$ continuous elements can often be viewed as one of a hierarchy of models in the sense described below. To obtain such models, the solution is essentially "expanded" in terms of polynomials (or other functions) in the $x_3$ direction. More precisely, let $n = (n_1, n_2, n_3)$, $n_i \geq 0$ integer and define

$$\begin{align*}
    n^1_D(\Omega) &= \{ \vec{u} \in H^1_D(\Omega) \mid u_1 &= \sum_{j=0}^{n_1} \psi_{i,j}(x_1, x_2) L_j \left( \frac{2x_3}{d} \right), i = 1, 2, \\
    u_3 &= \sum_{j=0}^{n_3} \psi_{j}(x_1, x_2) L_j \left( \frac{2x_3}{d} \right) \} \\
\end{align*}$$

where $L_j(x)$ are Legendre polynomials. Then by the hierarchy of plate models we will understand the problem of finding $\vec{u}^{(n)}$ such that

$$B(\vec{u}^{(n)}, \nu) = F(\vec{\nu}) \quad \forall \vec{\nu} \in n^1_D(\Omega).$$

Let us note that because of symmetry, the models $(2n_1+1, 2n_2+1, 2n_3)$ are the same as the models $(2n_1+2, 2n_2+2, 2n_3+1)$. Typical models used are $(n, n, n+1)$ models. As we note below, such hierarchies of models are easily implemented in the framework of the p version. The exact plate solution and the solution of the particular model depends on the thickness $d$. It is possible to show that the models $n = (n_1, n_2, n_3)$, $n_1 \geq 1$, $n_2 \geq 1$ and $n_3 \geq 2$ converge to a limiting solution as $d \to 0$ which is the same as the limiting solution as $d \to 0$ of the exact (three dimensional) solution introduced above. The model $(1, 1, 0)$ which is very often used in engineering (the so-called Reissner-Mindlin model) does not have the above property, namely, the same limiting solutions as the exact one, unless modified elastic constants are used (see, e.g. [BL1]). We note that for $\nu = 0$, no modification is needed. Further, it
is possible to show that the polynomial choice in \( x_3 \) used in (7.2) is asymptotically optimal.

For fixed \( d \), the solutions of the models converge to the exact 3-d solution as \( \min(n_1,n_2,n_3) \to \infty \). If \( g(x) \) in (7.1) satisfies some smoothness conditions, then, for larger \( n \), \( n \) fixed, the rate with respect to \( d \) \( (d \to 0) \) is higher. The various models have different properties with respect to boundary layers, corner singularities, etc., see [BL2].

Assume now that the domain \( \omega \) is partitioned into the elements \( T_j \) based on triangles \( T \) or squares \( S \) (see Section 4). Then \( T_j = T_j \times (-d/2,d/2) \) is a three-dimensional element. Now denote \( Q_j(p,q) \) to be the set of polynomials of degree \( q \) in \( x_3 \) and of degree \( p \) in \( x_1, x_2 \) (this could include the case of \( P_0, Q_0 \) or \( Q_p \) type elements in \( x_1, x_2 \) — see Section 4). Then the finite element solution with elements \( Q_j(p,q) \) of the three-dimensional plate problem with degree \( q \) fixed can be interpreted as the finite element solution (using elements of degree \( p \)) of the plate model with \( n_1 = n_2 = n_3 = q \) in (7.2). Hence, the hierarchic modelling can be easily implemented in the frame of the \( p \) (or \( h-p \)) version of three-dimensional elasticity. (In fact, this implementation is available in the commercial program MSC/PROBE). Let us comment on three main aspects of the solution of the plate problem in this hierarchical framework.

a) **The locking problem.** Here, the theory addressed in Section 5 fully applies. The analog to constraint (5.6) for the \( (1,1,0) \) model is

\[
\varphi_{1,1} = -\frac{\partial \psi_0}{\partial x_i}, \quad i = 1,2.
\]

For more, see [SBS].

b) **The problem of optimal meshes.** In the neighborhood of corners,
refined meshes similar to the ones discussed in the previous sections (radical or geometric) should be used. In the neighborhood of edges, the use of "needle" elements analogous to those used in the neighborhood of edges in three dimensions is of importance to resolve boundary layers. Let us mention, however, that the singularity here is of different character than the edge singularity in three-dimensional problems. The problem of the optimal mesh has not been satisfactory resolved as yet.

c) The pollution problem. The error in the neighborhood of the boundary can now influence the accuracy of the solution far from the boundary. Essentially, there is pollution present.

In addition to the problem of plates discussed above, let us mention that a similar treatment is possible for other problems over "thin" domains, like those arising, for instance, in shell theory.

8. IMPLEMENTATIONAL ASPECTS

In this paper, our emphasis has been mainly on theoretical aspects of the p and h-p versions. In this section, we briefly focus on implementational aspects which, of course, are of essential importance in terms of practical use.

As mentioned in the introduction, currently there are only a few large research and commercial codes based on the p and h-p versions. STRIPE (Aeronautical Research Institute of Sweden) is a 3-d adaptive p version research code, Applied Structure (Rasna Corp. CA, USA) and MSC/PROBE (MacNeal Schwendler, CA, USA) are p version commercial codes and PHLEX (Computational Mechanics, TX, USA) is an adaptive h-p version commercial program. The implementation of such p and h-p version codes is significantly different from
that of h version codes. Let us summarize some of the basic differences.

a) **Mesh construction.** This is much simpler for the p version, which uses a relatively small number of elements of large size. With p version codes, it is important that the user select a good mesh to obtain a high rate of convergence (See Remark 4.1).

b) **Sparsity.** The local stiffness matrices are much less sparse for the p version than the h version. Hence, the computation of stiffness matrices has to be given special care, for instance by taking into account the computer architecture used. For example, the program STRIPE successfully exploits the architecture of the CRAY computer to deal with this problem. Reduced sparsity, especially in the case of three-dimensional problems with more than 100,000 - 500,000 unknowns, has serious implications for the I/O time of communication between disk and fast memory.

c) **Adaptivity.** In the case of the p version, the structure of an adaptive program is much simpler than the h version, since only new shape functions are added, without changing the mesh. The architecture for the adaptive h-p version (e.g. PHLEX) is much more complex (see [DO]). The adaptive selection of shape functions is usually determined by considerations based on the error in the energy. Essentially, those shape functions are added which maximally change the computed energy. The reason that adaptive principles for the h-p version are more complex is that simultaneous decisions about mesh changes and shape function selection have to be made.

d) **Accuracy assessment.** The p version allows the user to control the desired accuracy of the data of interest. This control is essentially based on comparison of the computed data by increasing the degree $p$, so that the user can a priori specify the required accuracy. This feature has been
incorporated e.g. in Applied Structure—see also [AB].

e) **Linear solver.** This can either be direct, when many right hand sides (loads) are present, or iterative, when the number of unknowns is large (> 100k, say) (see e.g. [Mn]). An effective preconditioner can be based on the low p (p = 1,2) discretization.

f) **Hierarchic bases.** As mentioned in Section 2.2, the shape functions for the p version have a hierarchic character. This is exploited in adaptive computations and in the construction of preconditioners.

g) **Overall effectivity.** The effectivity of any methods, h, p or h-p, depends upon many factors. In this paper, we have focused primarily on one factor, the asymptotic rate of convergence (with respect to the number of degrees of freedom). Some other important factors are computational cost, man-power cost for data preparation, and proper usage. Further, estimations of efficiency must also take into account various features that may be present, such as adaptive features, a posteriori error estimation, etc. (The capability of a posteriori error estimations is essential if a program is to give reliable results.)

Let us mention with respect to overall effectivity that the complex question of comparison of different methods and codes is usually accomplished through various benchmark problems. In this regard, an interesting engineering comparison study is presented in [CM], where the p version is compared to a commercial h version code for two test problems of typical industrial complexity. The results presented underscore the viability of the p version not only in comparisons based on the number of degrees of freedom, but also in those based on the total CPU time.
REFERENCES


83


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