## University of Alberta

## WAVELETS AND APPLICATIONS

by
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in
Mathematics

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## Chapter 1

## Wavelet Bases of Hermite Cubic Splines on an Interval

### 1.1 Introduction

In this chapter we shall construct wavelet bases of Hermite cubic splines on the interval. These wavelet bases are suitable for numerical solutions of differential equations.

By $L_{2}(\mathbb{R})$ we denote the linear space of all square-integrable real-valued functions on $\mathbb{R}$. The inner product in $L_{2}(\mathbb{R})$ is defined as

$$
\langle u, v\rangle:=\int_{\mathbb{R}} u(x) v(x) d x, \quad u, v \in L_{2}(\mathbb{R})
$$

If $\langle u, v\rangle=0$, then we say that $u$ and $v$ are orthogonal. The norm of a function $f$ in $L_{2}(\mathbb{R})$ is given by $\|f\|_{2}:=\sqrt{\langle f, f\rangle}$.

Smooth orthogonal wavelets with compact support were constructed by Daubechies (see [22]). The Daubechies orthogonal wavelets were adapted to the interval $[0,1]$ by Cohen, Daubechies, and Vial ([17]). Semi-orthogonal spline wavelets were constructed by Chui and Wang [16]. These spline wavelets were adapted to the interval [0,1] by Chui and Quak [15]. In [50] Wang constructed cubic spline wavelet bases for Sobolev spaces.

Orthogonal multi-wavelets were constructed by Donovan, Geronimo, Hardin,
and Massopust [23]. In [30], Heil, Strang, and Strela considered the possibility of construction of wavelets on the basis of Hermite cubic splines.

Let $\phi_{1}$ and $\phi_{2}$ be the cubic splines given by

$$
\phi_{1}(x):=\left\{\begin{array}{lll}
(x+1)^{2}(1-2 x) & \text { for } & x \in[-1,0] \\
(1-x)^{2}(2 x+1) & \text { for } & x \in[0,1] \\
0 & \text { for } & x \in \mathbb{R} \backslash[-1,1]
\end{array}\right.
$$

and

$$
\phi_{2}(x):=\left\{\begin{array}{lll}
x(x+1)^{2} & \text { for } & x \in[-1,0] \\
x(x-1)^{2} & \text { for } & x \in[0,1] \\
0 & \text { for } & x \in \mathbb{R} \backslash[-1,1]
\end{array}\right.
$$

In [20], Dahmen, Han, Jia, and Kunoth constructed biorthogonal multiwavelets on the basis of the Hermite cubic splines $\phi_{1}$ and $\phi_{2}$. These wavelets were adapted to the interval $[0,1]$. However, their construction for the wavelet basis on the interval $[0,1]$ was quite complicated.

In this chapter we take a new approach to the construction of wavelet bases of Hermite cubic splines. In contrast to the semi-orthogonal wavelets of Chui and Wang, the wavelets at different levels are orthogonal with respect to the inner product $\left\langle u^{\prime}, v^{\prime}\right\rangle$, rather than $\langle u, v\rangle$. This requirement of orthogonality is more pertinent to applications of wavelets to numerical solutions of differential equations.

As is well-known, the semi norm is a norm in the underlying Sobolev space for the second order elliptic problems with zero Dirichlet boundary condition. Hence, the wavelets with the above orthogonality form a Riesz basis in Sobolev space and thus stiffness matrices arising from the discretization of the problems by the wavelets have the uniformly bounded condition numbers. this, in turn, ensures the efficiency of iterative methods applied to solving the discretized linear system.

On the other hand, Hermite cubic splines, unlike Daubechie's scaling functions, have explicit expressions with short supports, which are favorite in numerical solutions of partial differential equations. Furthermore, our wavelets
have the same short supports as those of Hermite cubic splines, and this guarantees the efficiency in algorithm. The potential use of such wavelets maybe the numerical solutions of differential equations, and the tensor-product counterparts of our wavelets may serve well for solving partial differential equations in multidimensional spaces. Moreover, changing the orthogonality property with different inner products results in wavelets suitable for numerical solutions of higher order differential equations or integral equations. This is also the motivation of constructing such wavelets.

In Section 1.2 we will give two wavelets $\psi_{1}$ and $\psi_{2}$ as follows:

$$
\begin{aligned}
& \psi_{1}(x)=-2 \phi_{1}(2 x+1)+4 \phi_{1}(2 x)-2 \phi_{1}(2 x-1)-21 \phi_{2}(2 x+1)+21 \phi_{2}(2 x-1), \\
& \psi_{2}(x)=\phi_{1}(2 x+1)-\phi_{1}(2 x-1)+9 \phi_{2}(2 x+1)+12 \phi_{2}(2 x)+9 \phi_{2}(2 x-1)
\end{aligned}
$$

Clearly, $\psi_{1}$ and $\psi_{2}$ are supported on $[-1,1] ; \psi_{1}$ is symmetric and $\psi_{2}$ is antisymmetric. Moreover,

$$
\left\langle\psi_{1}^{\prime}, \phi_{m}^{\prime}(\cdot-j)\right\rangle=\left\langle\psi_{2}^{\prime}, \phi_{m}^{\prime}(\cdot-j)\right\rangle=0, \quad m=1,2, \quad \forall \quad j \in \mathbb{Z}
$$

These wavelets can be easily adapted to the interval $[0,1]$.
By $L_{2}(0,1)$ we denote the space of all square-integrable real-valued functions on $(0,1)$. The inner product in $L_{2}(0,1)$ is defined as

$$
\langle u, v\rangle:=\int_{0}^{1} u(x) v(x) d x, \quad u, v \in L_{2}(0,1) .
$$

Let $H^{1}(0,1)$ be the space of all functions $u$ in $L_{2}(0,1)$ for which (the distributional derivative) $u^{\prime} \in L_{2}(0,1)$. Let $H_{0}^{1}(0,1)$ be the closure of the set

$$
\left\{u \in C[0,1] \cap C^{1}(0,1): u(0)=u(1)=0\right\}
$$

in the space $H^{1}(0,1)$, where $C[0,1]$ denotes the space of all continuous functions on $[0,1]$, and $C^{1}(0,1)$ denotes the space of those continuous functions on $(0,1)$ whose derivatives are also continuous.

For a nonnegative integer $k$, we denote by $\Pi_{k}$ the set of all polynomials of degree at most $k$. For $n \geq 1$, let $V_{n}$ be the space of those cubic splines $v \in C^{1}(0,1) \cap C[0,1]$ for which $v(0)=v(1)=0$ and

$$
\left.\left.v\right|_{\left(j / 2^{n},(j+1) / 2^{n}\right)} \in \Pi_{3}\right|_{\left(j / 2^{n},(j+1) / 2^{n}\right)} \quad \text { for } j=0, \ldots, 2^{n}-1
$$

The dimension of $V_{n}$ is $2^{n+1}$. It is easily seen that the set

$$
\begin{equation*}
\Phi_{n}:=\left\{\phi_{1}\left(2^{n} \cdot-j\right): j=1, \ldots, 2^{n}-1\right\} \cup\left\{\left.\phi_{2}\left(2^{n} \cdot-j\right)\right|_{(0,1)}: j=0, \ldots, 2^{n}\right\} \tag{1.1}
\end{equation*}
$$

is a basis for $V_{n}$. We label the elements in $\Phi_{n}$ as $\left\{v_{1}, v_{2}, \ldots, v_{2^{n+1}}\right\}$.
Let $\Psi_{n}$ be the set of wavelets given by

$$
\begin{equation*}
\Psi_{n}:=\left\{\psi_{1}\left(2^{n} \cdot-j\right): j=1, \ldots, 2^{n}-1\right\} \cup\left\{\left.\psi_{2}\left(2^{n} \cdot-j\right)\right|_{(0,1)}: j=0, \ldots, 2^{n}\right\} \tag{1.2}
\end{equation*}
$$

Let $W_{n}$ be the linear span of $\Psi_{n}$. It is easily seen that $\Psi_{n}$ is a basis for $W_{n}$. Consequently, the dimension of $W_{n}$ is $2^{n+1}$. In Section 1.3 we shall show that

$$
\int_{0}^{1} w^{\prime}(x) v^{\prime}(x) d x=0 \quad \forall w \in \Psi_{n} \text { and } v \in \Phi_{n}
$$

It follows that $V_{n} \cap W_{n}=\{0\}$. Moreover, we have $V_{n+1} \supseteq V_{n}+W_{n}$ and

$$
\operatorname{dim}\left(V_{n+1}\right)=\operatorname{dim}\left(V_{n}\right)+\operatorname{dim}\left(W_{n}\right)
$$

This shows that $V_{n+1}$ is the direct sum of $V_{n}$ and $W_{n}$. Therefore, we have the following decomposition of $H_{0}^{1}(0,1)$ :

$$
H_{0}^{1}(0,1)=V_{1}+W_{1}+W_{2}+\cdots
$$

Recall that $\Phi_{1}=\left\{v_{1}, v_{2}, v_{3}, v_{4}\right\}$. For $n=1,2, \ldots$, we label the elements in $\Psi_{n}$ as follows:

$$
\Psi_{n}=\left\{w_{2^{n+1}+1}, \cdots, w_{2^{n+2}}\right\} .
$$

Let $g_{k}:=v_{k} /\left\|v_{k}^{\prime}\right\|_{2}$ for $k=1,2,3,4$ and $g_{k}:=w_{k} /\left\|w_{k}^{\prime}\right\|_{2}$ for $k>4$. Thus, $\left\|g_{k}^{\prime}\right\|_{2}=1$ for $k=1,2, \ldots$ In Section 1.3 we will show that $\left(g_{k}^{\prime}\right)_{k=1,2 \ldots .}$ is a Riesz sequence in $L_{2}(0,1)$.

In Section 1.4 we shall apply the wavelets constructed in Section 1.3 to numerical solutions of the Sturm-Liouville equation of the form

$$
\begin{equation*}
-\frac{d}{d x}\left(p(x) \frac{d u}{d x}\right)+q(x) u(x)=f(x), \quad x \in(0,1) \tag{1.3}
\end{equation*}
$$

with the Dirichlet boundary condition $u(0)=u(1)=0$. We assume that $p$ and $q$ are continuous functions on $[0,1]$ and

$$
p(x)>0, q(x) \geq 0 \text { for all } x \in[0,1] \text {. Let }
$$

$$
\begin{equation*}
a(u, v):=\int_{0}^{1} p(x) u^{\prime}(x) v^{\prime}(x) d x+\int_{0}^{1} q(x) u(x) v(x) d x, \quad u, v \in H_{0}^{1}(0,1) \tag{1.4}
\end{equation*}
$$

Then the variational form of the above equation with the Dirichlet boundary condition is

$$
a(u, v)=\langle f, v\rangle \quad \forall v \in H_{0}^{1}(0,1)
$$

Wavelets have been used to discretize differential equations. In particular, Xu and Shann [52] successfully applied the wavelet method to numerical solutions of the Sturm-Liouville equation (1.3). The wavelet bases in their paper are anti-derivatives of the Daubechies orthogonal wavelets. Consequently, their basis functions are not locally supported and, in general, the corresponding stiffness matrix is full (not sparse). Furthermore, the condition number of the stiffness matrix is not uniformly bounded.

In application of the wavelet method one often encounters the difficulty that the boundary conditions are hard to impose on wavelets. In our construction, only two wavelets in $\Psi_{n}, \psi_{2}\left(2^{n}\right)$ and $\psi_{2}\left(2^{n} \cdot-2^{n}\right)$, needed to be adapted to the interval $(0,1)$ by means of restriction. This is in sharp contrast to the complexity of the construction of boundary wavelets given in [20].

Recall that $\left\{g_{k}: k=1, \ldots, 2^{n+1}\right\}$ is a wavelet basis for $V_{n}$. Let $A_{n}$ denote the stiffness matrix $\left(a\left(g_{j}, g_{k}\right)\right)_{j, k=1, \ldots, 2^{n+1}}$. In Section 1.4 we will prove that the condition number of $A_{n}$ is uniformly bounded (independent of $n$ ). In particular, for the case $p=1$ and $q=1$, numerical computation suggests that the condition number of $A_{n}$ be less than 3.75 for all $n$. By comparison, the condition number of the stiffness matrix with respect to the wavelet basis constructed in [20] is very large.

At the end of this chapter, we shall provide two numerical examples using the above wavelet basis. The computational results demonstrate the advantage of our wavelet basis.

### 1.2 Spline Wavelets

In this section we construct wavelets on the basis of Hermite cubic splines.
Let $\phi_{1}$ and $\phi_{2}$ be the cubic splines given in Section 1.1. The graphs of $\phi_{1}$ and $\phi_{2}$ are depicted in Figure 1.1. Clearly, both $\phi_{1}$ and $\phi_{2}$ belong to $C^{1}(\mathbb{R})$. Moreover, we have

$$
\phi_{1}(0)=1, \quad \phi_{1}^{\prime}(0)=0, \quad \phi_{2}(0)=0, \quad \phi_{2}^{\prime}(0)=1
$$

Hence, for a function $f \in C^{1}(\mathbb{R})$,

$$
u=\sum_{j \in \mathbb{Z}} f(j) \phi_{1}(\cdot-j)+\sum_{j \in \mathbb{Z}} f^{\prime}(j) \phi_{2}(\cdot-j)
$$

is a Hermite interpolant to $f$ on $\mathbb{Z}$, that is, $u(j)=f(j)$ and $u^{\prime}(j)=f^{\prime}(j)$ for all $j \in \mathbb{Z}$.


Figure 1.1: Hermit cubic splines on $\mathbb{R}$

Let $\Phi:=\left(\phi_{1}, \phi_{2}\right)^{T}$, the transpose of the $1 \times 2$ vector $\left(\phi_{1}, \phi_{2}\right)$. Then $\Phi$ satisfies the following vector refinement equation (see [30]):

$$
\Phi(x)=\sum_{j=-1}^{1} a(j) \Phi(2 x-j), \quad x \in \mathbb{R},
$$

where

$$
a(-1)=\left[\begin{array}{cc}
1 / 2 & 3 / 4 \\
-1 / 8 & -1 / 8
\end{array}\right], a(0)=\left[\begin{array}{cc}
1 & 0 \\
0 & 1 / 2
\end{array}\right], \text { and } a(1)=\left[\begin{array}{cc}
1 / 2 & -3 / 4 \\
1 / 8 & -1 / 8
\end{array}\right] .
$$

Let $S$ be the shift invariant space generated by $\phi_{1}$ and $\phi_{2}$. A function $g$ belongs to $S$ if and only if there are two sequences $b_{1}$ and $b_{2}$ on $\mathbb{Z}$ such that

$$
g=\sum_{j \in \mathbb{Z}}\left[b_{1}(j) \phi_{1}(\cdot-j)+b_{2}(j) \phi_{2}(\cdot-j)\right] .
$$

Let $S_{1}:=\{g(2 \cdot): g \in S\}$. Then $S \subset S_{1}$, since $\Phi$ is refinable. We look for a wavelet space $W$ such that $S_{1}$ is the direct sum of $S$ and $W$. We wish to find two wavelets $\psi_{1}$ and $\psi_{2}$ such that their shifts generate $W$. Moreover, we require

$$
\begin{equation*}
\left\langle\psi_{1}^{\prime}, \phi_{m}^{\prime}(\cdot-j)\right\rangle=\left\langle\psi_{2}^{\prime}, \phi_{m}^{\prime}(\cdot-j)\right\rangle=0, \quad m=1,2, \quad \forall j \in \mathbb{Z} \tag{1.5}
\end{equation*}
$$

For this purpose we need to calculate the inner product of the derivatives of shifts of $\phi_{1}$ and $\phi_{2}$. Note that

$$
\phi_{1}^{\prime}(x):= \begin{cases}-6 x^{2}-6 x & \text { for } x \in[-1,0] \\ 6 x^{2}-6 x & \text { for } x \in[0,1] \\ 0 & \text { otherwise }\end{cases}
$$

and

$$
\begin{gathered}
\phi_{2}^{\prime}(x):= \begin{cases}3 x^{2}+4 x+1 & \text { for } x \in[-1,0], \\
3 x^{2}-4 x+1 & \text { for } x \in[0,1], \\
0 & \text { otherwise }\end{cases} \\
\psi(x)=\sum_{k \in \mathbb{Z}}\left[b_{1}(k) \phi_{1}(2 x-k)+b_{2}(k) \phi_{2}(2 x-k)\right], \quad x \in \mathbb{R} .
\end{gathered}
$$

Then for $j \in \mathbb{Z}$ we have

$$
\begin{aligned}
\left\langle\psi^{\prime}, \phi_{1}^{\prime}(\cdot-j)\right\rangle= & \frac{1}{20}\left[-21 b_{1}(2 j-2)+42 b_{1}(2 j)-21 b_{1}(2 j+2)\right. \\
& \left.-3 b_{2}(2 j-2)+4 b_{2}(2 j-1)-4 b_{2}(2 j+1)+3 b_{2}(2 j+2)\right]
\end{aligned}
$$

and

$$
\begin{aligned}
\left\langle\psi^{\prime}, \phi_{2}^{\prime}(\cdot-j)\right\rangle= & \frac{1}{120}\left[33 b_{1}(2 j-2)-60 b_{1}(2 j-1)+60 b_{1}(2 j+1)\right. \\
& -33 b_{1}(2 j+2)+4 b_{2}(2 j-2)-12 b_{2}(2 j-1) \\
& \left.+28 b_{2}(2 j)-12 b_{2}(2 j+1)+4 b_{2}(2 j+2)\right]
\end{aligned}
$$

For $z \in \mathbb{C} \backslash\{0\}$, let

$$
\begin{array}{ll}
q_{11}(z):=\sum_{j \in \mathbb{Z}} b_{1}(2 j+1) z^{2 j+1}, & q_{12}(z):=\sum_{j \in \mathbb{Z}} b_{1}(2 j) z^{2 j} \\
q_{21}(z):=\sum_{j \in \mathbb{Z}} b_{2}(2 j+1) z^{2 j+1}, & q_{22}(z):=\sum_{j \in \mathbb{Z}} b_{2}(2 j) z^{2 j} .
\end{array}
$$

Then $\left\langle\psi^{\prime}, \phi_{m}^{\prime}(\cdot-j)\right\rangle=0$ for $m=1,2$ and all $j \in \mathbb{Z}$ if and only if

$$
B(z)\left(q_{11}(z), q_{12}(z), q_{21}(z), q_{22}(z)\right)^{T}=0 \quad \forall z \in \mathbb{C} \backslash\{0\}
$$

where

$$
B(z):=\left[\begin{array}{cccc}
0 & -21 z^{2}+42-21 z^{-2} & 4 z-4 z^{-1} & -3 z^{2}+3 z^{-2} \\
-60 z+60 z^{-1} & 33 z^{2}-33 z^{-2} & -12 z-12 z^{-1} & 4 z^{2}+28+4 z^{-2}
\end{array}\right] .
$$

We find two independent solutions as follows:

$$
\left[\begin{array}{c}
q_{11}(z) \\
q_{12}(z) \\
q_{21}(z) \\
q_{22}(z)
\end{array}\right]=\left[\begin{array}{c}
-2 z^{-1}-2 z \\
4 \\
-21 z^{-1}+21 z \\
0
\end{array}\right] \quad \text { and } \quad\left[\begin{array}{c}
q_{11}(z) \\
q_{12}(z) \\
q_{21}(z) \\
q_{22}(z)
\end{array}\right]=\left[\begin{array}{c}
z^{-1}-z \\
0 \\
9 z^{-1}+9 z \\
12
\end{array}\right]
$$

These two solutions induce two wavelets $\psi_{1}$ and $\psi_{2}$ given by

$$
\begin{aligned}
& \psi_{1}(x)=-2 \phi_{1}(2 x+1)+4 \phi_{1}(2 x)-2 \phi_{1}(2 x-1)-21 \phi_{2}(2 x+1)+21 \phi_{2}(2 x-1), \\
& \psi_{2}(x)=\phi_{1}(2 x+1)-\phi_{1}(2 x-1)+9 \phi_{2}(2 x+1)+12 \phi_{2}(2 x)+9 \phi_{2}(2 x-1)
\end{aligned}
$$

By our construction, $\psi_{1}$ and $\psi_{2}$ are supported on $[-1,1]$, they satisfy the conditions in (1.5), and their shifts generate the wavelet space $W$ such that $S_{1}$ is the direct sum of $S$ and $W$. Moreover, $\psi_{1}$ is symmetric and $\psi_{2}$ is antisymmetric (see Figure 1.2).


Figure 1.2: Wavelets $\psi_{1}$ and $\psi_{2}$

Let us take a look at $\psi_{1}^{\prime}$ and $\psi_{2}^{\prime}$. For $0 \leq x \leq 1 / 2$ we have

$$
\begin{aligned}
& \psi_{1}^{\prime}(x)=792 x^{2}-312 x, \quad \psi_{1}^{\prime}(x-1)=-408 x^{2}+120 x \\
& \psi_{2}^{\prime}(x)=552 x^{2}-288 x+24, \quad \psi_{2}^{\prime}(x-1)=168 x^{2}-48 x
\end{aligned}
$$

For $1 / 2 \leq x \leq 1$ we have

$$
\begin{aligned}
& \psi_{1}^{\prime}(x)=408 x^{2}-696 x+288, \quad \psi_{1}^{\prime}(x-1)=-792 x^{2}+1272 x-480 \\
& \psi_{2}^{\prime}(x)=168 x^{2}-288 x+120, \quad \psi_{2}^{\prime}(x-1)=552 x^{2}-816 x+288
\end{aligned}
$$

Hence, the shifts of $\psi_{1}^{\prime}$ and $\psi_{2}^{\prime}$ are linearly independent on the interval $(0,1)$. Because of shift invariance, the shifts of $\psi_{1}^{\prime}$ and $\psi_{2}^{\prime}$ are linear independent on the interval $(k, k+1)$ for every $k \in \mathbb{Z}$. Suppose $b_{1}$ and $b_{2}$ are two square summable sequences on $\mathbb{Z}$. Let

$$
u:=\sum_{j \in \mathbb{Z}}\left[b_{1}(j) \psi_{1}^{\prime}(\cdot-j)+b_{2}(j) \psi_{2}^{\prime}(\cdot-j)\right] .
$$

For $j<k$ or $j>k+1, \psi_{1}^{\prime}(\cdot-j)$ and $\psi_{2}^{\prime}(\cdot-j)$ vanish on $(k, k+1)$. Since the shifts of $\psi_{1}^{\prime}$ and $\psi_{2}^{\prime}$ are linearly independent on $(k, k+1)$, there exist two positive constants $C_{1}$ and $C_{2}$ independent of $k, b_{1}$, and $b_{2}$ such that

$$
C_{1}^{2} \sum_{j=k}^{k+1}\left[\left|b_{1}(j)\right|^{2}+\left|b_{2}(j)\right|^{2}\right] \leq \int_{k}^{k+1}|u(x)|^{2} d x \leq C_{2}^{2} \sum_{j=k}^{k+1}\left[\left|b_{1}(j)\right|^{2}+\left|b_{2}(j)\right|^{2}\right]
$$

It follows that

$$
2 C_{1}^{2} \sum_{j \in \mathbb{Z}}\left[\left|b_{1}(j)\right|^{2}+\left|b_{2}(j)\right|^{2}\right] \leq \int_{\mathbb{R}}|u(x)|^{2} d x \leq 2 C_{2}^{2} \sum_{j \in \mathbb{Z}}\left[\left|b_{1}(j)\right|^{2}+\left|b_{2}(j)\right|^{2}\right]
$$

In other words, the shifts of $\psi_{1}^{\prime}$ and $\psi_{2}^{\prime}$ are stable. See [35] for a study of stability of shifts of several functions.

### 1.3 Wavelets on the Interval

In this section we use the spline wavelets in the previous section to construct a wavelet basis for the space $H_{0}^{1}(0,1)$.

Recall that $V_{n}$ is the linear space of those cubic splines $v \in C^{1}(0,1) \cap C[0,1]$ for which $v(0)=v(1)=0$ and

$$
\left.\left.v\right|_{\left(j / 2^{n},(j+1) / 2^{n}\right)} \in \Pi_{3}\right|_{\left(j / 2^{n},(j+1) / 2^{n}\right)} \quad \text { for } j=0, \ldots, 2^{n}-1
$$

The dimension of $V_{n}$ is $2^{n+1}$. Moreover,
(a) $V_{1} \subset V_{2} \subset \cdots \subset H_{0}^{1}(0,1)$;
(b) $\cup_{n=1}^{\infty} V_{n}$ is dense in $H_{0}^{1}(0,1)$.

Let $\Phi_{n}$ and $\Psi_{n}$ be the sets defined in (1.1) and (1.2), respectively. Then $\Phi_{n}$ is a basis for $V_{n}$. Let $W_{n}$ be the linear span of $\Psi_{n}$. Clearly, $\Psi_{n}$ is a basis for $W_{n}$. Consequently, the dimension of $W_{n}$ is $2^{n+1}$.

We claim that

$$
\begin{equation*}
\int_{0}^{1} w^{\prime}(x) v^{\prime}(x) d x=0 \quad \forall w \in \Psi_{n} \text { and } v \in \Phi_{n} \tag{1.6}
\end{equation*}
$$

Suppose $w=\psi_{r}\left(2^{n} \cdot-j\right)$ for some $r \in\{1,2\}$ and $j \in\left\{1, \ldots, 2^{n}-1\right\}$. Then $\psi_{r}^{\prime}\left(2^{n} \cdot-j\right)$ is supported in the interval $[0,1]$. Hence, for $s=1,2$ and $k \in \mathbb{Z}$, we have

$$
\int_{0}^{1} \psi_{r}^{\prime}\left(2^{n} x-j\right) \phi_{s}^{\prime}\left(2^{n} x-k\right) d x=\int_{\mathbb{R}} \psi_{r}^{\prime}\left(2^{n} x-j\right) \phi_{s}^{\prime}\left(2^{n} x-k\right) d x=0
$$

where (1.5) has been used to derive the second equality. For the same reason, (1.6) is valid if $v=\phi_{s}\left(2^{n} .-k\right)$ for some $s \in\{1,2\}$ and $k \in\left\{1, \ldots, 2^{n}-1\right\}$. Thus, in order to complete the proof of (1.6), it remains to deal with the case
$w=\left.\psi_{2}\left(2^{n} \cdot-j\right)\right|_{(0,1)}$ and $v=\left.\phi_{2}\left(2^{n} \cdot-k\right)\right|_{(0,1)}$ for $j, k \in\left\{0,2^{n}\right\}$. We have $v^{\prime}(x) w^{\prime}(x)=0$ for $x \in(0,1)$ if $j=0$ and $k=2^{n}$, or if $j=2^{n}$ and $k=0$. Hence (1.6) is valid in this case. Suppose $j=k=0$. Since $\psi_{2}$ and $\phi_{2}$ are anti-symmetric, $\psi_{2}^{\prime}$ and $\phi_{2}^{\prime}$ are symmetric. It follows that

$$
\int_{-1}^{0} \psi_{2}^{\prime}(x) \phi_{2}^{\prime}(x) d x=\int_{0}^{1} \psi_{2}^{\prime}(x) \phi_{2}^{\prime}(x) d x
$$

But (1.5) tells us that

$$
\int_{-1}^{1} \psi_{2}^{\prime}(x) \phi_{2}^{\prime}(x) d x=0
$$

Therefore,

$$
\int_{0}^{1} \psi_{2}^{\prime}(x) \phi_{2}^{\prime}(x) d x=0
$$

Consequently,

$$
\int_{0}^{1} \psi_{2}^{\prime}\left(2^{n} x\right) \phi_{2}^{\prime}\left(2^{n} x\right) d x=2^{-n} \int_{0}^{2^{n}} \psi_{2}^{\prime}(x) \phi_{2}^{\prime}(x) d x=0
$$

This verifies (1.6) for $w=\left.\psi_{2}\left(2^{n} \cdot\right)\right|_{(0,1)}$ and $v=\left.\phi_{2}\left(2^{n} \cdot\right)\right|_{(0,1)}$. An analogous argument shows that (1.6) is valid for $w=\left.\psi_{2}\left(2^{n} \cdot-2^{n}\right)\right|_{(0,1)}$ and $v=\phi_{2}\left(2^{n} \cdot-\right.$ $\left.2^{n}\right)\left.\right|_{(0,1)}$. The proof of (1.6) is complete.

It follows from (1.6) that

$$
\int_{0}^{1} w^{\prime}(x) v^{\prime}(x) d x=0 \quad \forall w \in W_{n} \text { and } v \in V_{n}
$$

In particular, $V_{n} \cap W_{n}=\{0\}$. We have $V_{n+1} \supseteq V_{n}+W_{n}$ and

$$
\operatorname{dim}\left(V_{n}+W_{n}\right)=\operatorname{dim}\left(V_{n}\right)+\operatorname{dim}\left(W_{n}\right)=2^{n+1}+2^{n+1}=\operatorname{dim}\left(V_{n+1}\right)
$$

This shows that $V_{n+1}$ is the direct sum of $V_{n}$ and $W_{n}$. Consequently,

$$
V_{n+1}=V_{1}+W_{1}+\cdots+W_{n}
$$

Therefore, we have the following decomposition of $H_{0}^{1}(0,1)$ :

$$
H_{0}^{1}(0,1)=V_{1}+W_{1}+W_{2}+\cdots
$$

Suppose $v \in V_{1}$ and $w_{n} \in W_{n}$ for $n=1,2, \ldots$. The preceding discussion tells us that $\left\langle v^{\prime}, w_{n}^{\prime}\right\rangle=0$ for all $n$ and $\left\langle w_{m}^{\prime}, w_{n}^{\prime}\right\rangle=0$ for $m \neq n$. Hence,

$$
\begin{equation*}
\left\|v^{\prime}+\sum_{n=1}^{\infty} w_{n}^{\prime}\right\|_{L_{2}(0,1)}^{2}=\left\|v^{\prime}\right\|_{L_{2}(0,1)}^{2}+\sum_{n=1}^{\infty}\left\|w_{n}^{\prime}\right\|_{L_{2}(0,1)}^{2} \tag{1.7}
\end{equation*}
$$

For $n=1,2, \ldots$ and $x \in(0,1)$, let

$$
\begin{gathered}
\psi_{n, j}(x):=\left(2^{-n / 2} / \sqrt{729.6}\right) \psi_{1}\left(2^{n} x-j / 2\right) \quad \text { for } j=2,4, \ldots, 2^{n+1}-2 \\
\psi_{n, j}(x):=\left(2^{-n / 2} / \sqrt{153.6}\right) \psi_{2}\left(2^{n} x-(j-1) / 2\right) \quad \text { for } j=3,5, \ldots, 2^{n+1}-1
\end{gathered}
$$

and
$\psi_{n, 1}(x):=\left(2^{-n / 2} / \sqrt{76.8}\right) \psi_{2}\left(2^{n} x\right), \quad \psi_{n, 2^{n+1}}(x):=\left(2^{-n / 2} / \sqrt{76.8}\right) \psi_{2}\left(2^{n} x-2^{n}\right)$.
Note that $\psi_{n, j}$ are so normalized that $\left\|\psi_{n, j}^{\prime}\right\|_{L_{2}(0,1)}=1$ for $j=1, \ldots, 2^{n+1}$.
Theorem 1.1. The sequence $\left(\psi_{n, j}^{\prime}\right)_{n=1,2, \ldots, 1 \leq j \leq 2^{n+1}}$ is a Riesz sequence in $L_{2}(0,1)$. In other words, there exist two positive constants $A$ and $B$ such that

$$
A\left(\sum_{n=1}^{\infty} \sum_{j=1}^{2^{n+1}}\left|b_{n, j}\right|^{2}\right)^{1 / 2} \leq\left\|\sum_{n=1}^{\infty} \sum_{j=1}^{2^{n+1}} b_{n, j} \psi_{n, j}^{\prime}\right\|_{L_{2}(0,1)} \leq B\left(\sum_{n=1}^{\infty} \sum_{j=1}^{2^{n+1}}\left|b_{n, j}\right|^{2}\right)^{1 / 2}
$$

for every sequence $\left(b_{n, j}\right)_{n=1,2, \ldots, 1 \leq j \leq 2^{n+1}}$.
Proof. By (1.7) we have

$$
\left\|\sum_{n=1}^{\infty} \sum_{j=1}^{2^{n+1}} b_{n, j} \psi_{n, j}^{\prime}\right\|_{L_{2}(0,1)}^{2}=\sum_{n=1}^{\infty}\left\|\sum_{j=1}^{2^{n+1}} b_{n, j} \psi_{n, j}^{\prime}\right\|_{L_{2}(0,1)}^{2}
$$

In light of the discussion at the end of Section 1.3, we assert that the shifts of $\psi_{1}^{\prime}$ and $\psi_{2}^{\prime}$ are linearly independent on $(k, k+1)$ for every $k \in \mathbb{Z}$. Hence, there exist two positive constants $A$ and $B$ (independent of $n$ ) such that

$$
A^{2} \sum_{j=1}^{2^{n+1}}\left|b_{n, j}\right|^{2} \leq\left\|\sum_{j=1}^{2^{n+1}} b_{n, j} \psi_{n, j}^{\prime}\right\|_{L_{2}(0,1)}^{2} \leq B^{2} \sum_{j=1}^{2^{n+1}}\left|b_{n, j}\right|^{2}
$$

This completes the proof of the theorem.

For $x \in(0,1)$, let

$$
\begin{aligned}
\phi_{1,1}(x) & :=\sqrt{5 / 24} \phi_{1}(2 x-1), \\
\phi_{1,2}(x) & :=\sqrt{15 / 4} \phi_{2}(2 x), \\
\phi_{1,3}(x) & :=\sqrt{15 / 8} \phi_{2}(2 x-1), \\
\phi_{1,4}(x) & :=\sqrt{15 / 4} \phi_{2}(2 x-2) .
\end{aligned}
$$

Note that each $\phi_{1, j}$ is so normalized that $\left\|\phi_{1, j}^{\prime}\right\|_{L_{2}(0,1)}=1, j=1, \ldots, 4$. Clearly, $V_{1}$ is spanned by $\phi_{1, j}, j=1, \ldots, 4$. Consequently, $H_{0}^{1}(0,1)$ is spanned by $\phi_{1, j}$, $j=1, \ldots, 4$, together with $\psi_{n, j}, n=1,2, \ldots, j=1, \ldots, 2^{n+1}$. We relabel these functions as follows. Let $g_{j}:=\phi_{1, j}$ for $j=1, \ldots, 4$, and let $g_{2^{n+1}+j}:=\psi_{n, j}$ for $n=1,2, \ldots$ and $j=1, \ldots, 2^{n+1}$. With the same reasoning as in the proof of Theorem 1.1, we see that the sequence $\left(g_{k}^{\prime}\right)_{k=1,2, \ldots}$ is a a Riesz sequence in $L_{2}(0,1)$. In other words, there exist two positive constants $A$ and $B$ such that

$$
\begin{equation*}
A\left(\sum_{k=1}^{\infty}\left|b_{k}\right|^{2}\right)^{1 / 2} \leq\left\|\sum_{k=1}^{\infty} b_{k} g_{k}^{\prime}\right\|_{L_{2}(0,1)} \leq B\left(\sum_{k=1}^{\infty}\left|b_{k}\right|^{2}\right)^{1 / 2} \tag{1.8}
\end{equation*}
$$

for every square summable sequence $\left(b_{k}\right)_{k=1,2, \ldots}$.

### 1.4 Applications

In this section the wavelets constructed in the previous section are used to solve differential equations. We shall confine ourselves to the Sturm-Liouville equation of the form (1.3) with the Dirichlet boundary condition $u(0)=u(1)=$ 0 . We assume that $p$ and $q$ are continuous functions on $[0,1]$ and $p(x)>0$, $q(x) \geq 0$ for all $x \in[0,1]$.

For $u, v \in H_{0}^{1}(0,1)$, let $a(u, v)$ be the bilinear form given in (1.4). Then the variational form of equation (1.3) with the Dirichlet boundary condition is

$$
\begin{equation*}
a(u, v)=\langle f, v\rangle \quad \forall v \in H_{0}^{1}(0,1) \tag{1.9}
\end{equation*}
$$

The corresponding Galerkin approximation problem is the following: find $u_{n} \in$ $V_{n}$ such that

$$
\begin{equation*}
a\left(u_{n}, v\right)=\langle f, v\rangle \quad \forall v \in V_{n} \tag{1.10}
\end{equation*}
$$

By the Lax-Milgram lemma (see, e.g., [7, p. 60]), the approximation problem (1.10) has a unique solution.

We propose to use the wavelet set $G_{n}:=\left\{g_{1}, \ldots, g_{2^{n+1}}\right\}$ as a basis for $V_{n}$. Recall that $g_{j}:=\phi_{1, j}$ for $j=1, \ldots, 4$, and $g_{2^{n+1}+j}:=\psi_{n, j}$ for $n=1,2, \ldots$ and $j=1, \ldots, 2^{n+1}$, where $\phi_{1, j}(j=1, \ldots, 4)$ and $\psi_{n, j}\left(j=1, \ldots, 2^{n+1}\right)$ are the functions constructed in the previous section. With this basis for $V_{n}$, the Galerkin approximation problem (1.10) can be discretized as follows:

$$
\sum_{k=1}^{2^{n+1}} a\left(g_{j}, g_{k}\right) \eta_{k}=\left\langle g_{j}, f\right\rangle, \quad j=1, \ldots, 2^{n+1}
$$

The stiffness matrix

$$
\left(a\left(g_{j}, g_{k}\right)\right)_{1 \leq j, k \leq 2^{n+1}}
$$

is denoted by $A_{n}$. We will prove that the condition number of $A_{n}$ is uniformly bounded (independent of $n$ ). Therefore, the wavelet basis $G_{n}$ is a good tool for preconditioning.

Let us recall that the condition number of an invertible square matrix $A$ is defined by

$$
\operatorname{cond}(A):=\|A\|\left\|A^{-1}\right\|
$$

where $\|\cdot\|$ is a matrix norm. The spectral condition number of $A$ is defined as

$$
\frac{\max _{i}\left|\lambda_{i}(A)\right|}{\min _{i}\left|\lambda_{i}(A)\right|}
$$

where the numbers $\lambda_{i}(A)$ are eigenvalues of $A$. If $A$ is a (real) symmetric matrix, then its condition number with respect to the 2 -norm is equal to its spectral condition number (see [10, p. 51]).

Theorem 1.2. The condition number of the stiffness matrix $A_{n}$ is uniformly bounded (independent of $n$ ).

Proof. It suffices to show that there exist two positive constants $B$ and $C$ independent of $n$ such that

$$
\begin{equation*}
B\left(\sum_{j=1}^{4}\left|c_{j}\right|^{2}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}}\left|b_{m, j}\right|^{2}\right) \leq a(u, u) \leq C\left(\sum_{j=1}^{4}\left|c_{j}\right|^{2}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}}\left|b_{m, j}\right|^{2}\right) \tag{1.11}
\end{equation*}
$$

for any

$$
u=\sum_{j=1}^{4} c_{j} \phi_{1, j}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}} b_{m, j} \psi_{m, j}
$$

By (1.8) there exists a positive constant $C_{1}$ independent of $n$ such that

$$
\begin{aligned}
\left\|u^{\prime}\right\|_{L_{2}(0,1)} & =\left\|\sum_{j=1}^{4} c_{j} \phi_{1, j}^{\prime}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}} b_{m, j} \psi_{m, j}^{\prime}\right\|_{L_{2}(0,1)} \\
& \geq C_{1}\left(\sum_{j=1}^{4}\left|c_{j}\right|^{2}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}}\left|b_{m, j}\right|^{2}\right)^{1 / 2}
\end{aligned}
$$

But

$$
a(u, u) \geq\left\langle p u^{\prime}, u^{\prime}\right\rangle \geq \mu\left\langle u^{\prime}, u^{\prime}\right\rangle=\mu\left\|u^{\prime}\right\|_{L_{2}(0,1)}^{2}
$$

where $\mu:=\min _{x \in[0,1]} p(x)>0$. This establishes the first inequality in (1.11). Furthermore, we observe that

$$
a(u, u) \leq \nu\left(\|u\|_{L_{2}(0,1)}^{2}+\left\|u^{\prime}\right\|_{L_{2}(0,1)}^{2}\right)
$$

where $\nu:=\max _{0 \leq x \leq 1}\{p(x), q(x)\}<\infty$. By (1.8) there exists a positive constant $C_{2}$ independent of $n$ such that

$$
\left\|u^{\prime}\right\|_{L_{2}(0,1)} \leq C_{2}\left(\sum_{j=1}^{4}\left|c_{j}\right|^{2}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}}\left|b_{m, j}\right|^{2}\right)^{\frac{1}{2}}
$$

Moreover,

$$
\|u\|_{L_{2}(0,1)} \leq\left\|\sum_{j=1}^{4} c_{j} \phi_{1, j}\right\|_{L_{2}(0,1)}+\sum_{m=1}^{n-1}\left\|\sum_{j=1}^{2^{m+1}} b_{m, j} \psi_{m, j}\right\|_{L_{2}(0,1)}
$$

Note that $\left\|\psi_{m, j}\right\|_{L_{2}(0,1)}=O\left(2^{-m}\right)$ as $m \rightarrow \infty$. Hence, there exists a positive constant $C_{3}$ independent of $n$ such that

$$
\|u\|_{L_{2}(0,1)} \leq C_{3}\left[\left(\sum_{j=1}^{4}\left|c_{j}\right|^{2}\right)^{1 / 2}+\sum_{m=1}^{n-1} 2^{-m}\left(\sum_{j=1}^{2^{m+1}}\left|b_{m, j}\right|^{2}\right)^{1 / 2}\right]
$$

With the help of the Schwarz inequality we see that there exists a positive constant $C_{4}$ independent of $n$ such that

$$
\|u\|_{L_{2}(0,1)}^{2} \leq C_{4}\left(\sum_{j=1}^{4}\left|c_{j}\right|^{2}+\sum_{m=1}^{n-1} \sum_{j=1}^{2^{m+1}}\left|b_{m, j}\right|^{2}\right) .
$$

The second inequality in (1.11) follows. The proof of the theorem is complete.

In what follows we apply the wavelet basis $G_{n}$ to two numerical examples. Example 1. Consider the Dirichlet problem:

$$
\left\{\begin{array}{l}
-u^{\prime \prime}=f \quad \text { on } \quad(0,1) \\
u(0)=u(1)=0
\end{array}\right.
$$

where $f$ is given by

$$
f(x)=(53.7 \pi)^{2} \sin (53.7 \pi x)+(2.3 \pi)^{2} \sin (2.3 \pi x), \quad x \in(0,1)
$$

The exact solution of the problem is

$$
\begin{equation*}
u(x)=\sin (53.7 \pi x)+\sin (2.3 \pi x), \quad x \in(0,1) \tag{1.12}
\end{equation*}
$$

which could be regarded as the sum of a high-frequency component and a low-frequency component. Let us use the wavelet basis $G_{n}:=\left\{g_{1}, \ldots, g_{2^{n+1}}\right\}$ to solve the Dirichlet problem. With $u_{n}=\sum_{k=1}^{2^{n+1}} \eta_{k} g_{k}$, the Galerkin approximation problem (1.10) is discretized as

$$
\begin{equation*}
\sum_{k=1}^{2^{n+1}}\left\langle g_{j}^{\prime}, g_{k}^{\prime}\right\rangle \eta_{k}=\left\langle g_{j}, f\right\rangle, \quad j=1, \ldots, 2^{n+1} \tag{1.13}
\end{equation*}
$$

The stiffness matrix $A_{n}:=\left(\left\langle g_{j}^{\prime}, g_{k}^{\prime}\right\rangle\right)_{1 \leq j, k \leq 2^{n+1}}$ is block diagonal. Moreover, each block is a banded matrix. By Theorem 1.2, the condition number of the matrix $A_{n}$ is uniformly bounded (independent of $n$ ). This assertion is confirmed by numerical computation of the maximal eigenvalue $\lambda_{\max }$, the minimal eigenvalue $\lambda_{\text {min }}$, and the condition number $\kappa=\lambda_{\max } / \lambda_{\text {min }}$ of the matrix $A_{n}$ for $n=6, \ldots, 12$ (see Table 1.1).

We use the CG (conjugate gradient) method to solve the above system (1.13) of linear equations. Since the stiffness matrix $A_{n}$ is well conditioned, the CG method converges very fast. Up to $n=12$, only 21 iterations are needed for convergence to the solution of the system of linear equations. Here and in what follows, we take $10^{-10}$ as the threshold to stop the iterations. For

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max }$ | 1.5780 | 1.5787 | 1.5789 | 1.5789 | 1.5789 | 1.5789 | 1.5789 |
| $\lambda_{\min }$ | 0.4220 | 0.4213 | 0.4211 | 0.4211 | 0.4211 | 0.4211 | 0.4211 |
| $\kappa$ | 3.7397 | 3.7474 | 3.7494 | 3.7498 | 3.7498 | 3.7498 | 3.7498 |

Table 1.1: Condition number of the matrix $A_{n}$
$n=1,2, \ldots$, let $e_{n}:=\left\|u_{n}-u\right\|_{L_{2}(0,1)}$, where $u$ is the exact solution given in (1.12). For $n=6, \ldots, 12$, the following table lists the error $e_{n}$ and the rate of convergence $\log _{2} e_{n-1} / e_{n}$.

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $e_{n}$ | $1.21-2$ | $1.33-3$ | $1.08-4$ | $7.36-6$ | $4.71-7$ | $2.96-8$ | $1.85-9$ |
| $\log _{2}\left(\frac{e_{n-1}}{e_{n}}\right)$ | 4.10 | 3.19 | 3.62 | 3.88 | 3.97 | 3.99 | 4.00 |

Table 1.2: Error $e_{n}$ and its convergence rate

It is well known from approximation theory that the Hermite cubic splines provide approximation of order 4. The preceding computation confirms this assertion.

If we use the finite elements in $\Phi_{n}$ given in (1.1) to discretize the equation (1.10), then the resulting stiffness matrix is ill conditioned. For $n=12$, the system of linear equations has 8192 unknowns. Without preconditioning, it takes more than 2000 iterations for the $C G$ method to converge. The following graph depicts the error against the number of iterations.

In [6], Bramble, Pasciak, and Xu proposed the so-called BPX method for preconditioning. This method was developed on the nodal basis (piecewise linear functions). The corresponding spectral condition number (not necessarily the condition number) was shown to be uniformly bounded. For $n=6, \ldots, 12$, the following table gives the maximal eigenvalue $\lambda_{\max }$, the minimal eigenvalue $\lambda_{\min }$, and the spectral condition number of the corresponding matrix after preconditioning:


Figure 1.3: The error against the number of iterations without preconditioning

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max }$ | 4.390 | 4.725 | 5.004 | 5.238 | 5.437 | 5.607 | 5.753 |
| $\lambda_{\min }$ | 0.9311 | 0.9297 | 0.9291 | 0.9323 | 0.9316 | 0.9311 | 0.9308 |
| $\kappa$ | 4.715 | 5.082 | 5.385 | 5.619 | 5.836 | 6.021 | 6.180 |

Table 1.3: BPX preconditioning for nodal basis

We observe that piecewise linear functions only provide approximation of order 2. In order to achieve convergence of order 4, one may extend the BPX method (or additive Schwarz method) to Hermite cubic splines. We will prove that BPX method is still valid for Hermite cubic splines in Appendix A [37]. For $n=6, \ldots, 12$, the following table gives the maximal eigenvalue $\lambda_{\text {max }}$, the minimal eigenvalue $\lambda_{\text {min }}$, and the spectral condition number of the corresponding matrix after preconditioning:

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max }$ | 3.562 | 3.632 | 3.682 | 3.718 | 3.743 | 3.763 | 3.777 |
| $\lambda_{\min }$ | 0.7693 | 0.7696 | 0.7696 | 0.7696 | 0.7696 | 0.7696 | 0.7696 |
| $\kappa$ | 4.630 | 4.719 | 4.784 | 4.831 | 4.864 | 4.890 | 4.907 |

Table 1.4: BPX preconditioning for Hermite cubic splines

We see that the condition number induced by our wavelet basis is smaller than that given by the BPX method. For $n=12$, after preconditioning by the BPX method, it takes 26 iterations for the PCG (preconditioned conjugate gradient) method to converge. Hence, the preconditioning method induced by our wavelet basis is competitive.

Example 2. Consider the Dirichlet problem

$$
\left\{\begin{array}{l}
-u^{\prime \prime}+u=f \quad \text { on } \quad(0,1) \\
u(0)=u(1)=0
\end{array}\right.
$$

where

$$
f(x)=\left[(53.7 \pi)^{2}+1\right] \sin (53.7 \pi x)+\left[(2.3 \pi)^{2}+1\right] \sin (2.3 \pi x), \quad x \in(0,1)
$$

The function $u$ given in (1.12) is the exact solution of the problem.
In this case, the bilinear form $a(u, v)$ is given by

$$
a(u, v)=\left\langle u^{\prime}, v^{\prime}\right\rangle+\langle u, v\rangle, \quad u, v \in H_{0}^{1}(0,1) .
$$

With the wavelet basis $G_{n}$ the Galerkin approximation problem (1.10) is discretized as

$$
\begin{equation*}
\sum_{k=1}^{2^{n+1}}\left(\left\langle g_{j}^{\prime}, g_{k}^{\prime}\right\rangle+\left\langle g_{j}, g_{k}\right\rangle\right) \eta_{k}=\left\langle g_{j}, f\right\rangle, \quad j=1, \ldots, 2^{n+1} \tag{1.14}
\end{equation*}
$$

The stiffness matrix

$$
A_{n}:=\left(\left\langle g_{j}^{\prime}, g_{k}^{\prime}\right\rangle+\left\langle g_{j}, g_{k}\right\rangle\right)_{1 \leq j, k \leq 2^{n+1}}
$$

is still a sparse matrix. By Theorem 1.2, the condition number of the matrix $A_{n}$ is uniformly bounded (independent of $n$ ). This assertion is confirmed by

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max }$ | 1.5780 | 1.5787 | 1.5789 | 1.5789 | 1.5789 | 1.5789 | 1.5789 |
| $\lambda_{\min }$ | 0.4220 | 0.4213 | 0.4211 | 0.4211 | 0.4211 | 0.4211 | 0.4211 |
| $\kappa$ | 3.7396 | 3.7474 | 3.7494 | 3.7498 | 3.7498 | 3.7498 | 3.7498 |

Table 1.5: Condition number of the matrix $A_{n}$
numerical computation of the maximal eigenvalue $\lambda_{\max }$, the minimal eigenvalue $\lambda_{\text {min }}$, and the condition number $\kappa$ of $A_{n}$ for $n=6, \ldots, 12$ (see Table 1.5).

We use the CG method to solve the above system (1.14) of linear equations. The computational results are similar to those in Example 1. Up to $n=12$, only 19 iterations are needed for convergence to the solution of the system of linear equations. For $n=6, \ldots, 12$, the following table lists the error $e_{n}$ and the rate of convergence $\log _{2} e_{n-1} / e_{n}$.

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $e_{n}$ | $1.21-2$ | $1.33-3$ | $1.08-4$ | $7.36-6$ | $4.71-7$ | $2.97-8$ | $1.92-9$ |
| $\log _{2}\left(\frac{e_{n-1}}{e_{n}}\right)$ | 4.10 | 3.19 | 3.62 | 3.88 | 3.97 | 3.99 | 3.95 |

Table 1.6: Error $e_{n}$ and its convergence rate

BPX method is applied for the comparison. Up to $n=12,21$ iterations are needed for convergence to the solution of the system of linear equations. The maximal and minimal eigenvalues of the preconditioned system, as well as spectral condition numbers, are listed in Table 1.7.

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max }$ | 3.562 | 3.632 | 3.682 | 3.718 | 3.743 | 3.763 | 3.777 |
| $\lambda_{\min }$ | 0.7696 | 0.7696 | 0.7696 | 0.7696 | 0.7696 | 0.7696 | 0.7696 |
| $\kappa$ | 4.628 | 4.719 | 4.784 | 4.831 | 4.864 | 4.890 | 4.907 |

Table 1.7: BPX preconditioning for Hermite cubic splines

Finally, we remark that our wavelet basis can also be used to solve inte-
gral equations numerically. A discrete wavelet Petrov-Galerkin method was developed by Chen, Micchelli, and Xu [12] for numerical solutions of integral equations of the second kind with weakly singular kernels. Recently, Shen and Lin [45] used the wavelet basis $G_{n}$ constructed in this chapter to find numerical solutions of integral equations on the upper half plane.

## Chapter 2

## Modified Hierarchy Basis For Solving Singular Boundary <br> Value Problems

### 2.1 Background

Our investigations in this chapter is concerned with the preconditioning method on the basis of the modified hierarchy basis for the numerical solution of the singular boundary value problem arising from the radically symmetric elliptic partial differential equations, a problem with numerous applications (see, e.g., [44]).

When the Dirichlet problem

$$
\begin{aligned}
-\Delta u(\mathbf{x})+q(\mathbf{x}) u(\mathbf{x}) & =f(\mathbf{x}), \quad \text { in } \quad \mathbf{B} \\
u & =0, \quad \text { on } \quad \partial \mathbf{B}
\end{aligned}
$$

is defined on a unit ball $\mathbf{B}:=B_{1}(0)$ in $\mathbb{R}^{d}$ and the data depend only on the radical coordinate, then after a change of variables, the problem will reduce to a one dimensional singular boundary value problem,

$$
\begin{aligned}
& -u^{\prime \prime}(x)-\frac{d-1}{x} u^{\prime}(x)+q(x) u(x)=f(x), \quad x \in(0,1) \\
& u^{\prime}(0)=u(1)=0
\end{aligned}
$$

where $q(x) \geq 0$ and $q(x) \in L_{\infty}(0,1)$.
For the smooth data, it has been proven (see, e.g., [24, 32, 43, 44]) that the (smooth) solution can be approximated with high order accuracy by the Galerkin method with a piecewise polynomial subspace. Therefore, no special functions are required in the basis.

Convergence results of the Galerkin method for the singular boundary values problems have been studied for the case $q(x)=0$ in detail in [32]. In [24], Eriksson and Thomee established the general optimal order error estimates and even generalized their results to the corresponding time dependent problems. It shows that the Galerkin method would give the same convergence results for the singular problems as for the nonsingular problems.

For the solution with certain smoothness (such as in $H^{2}$ ), the simple piecewise linear nodal basis shall satisfy the approximation requirement. By the error estimates provided in Section 2.3, we show that a slightly modified piecewise linear nodal basis provides the suitable approximation order.

However, it is still a challenging problem to efficiently solve the large system of linear equations arising from the Galerkin method for the singular boundary value problems. Like its counterpart for the regular elliptic problems, the linear system arising from the Galerkin method for the singular boundary value problems is also ill conditioned. For the regular elliptic boundary value problem, multigrid methods (see, e.g., $[5,4]$ ), and numerous other preconditioning methods (see, e.g., [6]), were successfully developed. Nevertheless, to our best knowledge, presently there are few references about preconditioning methods of the Galerkin method for the singular problems. To design an easily implemented preconditioning method through the construction of the modified hierarchy basis shall be the principle goal of this chapter.

The hierarchy basis has been discussed extensively in [55,56], and has been proven to be an efficient preconditioning method for low dimensional regular elliptic problems. In this chapter, we construct a modified hierarchy basis based on the concept of "stability" (see, e.g., [33, 35, 41]), and the "norm equivalence" for the Sobolev space (see, e.g., [2, 34, 41, 29]). Such basis is
then adapted to the nodal basis introduced in section 2.2 for the singular boundary value problem, and thus the preconditioning can be achieved. It will be shown later that after applying the preconditioning method based on the modified hierarchy basis, the condition number of the stiffness matrix arising from the Galerkin method will be uniformly bounded. In particular, the condition number is nicely bounded by 2 for the case $q(x)=0$.

This chapter is divided into three parts. In section 2.2, we propose the preconditioning method on the basis of the modified hierarchy basis for the singular boundary value problem, and show the connection between the concept of norm equivalence and stability of the modified hierarchy basis. The condition number of the preconditioned stiffness matrix is proven to be uniformly bounded. In section 2.3, we provide basic error estimates for the Galerkin approximation from the piecewise linear nodal basis subspace $V_{h}$ with its element $v$ satisfying the boundary conditions $v^{\prime}(0)=v(1)=0$. We will show such subspace provides the same approximation order as the linear nodal basis subspace without the condition $v^{\prime}(0)=0$. Numerical examples are computed to confirm our results in section 2.4.

### 2.2 The Galerkin method and the modified hierarchy basis

We consider the boundary value problem of the form

$$
\begin{align*}
& -\left(x^{\alpha} u^{\prime}(x)\right)^{\prime}+x^{\alpha} q(x) u(x)=x^{\alpha} f(x), \quad x \in(0,1)  \tag{2.1}\\
& u^{\prime}(0)=u(1)=0 \tag{2.2}
\end{align*}
$$

where $\alpha=d-1$.
Let $v$ be a real-valued Lebesgue measurable function on $\mathbb{R}$. We define the $L_{2}(0,1)$ inner product by

$$
\langle u, v\rangle:=\int_{0}^{1} u(x) v(x) d x
$$

and $L_{2}(0,1)$ space by

$$
L_{2}(0,1):=\left\{v: \quad\|v\|_{L_{2}(0,1)}<\infty\right\} .
$$

The weighted $L_{2}$ space $\dot{L}_{2}(0,1)$ is defined by

$$
\dot{L}_{2}(0,1):=\left\{v: \int_{0}^{1}\left|x^{\frac{\alpha}{2}} v(x)\right|^{2} d x<\infty\right\}
$$

The weighted Sobolev space $\dot{H}_{0}^{1}(0,1)$ is the closure of the set $\{v: v \in$ $\left.C([0,1]) \cap C^{1}(0,1), v(1)=0\right\}$ in the sense of the weighted Sobolev norm

$$
\|v\|_{\dot{H}^{1}}:=\left(\int_{0}^{1} x^{\alpha}\left(|v(x)|^{2}+\left|v^{\prime}(x)\right|^{2}\right) d x\right)^{1 / 2} .
$$

Define the symmetric bilinear form $a(\cdot, \cdot)$ as follows: for $u, v \in \dot{H}_{0}^{1}(0,1)$,

$$
\begin{equation*}
a(u, v):=\int_{0}^{1} x^{\alpha} u^{\prime}(x) v^{\prime}(x) d x+\int_{0}^{1} q(x) x^{\alpha} u(x) v(x) d x \tag{2.3}
\end{equation*}
$$

Then the solution $u$ of the singular boundary value problem also solves the variational problem

$$
\begin{equation*}
a(u, v)=\left\langle x^{\alpha} f(x), v(x)\right\rangle, \quad \forall v \in \dot{H}_{0}^{1}(0,1) . \tag{2.4}
\end{equation*}
$$

Here, with some ambiguity, we also use $x^{\alpha}$ to denote function $x \mapsto x^{\alpha}, x \in$ $(0,1)$, and we assume that $f \in \dot{L}_{2}(0,1)\left(x^{\frac{\alpha}{2}} f(x) \in L_{2}(0,1)\right)$.

We have the following Poincare-type inequality ([32]).

## Lemma 2.1.

$$
\left\|x^{\alpha / 2} v\right\|_{L_{2}} \leq \frac{2}{\alpha+1}\left\|x^{\alpha / 2} v^{\prime}\right\|_{L_{2}}, \quad v \in \dot{H}_{0}^{1}
$$

Proof. We have

$$
\begin{aligned}
\int_{0}^{1} x^{\alpha} v^{2}(x) d x & =\int_{0}^{1}\left(\frac{x^{\alpha+1}}{\alpha+1}\right)^{\prime} v^{2}(x) d x \\
& =-\int_{0}^{1}\left(\frac{x^{\alpha+1}}{\alpha+1}\right) 2 v(x) v^{\prime}(x) d x+\left.\left(\frac{x^{\alpha+1}}{\alpha+1}\right) v^{2}(x)\right|_{0} ^{1} \\
& \leq \frac{2}{\alpha+1}\left\|x^{\frac{\alpha}{2}} v\right\|_{L_{2}(0,1)}\left\|x^{\frac{\alpha+2}{2}} v^{\prime}\right\|_{L_{2}(0,1)} \\
& \leq \frac{2}{\alpha+1}\left\|x^{\frac{\alpha}{2}} v\right\|_{L_{2}(0,1)}\left\|x^{\frac{\alpha}{2}} v^{\prime}\right\|_{L_{2}(0,1)}\|x\|_{L_{\infty}(0,1)} .
\end{aligned}
$$

Since $\|x\|_{L_{\infty}(0,1)} \leq 1$, this completes the lemma.

Now we define another inner product for $\dot{H}_{0}^{1}(0,1)$ by

$$
\begin{equation*}
\langle u, v\rangle_{E}:=\int_{0}^{1} x^{\alpha} u^{\prime}(x) v^{\prime}(x) d x, \quad u, v \in \dot{H}_{0}^{1} \tag{2.5}
\end{equation*}
$$

By Lemma 2.1, we have the following inequalities:

$$
\begin{equation*}
\langle v, v\rangle_{E} \leq a(v, v) \leq\left(1+\left(\frac{2}{\alpha+1}\right)^{2}\|q\|_{L_{\infty}(0,1)}\right)\langle v, v\rangle_{E}, \quad v \in \dot{H}_{0}^{1} \tag{2.6}
\end{equation*}
$$

Hereafter, we fix $\alpha=1$ for simplicity. The case $\alpha>1$ can be handled in the same way without any extra difficulty.

For the uniform partition of $[0,1], 0=x_{0}<x_{1}<\ldots<x_{2^{n}}=1, x_{j}=2^{-n} j$, $j=0, \ldots, 2^{n}$, let $\phi$ be the hat function $\phi(x):=\max \{0,1-|x|\}$, and

$$
\begin{align*}
\phi_{n, 1} & :=\left(\phi\left(2^{n} \cdot\right)+\phi\left(2^{n} \cdot-1\right)\right) \chi_{[0,1]}  \tag{2.7}\\
\phi_{n, j} & :=\phi\left(2^{n} \cdot-j\right), j=2, \ldots, 2^{n}-1, \tag{2.8}
\end{align*}
$$

where $\chi_{[a, b]}, a<b$, is the characteristic function on the interval $[a, b]$. Let

$$
V_{n}:=\operatorname{span}\left\{\phi_{n, j}: j=1, \ldots, 2^{n}-1\right\} .
$$

It is easily seen that $V_{n} \subset V_{n+1}$ for $n=1,2, \ldots$.
The Galerkin method is defined as seeking the element $u_{n} \in V_{n}$ such that

$$
\begin{equation*}
a\left(u_{n}, v\right)=\langle x f, v\rangle, \quad v \in V_{n} . \tag{2.9}
\end{equation*}
$$

Lemma 2.1 shows that $a(\cdot, \cdot)$ is elliptic, and by the Lax-Milgram theorem, existence and uniqueness of the solution are guaranteed for both (2.4) and (2.9).

Taking

$$
u_{n}=\sum_{j=1}^{2^{n}-1} c_{n, j} \phi_{n, j}
$$

we can rewrite (2.9) as

$$
\begin{equation*}
\sum_{j=1}^{2^{n}-1} a\left(\phi_{n, j}, \phi_{n, l}\right) c_{n, j}=\left\langle x f, \phi_{n, l}\right\rangle, \quad l=1, \ldots, 2^{n}-1 \tag{2.10}
\end{equation*}
$$

or more briefly,

$$
\begin{equation*}
A_{n} C_{n}=F_{n}, \tag{2.11}
\end{equation*}
$$

where $(j, l)$ entry of the $2^{n}-1$ by $2^{n}-1$ stiffness matrix $A_{n}$ is $a\left(\phi_{n, j}, \phi_{n, l}\right)$, $C_{n}:=\left(c_{n, 1}, \ldots, c_{n, 2^{n}-1}\right)^{T}$, and $F_{n}:=\left(\left\langle x f, \phi_{n, 1}\right\rangle, \ldots,\left\langle f, \phi_{n, 2^{n}-1}\right\rangle\right)^{T}$. Here, the superscript $T$ denotes the transpose of a vector or a matrix.

The condition number of a nonsingular $M$ by $M$ matrix $A$ is defined by

$$
\kappa(A):=\|A\|\left\|A^{-1}\right\|
$$

where $\|A\|:=\sup _{\mathbf{x} \in \mathbb{R}^{M}} \frac{\|A \mathbf{x}\|}{\|\mathbf{x}\|}, \mathbf{x}:=\left(x_{1}, \ldots, x_{M}\right)^{T}$, and $\|\mathbf{x}\|:=\left(\sum_{i=1}^{M} x_{i}^{2}\right)^{\frac{1}{2}}$.
When $A$ is positive definite and symmetric, we have

$$
\kappa(A)=\frac{\lambda_{\max , A}}{\lambda_{\min , A}}
$$

where $\lambda_{\max , A}, \lambda_{\min , A}$ are the maximum and the minimum eigenvalues of the matrix $A$, respectively.

The following error estimate will be established in the next section:

$$
\left\|x^{1 / 2}\left(u-u_{n}\right)\right\|_{L_{2}} \leq C\left(2^{-n}\right)^{2}\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}}
$$

Consequently, the subspace $V_{n}$ has to be large enough to guarantee that the error $u-u_{n}$ is sufficiently small. However, increasing the number $n$ will dramatically increase the condition number of the associated stiffness matrix $A_{n}$ (see, e.g., [5]), which makes solving $u_{n}$ numerically difficult. It is well-known in the literature that for an ill-conditioned large linear system, without any preconditioning, it's impossible to find an efficient solver. Therefore, seeking a suitable preconditioning method will be important for solving the discretized system numerically. There is an abundance of literature contributed to this topic for the regular elliptic boundary problems, such as $[6,29,55,56]$. Recently, wavelet methods have been introduced to serve as new preconditioning methods (see, e.g., [18, 28, 41, 51, 52]). Stability plays the key role in the
wavelet preconditioning method. In other words, if one is able to find a basis which is stable in the corresponding Sobolev space, then the condition number of the associated stiffness matrix is uniformly bounded. A basis, say $\left\{\psi_{i}\right\}_{i=1}^{\infty}$, is stable if it satisfies,

$$
C_{0}\left(\sum_{i=1}^{\infty} c_{i}^{2}\right) \leq\left\|\sum_{i=1}^{\infty} c_{i} \psi_{i}\right\|^{2} \leq C_{1}\left(\sum_{i=1}^{\infty} c_{i}^{2}\right)
$$

where $C_{0}, C_{1}$ are two positive constants independent of $\left\{c_{i}\right\}_{i=1}^{\infty}$, and $\|\cdot\|$ refers to the norm for the space in which we are interested. Stability of the shift invariant space has been studied extensively in [33, 35].

To find a proper preconditioning matrix for $A_{n}$ in (2.11), we introduce the following lemma.

Lemma 2.2. If two positive definite symmetric $M \times M$ matrices $A, B$ satisfy the following condition

$$
C_{0} \mathbf{x}^{T} B \mathbf{x} \leq \mathbf{x}^{T} A \mathbf{x} \leq C_{1} \mathbf{x}^{T} B \mathbf{x}, \quad \forall \mathbf{x} \in \mathbb{R}^{M}
$$

Then for any $M \times M$ matrix $S$,

$$
\kappa\left(S A S^{T}\right) \leq \frac{C_{1}}{C_{0}} \kappa\left(S B S^{T}\right)
$$

Proof. Since

$$
\begin{aligned}
\lambda_{\max , S A S^{T}} & =\sup _{\mathbf{x} \in \mathbb{R}^{M}} \frac{\mathbf{x}^{T} S A S^{T} \mathbf{x}}{\mathbf{x}^{T} \mathbf{x}}=\sup _{\mathbf{x} \in \mathbb{R}^{M}} \frac{\left(S^{T} \mathbf{x}\right)^{T} A\left(S^{T} \mathbf{x}\right)}{\mathbf{x}^{T} \mathbf{x}} \\
& \leq C_{1} \sup _{\mathbf{x} \in \mathbb{R}^{M}} \frac{\mathbf{x}^{T} S B S^{T} \mathbf{x}}{\mathbf{x}^{T} \mathbf{x}}=C_{1} \lambda_{\max , S B S^{T}}
\end{aligned}
$$

Likewise we obtain

$$
\lambda_{\min , S A S^{T}} \geq C_{0} \lambda_{\min , S B S^{T}}
$$

and hence

$$
\kappa\left(S A S^{T}\right) \leq \frac{C_{1}}{C_{0}} \kappa\left(S B S^{T}\right)
$$

Lemma 2.2 tells that once one finds a good preconditioning matrix for $B$, then it is also a good preconditioning matrix for $A$ provided that the ratio $C_{1} / C_{0}$ is not large. Basic properties of positive definite matrices and their condition numbers maybe found in ([31], chapter 7).

Lemma 2.3. For $n=1,2, \ldots$, let $\chi_{n}:=\sum_{k=1}^{n} 2^{-k / 2} \chi_{\left[2^{-k}, 2^{-k+1}\right]}$, and $g_{n, j}:=$ $\chi_{n} \phi_{n, j}^{\prime}, \quad j=1, \ldots, 2^{n}-1$. Let $u=\sum_{j=1}^{2^{n}-1} c_{n, j} \phi_{n, j}$. Then

$$
\begin{equation*}
\int_{0}^{1}\left|\sum_{j=1}^{2^{n}-1} c_{n, j} g_{n, j}(x)\right|^{2} d x \leq\langle u, u\rangle_{E} \leq 2 \int_{0}^{1}\left|\sum_{j=1}^{2^{n}-1} c_{n, j} g_{n, j}(x)\right|^{2} d x \tag{2.12}
\end{equation*}
$$

where $\langle\cdot, \cdot\rangle_{E}$ is defined in (2.5) with $\alpha=1$.
Note: We may think of $g_{n, j}$ as the weighted derivative of $\phi_{n, j}$, and the weights are $2^{-k}, k=1, \ldots, n$, on the subintervals $\left(2^{-k}, 2^{-k+1}\right), k=n, \ldots, 1$. In other words, we discretize the weight $x$ in the inner product form $\langle\cdot, \cdot\rangle_{E}$ through $\chi_{n}$.

Proof. Noting that $g_{n, j}=0, j=1, \ldots, 2^{n}-1$ on $\left(0,2^{-n}\right)$, we get

$$
\langle u, u\rangle_{E}=\sum_{k=1}^{n} \int_{2^{-k}}^{2^{-k+1}} x\left|u^{\prime}(x)\right|^{2} d x .
$$

Accordingly,

$$
\langle u, u\rangle_{E}=\sum_{k=1}^{n} \int_{2^{-k}}^{2^{-k+1}} x\left|\sum_{j \in I_{k}} c_{n, j} \phi_{n, j}^{\prime}(x)\right|^{2} d x
$$

where $I_{k}:=\left\{2^{n-k}, \ldots, 2^{n-k+1}\right\}$.
Now we have

$$
\begin{align*}
\langle u, u\rangle_{E} & \geq \sum_{k=1}^{n} \int_{2^{-k}}^{2^{-k+1}} 2^{-k}\left|\sum_{j \in I_{k}} c_{n, j} \phi_{n, j}^{\prime}(x)\right|^{2} d x  \tag{2.13}\\
& =\sum_{k=1}^{n} \int_{2^{-k}}^{2^{-k+1}}\left|\sum_{j \in I_{k}}\left(c_{n, j} 2^{-\frac{k}{2}} \phi_{n, j}^{\prime}(x)\right)\right|^{2} d x \tag{2.14}
\end{align*}
$$

By the definition of $\phi_{n, j}$ in $(2.7,2.8)$, we have

$$
\phi_{n, j}^{\prime}(x)=\left\{\begin{array}{cl}
2^{n}, & (j-1) 2^{-n}<x<j 2^{-n} \\
-2^{n}, & j 2^{-n}<x<(j+1) 2^{-n} \\
0, & \text { otherwise }
\end{array}\right.
$$

Then, on each subinterval $\left(2^{-k}, 2^{-k+1}\right), k=n, n-1, \ldots, 1$, it follows that

$$
\sum_{j \in I_{k}} c_{n, j} 2^{-\frac{k}{2}} \phi_{n, j}^{\prime}=\sum_{j \in I_{k}} c_{n, j} \chi_{n} \phi_{n, j}^{\prime}=\sum_{j \in I_{k}} c_{n, j} g_{n, j} .
$$

This together with (2.14) yields

$$
\langle u, u\rangle_{E} \geq \int_{0}^{1}\left|\sum_{j=1}^{2^{n}-1}\left(c_{n, j} g_{n, j}(x)\right)\right|^{2} d x
$$

The proof of the right inequality of (2.12) is similar and is omitted.

Combining Lemma 2.3 with inequality (2.6), we have
Lemma 2.4. Denote by $A_{n}$ the matrix $\left(a\left(\phi_{n, j}, \phi_{n, l}\right)\right)_{j, l=1, \ldots, n}, A_{E, n}$ the matrix $\left(\left\langle\phi_{n, j}, \phi_{n, l}\right\rangle_{E}\right)_{j, l=1, \ldots, n}$ and by $\tilde{A}_{n}$ the matrix $\left(\left\langle g_{n, j}, g_{n, l}\right\rangle_{)_{j, l=1, \ldots, n}}\right.$. Then the inequalities

$$
\begin{equation*}
\mathbf{x}^{T} A_{E, n} \mathbf{x} \leq \mathbf{x}^{T} A_{n} \mathbf{x} \leq\left(1+\|q\|_{L_{\infty}(0,1)}\right) \mathbf{x}^{T} A_{E, n} \mathbf{x}, \quad \forall \mathbf{x} \in \mathbb{R}^{2^{n}-1} \tag{2.15}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathbf{x}^{T} \tilde{A}_{n} \mathbf{x} \leq \mathbf{x}^{T} A_{E, n} \mathbf{x} \leq 2 \mathbf{x}^{T} \tilde{A}_{n} \mathbf{x}, \quad \forall \mathbf{x} \in \mathbb{R}^{2^{n}-1} \tag{2.16}
\end{equation*}
$$

hold.

The following theorem is a simple consequence of Lemma 2.2 and Lemma 2.4.

Theorem 2.1. For any matrix $S$ of the same size as $A_{n}$,

$$
\kappa\left(S A_{n} S^{T}\right) \leq 2\left(1+\|q\|_{L_{\infty}(0,1)}\right) \kappa\left(S \tilde{A}_{n} S^{T}\right)
$$

By Theorem 2.1, we reduce the problem to preconditioning the much simpler matrix $\tilde{A}_{n}$ instead of $A_{n}$. Due to the similarity between the basis $\left\{g_{l, j}\right\}$ and the derivative of the basis $\left\{\phi_{n, i}\right\}$, it's natural to construct another orthogonal basis similar to the hierarchy basis to preconditioning $\tilde{A}_{n}$ (see, e.g., $[55,56])$. We will construct such a basis in the rest of this section.

Proposition 2.1. Let $\tilde{V}_{n}$ be the linear span of $g_{n, j}, j=1, \ldots, 2^{n-1}$. The sequence $\left\{\tilde{V}_{n}\right\}_{n=1,2, \ldots}$ of subspaces is nested, that is, $\tilde{V}_{n} \subset \tilde{V}_{n+1}$ for all $n$.

Proof. We shall show that the following relation is valid almost everywhere:

$$
g_{n-1, j}=\left\{\begin{array}{cl}
g_{n, 1}+g_{n, 2}+\frac{1}{2} g_{n, 3}, & j=1,  \tag{2.17}\\
\frac{1}{2} g_{n, 2 j-1}+g_{n, 2 j}+\frac{1}{2} g_{n, 2 j+1}, & j=2, \ldots, 2^{n-1}-1 .
\end{array}\right.
$$

For brevity, we define $\eta:=2^{-n}$. From the definition of $g_{n-1,1}$, we have

$$
g_{n-1,1}=\left\{\begin{array}{cl}
-2^{\frac{n-1}{2}}, & 2 \eta<x<4 \eta \\
0, & \text { otherwise }
\end{array}\right.
$$

Note that

$$
\begin{gathered}
g_{n, 1}(x)=\left\{\begin{array}{cl}
-2^{\frac{n}{2}}, & \eta<x<2 \eta \\
0, & \text { otherwise }
\end{array}\right. \\
g_{n, 2}(x)=\left\{\begin{array}{cl}
2^{\frac{n}{2}}, & \eta<x<2 \eta \\
-2^{\frac{n+1}{2}}, & 2 \eta<x<3 \eta \\
0, & \text { otherwise }
\end{array}\right.
\end{gathered}
$$

and

$$
g_{n, 3}(x)=\left\{\begin{array}{cl}
2^{\frac{n+1}{2}}, & 2 \eta<x<3 \eta \\
-2^{\frac{n+1}{2}}, & 3 \eta<x<4 \eta \\
0, & \text { otherwise }
\end{array}\right.
$$

Hence, for $x \in(0,4 \eta) \backslash\{\eta, 2 \eta, 3 \eta\}$, we have

$$
g_{n-1,1}(x)=g_{n, 1}(x)+g_{n, 2}(x)+\frac{1}{2} g_{n, 3}(x)
$$

To verify the second equation in (2.17), first we recall that

$$
\phi_{n-1, j}^{\prime}=\frac{1}{2} \phi_{n, 2 j-1}^{\prime}+\phi_{n, 2 j}^{\prime}+\frac{1}{2} \phi_{n, 2 j+1}^{\prime}, \text { a.e. } j=2, \ldots, 2^{n-1}-1 .
$$

Moreover, $\chi_{n-1}$ and $\chi_{n}$ agree on the interval $\left[2^{-n+1}, 1\right]$ and, for $j=2, \ldots, 2^{n-1}-$ $1, \phi_{n-1, j}$ is supported in $\left[2^{-n+1}, 1\right]$. Therefore,

$$
\begin{aligned}
g_{n-1, j} & =\chi_{n-1} \phi_{n-1, j}^{\prime}=\chi_{n} \phi_{n-1, j}^{\prime} \\
& =\chi_{n}\left(\frac{1}{2} \phi_{n, 2 j-1}^{\prime}+\phi_{n, 2 j}^{\prime}+\frac{1}{2} \phi_{n, 2 j+1}^{\prime}\right) \\
& =\frac{1}{2} g_{n, 2 j-1}+g_{n, 2 j}+\frac{1}{2} g_{n, 2 j+1}
\end{aligned}
$$

This proves the proposition.

Similar to the construction of the hierarchy basis, let

$$
\begin{equation*}
\tilde{\psi}_{l-1, j}:=g_{l, 2 j-1}, \quad j=1, \ldots, 2^{l-1}, l=n, n-1, \ldots, 1 \tag{2.18}
\end{equation*}
$$

and

$$
\tilde{W}_{l-1}:=\operatorname{span}\left\{\tilde{\psi}_{l-1, j}: j=1, \ldots, 2^{l-1}\right\}
$$

Then we have
Proposition 2.2. $\left\{\tilde{\psi}_{l, j}: l=1, \ldots, n, j=1, \ldots, 2^{l-1}\right\}$ is an orthogonal basis for $\tilde{V}_{n}$.

Proof. We shall verify the following properties:
i) $\left\langle\tilde{\psi}_{l-1, j}, g_{l-1, j^{\prime}}\right\rangle=0, \quad j=1, \ldots, 2^{l-1}, j^{\prime}=1, \ldots, 2^{l-1}-1$,
ii) $\left\langle\tilde{\psi}_{l-1, j}, \tilde{\psi}_{l-1, j^{\prime}}\right\rangle=0, \quad j \neq j^{\prime}$,
iii) $\tilde{V}_{n}=\tilde{W}_{0}+\tilde{W}_{1}+\cdots+\tilde{W}_{n-1}$.

Considering i), for $j \neq 1$, there exists $k$ such that $2 j-1 \in\left\{2^{l-k}+\right.$ $\left.1, \ldots, 2^{l-k+1}-1\right\}$, and

$$
\tilde{\psi}_{l-1, j}=g_{l, 2 j-1}=2^{l-\frac{k}{2}} \begin{cases}1, & (2 j-2) 2^{-l}<x<(2 j-1) 2^{-l} \\ -1, & (2 j-1) 2^{-l}<x<(2 j) 2^{-l} \\ 0, & \text { otherwise }\end{cases}
$$

Since $g_{l-1, j^{\prime}}$ is a constant on $\operatorname{supp}\left\{\tilde{\psi}_{l-1, j}\right\}=\left[(2 j-2) 2^{-l},(2 j) 2^{-l}\right]$ for $j^{\prime}=$ $\left.1, \ldots, 2^{l-1}-1, i\right)$ is true. For the case $j=1$, we obtain that $\tilde{\psi}_{l-1,1}\left(=g_{l, 1}\right)$ is orthogonal to $\tilde{V}_{l-1}$ because $g_{l-1, j^{\prime}}, j^{\prime}=1, \ldots, 2^{l-1}-1$, have no overlapped support with $\tilde{\psi}_{l-1,1}$.
ii) follows from

$$
\operatorname{supp}\left\{\tilde{\psi}_{l-1, j}\right\} \bigcap \operatorname{supp}\left\{\tilde{\psi}_{l-1, j^{\prime}}\right\}=\emptyset
$$

Finally, we turn to iii). First, $\left\{g_{l, j}\right\}_{j=1}^{2^{l}-1}$ is defined to be a basis for $\tilde{V}_{l}$.
Second, by i) and ii), we have

$$
\begin{equation*}
\tilde{V}_{n}=\tilde{V}_{1}+\tilde{W}_{1}+\tilde{W}_{2}+\cdots+\tilde{W}_{n-1} \tag{2.19}
\end{equation*}
$$

According to definitions, $\tilde{\psi}_{0,1}=g_{1,1}$ by (3.21), $\tilde{W}_{0}=\operatorname{span}\left\{\tilde{\psi}_{0,1}\right\}$, and $\tilde{V}_{1}=\operatorname{span}\left\{g_{1,1}\right\}$. Therefore, $\tilde{V}_{1}$ can be replaced by $\tilde{W}_{0}$ in (2.19).

This completes the proof.

In what follows we shall provide the preconditioning method for $\tilde{A}_{n}$ in (2.16). More precisely, we can find two sparse matrices $P$ and $H$ based on the change of bases from $\left\{g_{n, j}\right\}_{j}$ to $\left\{\tilde{\psi}_{l, j}\right\}_{l, j}$ such that $(P H) \tilde{A}_{n}(P H)^{T}$ is an identity matrix. By Theorem 2.1, it is clear that $(P H)$ is also a good preconditioner for the stiffness matrix $A_{n}$. To find the matrices $P$ and $H$, we shall write $(2.17,3.21)$ into the matrix form for the convenience of explanation.

Denote by $G_{l}, \tilde{\Psi}_{l}$ the vectors of functions $\left(g_{l, 1}, \ldots, g_{l, 2^{l}-1}\right)^{T},\left(\tilde{\psi}_{l, 1}, \ldots, \tilde{\psi}_{l, 2^{l}}\right)^{T}$, respectively. Let $\tilde{\Psi}:=\left(\tilde{\Psi}_{0}^{T}, \ldots, \tilde{\Psi}_{n-1}^{T}\right)^{T}$, and denote by $\tilde{A}_{\tilde{\Psi}, n}$ the matrix $\left\langle\tilde{\Psi},(\tilde{\Psi})^{T}\right\rangle$. Then $\tilde{A}_{\tilde{\Psi}, n}$ is a diagonal matrix by Proposition 2.2. Furthermore, one can find a diagonal matrix $P$ such that

$$
\begin{equation*}
I_{2^{n}-1} \equiv P \tilde{A}_{\tilde{\Psi}, n} P^{T} \tag{2.20}
\end{equation*}
$$

where $(l, l)$ entry of the matrix $P$ is $\|\tilde{\Psi}(l)\|_{L_{2}}^{-1}$, and $\tilde{\Psi}(l)$ denotes the $l$-th entry of the vector $\tilde{\Psi}$.

Clearly, $P \tilde{\Psi}$ is a stable (orthonormal) basis for $\tilde{V}_{n}$, and due to the simple transformation from the basis $G_{n}$ to the basis $P \tilde{\Psi}$ (see (2.17), (3.21)), $\tilde{A}_{n}$ can be preconditioned through a basis transformation from $G_{n}$ to $P \tilde{\Psi}$.

By (2.17), we have

$$
\begin{equation*}
G_{l-1}=B_{g, l-1} G_{l}, \tag{2.21}
\end{equation*}
$$

where $B_{g, l-1}$ is a $2^{l-1}-1$ by $2^{l}-1$ matrix (only nonzero entries are listed.)

$$
\left(\begin{array}{ccccccccc}
1 & 1 & \frac{1}{2} & & & & & & \\
& & \frac{1}{2} & 1 & \frac{1}{2} & & & & \\
& & & & & \ddots & & & \\
& & & & & & & & \\
& & & & & & \frac{1}{2} & 1 & \frac{1}{2}
\end{array}\right) .
$$

Denote by $B_{\tilde{\psi}, l-1}$ the $2^{l-1}$ by $2^{l}-1$ matrix

$$
\left(\begin{array}{ccccccc}
1 & 0 & 0 & & & & \\
0 & 0 & 1 & 0 & 0 & & \\
& & & & & \ddots & \\
& & & & & & 1
\end{array}\right)
$$

Then, (3.21) becomes

$$
\begin{equation*}
\tilde{\Psi}_{l-1}=B_{\bar{\psi}, l-1} G_{l} . \tag{2.22}
\end{equation*}
$$

Thus, (2.21) and (2.22) together yield

$$
\binom{G_{l-1}}{\tilde{\Psi}_{l-1}}=\binom{B_{g, l-1}}{B_{\tilde{\psi}, l-1}} G_{l} .
$$

By $H_{l-1}$ we denote the $2^{n}-1$ by $2^{n}-1$ transformation matrix

$$
H_{l-1}:=\left(\begin{array}{cc}
B_{g, l-1} & 0 \\
B_{\tilde{\psi}, l-1} & 0 \\
0 & I_{2^{n}-2^{l}}
\end{array}\right)
$$

Then we have

$$
\tilde{\Psi}=H_{1} \cdots H_{n-1} G_{n}
$$

For brevity, let

$$
H:=H_{1} \cdots H_{n-1},
$$

and thus we have the transformation between two bases

$$
\tilde{\Psi}=H G_{n}
$$

Note that $\tilde{A}_{\tilde{\Psi}, n}=H \tilde{A}_{n} H^{T}$. By (2.20), we have

$$
\begin{equation*}
I_{2^{n}-1} \equiv(P H) \tilde{A}_{n}(P H)^{T} \tag{2.23}
\end{equation*}
$$

Let $S$ in Theorem 2.1 be $P H$ in (2.23). Then

$$
\kappa\left(S A_{n} S^{T}\right) \leq 2\left(1+\|q\|_{L_{\infty}(0,1)}\right)
$$

Consequently, $(P H)$ is a suitable preconditioner for $A_{n}$. Furthermore, it's easily seen that $(P H)$ has $O(N)$ nonzero entries, where $N=2^{n}-1$ is the size of the basis functions for $V_{n}$. Therefore, implementation of the preconditioning shall be efficient. Detail discussion may be found in [52, Prop. 4.6].

Corollary 2.1. For the case $q(x)=0$, the condition number of the matrix $(P H) A_{n}(P H)^{T}$ is bounded by 2 for all $n$.

Now we provide a preconditioning algorithm for solving (2.11). Notice that

$$
A_{n} C_{n}=F_{n} \Leftrightarrow(P H) A_{n}(P H)^{T}\left((P H)^{T}\right)^{-1} C_{n}=(P H) F_{n}
$$

Then (2.11) is equivalent to the following linear equations with $\mathbf{x}=\left((P H)^{T}\right)^{-1} C_{n}$,

$$
\begin{equation*}
\left[(P H) A_{n}(P H)^{T}\right] \mathbf{x}=(P H) F_{n} \tag{2.24}
\end{equation*}
$$

To solve (2.11) for $C_{n}$, we first solve (2.24) for $\mathbf{x}$, and the solution of (2.11) is

$$
C_{n}=(P H)^{T} \mathbf{x}
$$

Note that the matrix $\left[(P H) A_{n}(P H)^{T}\right]$ is well conditioned. Therefore it's efficient to solve (2.24) for $\mathbf{x}$ numerically.

### 2.3 Error Estimates

We provide basic error estimates in this section and show that finite dimensional subspaces used in section 2.2 do provide the suitable approximation order.

To keep the practical applicability, and for the convenience of stating the results, we restrict ourselves to the uniform partition case in the previous section. Under such setting, it's easier to describe the preconditioning method based on the multi-level nested subspaces.

However, error estimates stated in this section hold for the general nonuniform partition case. Furthermore, the preconditioning method developed in the previous section is readily generalized to the non-uniform partition case as long as the sequence of subspaces are nested.

For the general non-uniform partition defined by $0=x_{0}<x_{1}<\ldots<x_{M}=$ 1, let

$$
\phi_{j}:=\left\{\begin{array}{cl}
\frac{x-x_{j-1}}{x_{j}-x_{j-1}}, & x \in\left[x_{j-1}, x_{j}\right], \\
\frac{x_{1+}-x}{x_{j+1}-x_{j}}, & x \in\left[x_{j}, x_{j+1}\right], \\
0, & \text { otherwise. }
\end{array}\right.
$$

We also let $h_{j}:=x_{j}-x_{j-1}$, and $h:=\max _{1 \leq j \leq M}\left\{h_{j}\right\}$, where the later quantity measures the mesh size. The finite dimensional space is spanned by the nodal basis functions $\left\{\phi_{j}\right\}$,

$$
V_{h}:=\operatorname{span}\left\{\phi_{0}+\phi_{1}, \phi_{2}, \ldots, \phi_{M-1}\right\} .
$$

Then the Galerkin method is to find $u_{h} \in V_{h}$ such that

$$
a\left(u_{h}, v\right)=\langle x f, v\rangle, \quad \forall v \in V_{h} .
$$

We will follow several lemmas to obtain error estimates in this section.
In the following, the solution $u$ is assumed to be smooth ( $u \in H^{2}$, where $H^{2}$ denotes the usual Sobolev space of functions with the second weak derivative on ( 0,1 ) ) with the boundary conditions $u^{\prime}(0)=u(1)=0$. We let the same letter $C$ which is independent of $h$ denote the different constants in the different inequalities.

Lemma 2.5. There exists a constant $h_{0}$ such that for all $h<h_{0}$,

$$
\left\|x^{1 / 2}\left(u^{\prime}-u_{I}^{\prime}\right)\right\|_{L_{2}} \leq C h \mid\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}},
$$

where $C$ is a constant depending on $\max _{i>1}\left\{x_{i} / x_{i-1}\right\}$ and that $u_{I} \in V_{h}$ is the interpolant of $u$ defined by

$$
u_{I}(x):=u\left(x_{1}\right)\left(\phi_{0}(x)+\phi_{1}(x)\right)+\sum_{i=2}^{M-1} u\left(x_{j}\right) \phi_{j}(x) .
$$

Proof. On the interval $I_{i}:=\left(x_{i-1}, x_{i}\right), i>1$, similar to the proof of Lemma 2 in [24], we have

$$
\int_{x_{i-1}}^{x_{i}} x\left(u^{\prime}-u_{I}^{\prime}\right)^{2} d x \leq x_{i} \int_{x_{i-1}}^{x_{i}}\left(u^{\prime}-u_{I}^{\prime}\right)^{2} d x
$$

By the well-known result (see, e.g., [7], p. 7)

$$
\left\|u^{\prime}-u_{I}^{\prime}\right\|_{L_{2}\left(I_{i}\right)} \leq C h\left\|u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)}
$$

it follows that

$$
\begin{aligned}
\int_{x_{i-1}}^{x_{i}} x\left(u^{\prime}-u_{I}^{\prime}\right)^{2} d x & \leq x_{i} C h^{2} \int_{x_{i-1}}^{x_{i}}\left(u^{\prime \prime}\right)^{2} d x \leq C h^{2} \frac{x_{i}}{x_{i-1}} \int_{x_{i-1}}^{x_{i}} x\left(u^{\prime \prime}\right)^{2} d x \\
& \leq C h^{2}\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)}^{2}
\end{aligned}
$$

which implies,

$$
\begin{equation*}
\left\|x^{1 / 2}\left(u^{\prime}-u_{I}^{\prime}\right)\right\|_{L_{2}\left(I_{i}\right)} \leq C h\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)} \tag{2.25}
\end{equation*}
$$

On the interval $I_{1}=\left(0, x_{1}\right)$, let $e(x):=u^{\prime}(x)-u^{\prime}\left(x_{1}\right)$. Then $e\left(x_{1}\right)=0$ and $e^{\prime}(x)=u^{\prime \prime}(x)$. Following the idea in Lemma 2.1, we have

$$
\begin{aligned}
\int_{0}^{x_{1}} x|e(x)|^{2} & =\int_{0}^{x_{1}} x\left|\int_{x}^{x_{1}} e^{\prime}(t) d t\right|^{2} d x \leq\left.\int_{0}^{x_{1}} x\left|\int_{x}^{x_{1}}\left(\frac{t}{x}\right)^{1 / 2}\right| e^{\prime}(t)\right|^{2} d t \\
& \leq\left\|t^{2 / 2} e^{\prime}\right\|_{L_{2}\left(I_{1}\right)}^{2} h^{2}
\end{aligned}
$$

and hence

$$
\begin{equation*}
\int_{0}^{x_{1}} x\left|u^{\prime}(x)-u^{\prime}\left(x_{1}\right)\right|^{2} d x \leq h^{2} \int_{0}^{x_{1}} x\left|u^{\prime \prime}(x)\right|^{2} d x \tag{2.26}
\end{equation*}
$$

Since

$$
\int_{0}^{x_{1}} x\left|u^{\prime}(x)-u^{\prime}\left(x_{1}\right)\right|^{2} d x=\Gamma_{1}+\Gamma_{2}-2 \int_{0}^{x_{1}} x u^{\prime}(x) u^{\prime}\left(x_{1}\right) d x
$$

where $\Gamma_{1}:=\int_{0}^{x_{1}} x\left|u^{\prime}(x)\right|^{2} d x$, and $\Gamma_{2}:=\left|u^{\prime}\left(x_{1}\right)\right|^{2} \int_{0}^{x_{1}} x d x$, we have

$$
\int_{0}^{x_{1}} x\left|u^{\prime}(x)-u^{\prime}\left(x_{1}\right)\right|^{2} d x \geq \Gamma_{1}-\alpha \Gamma_{1}-\frac{1}{\alpha} \Gamma_{2}+\Gamma_{2} .
$$

Let $\alpha=1 / 2$. Note that $u_{I}^{\prime}(x)=0, x \in I_{1}$. Then together with (2.26), we have

$$
\begin{equation*}
\left\|x^{1 / 2}\left(u^{\prime}-u_{I}^{\prime}\right)\right\|_{L_{2}\left(I_{1}\right)}=\left\|x^{1 / 2} u^{\prime}\right\|_{L_{2}\left(I_{1}\right)}^{2}=\Gamma_{1} \leq 2 h^{2}\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}\left(I_{1}\right)}^{2}+2 \Gamma_{2} . \tag{2.27}
\end{equation*}
$$

Combing (2.25) with (2.27) yields

$$
\left\|x^{1 / 2}\left(u^{\prime}-u_{I}^{\prime}\right)\right\|_{L_{2}(0,1)}^{2} \leq C h^{2}\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}(0,1)}^{2}+2 \Gamma_{2}
$$

The proof will be completed by estimating $\Gamma_{2}$ :

$$
2 \Gamma_{2}=\left|u^{\prime}\left(x_{1}\right)\right|^{2}\left(x_{1}\right)^{2} \leq h^{2}\left|u^{\prime}\left(x_{1}\right)\right|^{2}=h^{2}\left|\int_{0}^{x_{1}} u^{\prime \prime}(t) d t\right|^{2} \leq \frac{h^{3}}{2}\left\|u^{\prime \prime}\right\|_{L_{2}\left(I_{1}\right)}^{2}
$$

Hence there exists a constant $h_{0}$ such that

$$
h\left\|u^{\prime \prime}\right\|_{L_{2}\left(0, x_{1}\right)}^{2} \leq\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}(0,1)}^{2}, \quad \forall h<h_{0},
$$

and thus completes the proof.
Theorem 2.2. There exists a constant $h_{0}$ such that for any $h<h_{0}$,

$$
\begin{equation*}
\left\|x^{1 / 2}\left(u^{\prime}-u_{h}^{\prime}\right)\right\|_{L_{2}(0,1)} \leq C h\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}(0,1)} \tag{2.28}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}(0,1)} \leq C h^{2}\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}(0,1)} \tag{2.29}
\end{equation*}
$$

hold.

Proof. (2.28) is a standard error estimate.

$$
\begin{aligned}
\left\|x^{1 / 2}\left(u^{\prime}-u_{h}^{\prime}\right)\right\|_{L_{2}}^{2} & \leq C a\left(u-u_{h}, u-u_{h}\right) \leq C a\left(u-u_{h}, u-u_{I}\right) \\
& \leq C\left\|x^{1 / 2}\left(u^{\prime}-u_{h}^{\prime}\right)\right\|_{L^{2}}\left\|x^{1 / 2}\left(u^{\prime}-u_{I}^{\prime}\right)\right\|_{L_{2}}
\end{aligned}
$$

and hence,

$$
\left\|x^{1 / 2}\left(u^{\prime}-u_{h}^{\prime}\right)\right\|_{L_{2}} \leq C\left\|x^{1 / 2}\left(u^{\prime}-u_{I}^{\prime}\right)\right\|_{L_{2}} \leq C h\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}}
$$

where the last step used Lemma 2.5.
(2.29) can be obtained through a duality argument.

Let $w$ solves $a(v, w)=\left\langle x\left(u-u_{h}\right), v\right\rangle, \quad \forall v \in \dot{H}_{0}^{1}$. Then we have

$$
a\left(u-u_{h}, w\right)=\left\langle x\left(u-u_{h}\right), u-u_{h}\right\rangle=\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}}^{2},
$$

and

$$
\begin{aligned}
a\left(u-u_{h}, w\right) & =a\left(u-u_{h}, w-w_{I}\right) \leq C\left\|x^{1 / 2}\left(u^{\prime}-u_{h}^{\prime}\right)\right\|_{L_{2}}\left\|x^{1 / 2}\left(w^{\prime}-w_{I}^{\prime}\right)\right\|_{L_{2}} \\
& \leq C h^{2}\left\|x^{1 / 2} u^{\prime \prime}\right\|_{L_{2}}\left\|x^{1 / 2} w^{\prime \prime}\right\|_{L_{2}} .
\end{aligned}
$$

Once we prove the regularity of $w$, i.e.,

$$
\left\|x^{1 / 2} w^{\prime \prime}\right\|_{L_{2}} \leq C\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}}
$$

(2.29) holds.
$w$ satisfies the following equation, (from $\left.-\left(x w^{\prime}\right)^{\prime}=x\left(u-u_{h}-q w\right)\right)$

$$
w(x)=\int_{x}^{1} \frac{1}{t} \int_{0}^{t}\left(u(s)-u_{h}(s)-q(s) w(s)\right) s d s d t
$$

Differentiating both sides of the above expression twice, we have

$$
w^{\prime \prime}(x)=\frac{1}{x^{2}} \int_{0}^{x}\left(u(s)-u_{h}(s)-q(s) w(s)\right) s d s-\left(u(x)-u_{h}(x)\right)+q(x) w(x)
$$

and thus

$$
\begin{aligned}
\left\|x^{1 / 2} w^{\prime \prime}\right\|_{L_{2}} & \leq\left\|\frac{1}{x} \int_{0}^{x}\left(\frac{s}{x}\right)^{1 / 2} s^{1 / 2}\left(u-u_{h}-q w\right) d s\right\|_{L_{2}} \\
& +\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}}+C\left\|x^{1 / 2} w\right\|_{L_{2}}
\end{aligned}
$$

By the Hardy's inequality

$$
\left\|\frac{1}{x} \int_{0}^{x} f(t) d t\right\|_{L_{2}} \leq 2\|f\|_{L_{2}}
$$

we get

$$
\begin{aligned}
\left\|x^{1 / 2} w^{\prime \prime}\right\|_{L_{2}} & \leq\left\|\frac{1}{x} \int_{0}^{x}\left(s^{1 / 2}\left(u-u_{h}\right)-s^{1 / 2} q w\right) d s\right\|_{L_{2}}+\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}} \\
& +C\left\|x^{1 / 2} w\right\|_{L_{2}} \leq C\left[\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}}+\left\|x^{1 / 2} w\right\|_{L_{2}}\right]
\end{aligned}
$$

It remains to prove $\left\|x^{1 / 2} w\right\|_{L_{2}} \leq C\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}}$ :

$$
\begin{aligned}
\left\|x^{1 / 2} w\right\|_{L_{2}}\left\|x^{1 / 2}\left(u-u_{h}\right)\right\|_{L_{2}} & \geq\left\langle x\left(u-u_{h}\right), w\right\rangle=a(w, w) \\
& \geq C\left\|x^{1 / 2} w^{\prime}\right\|_{L_{2}}^{2} \geq C\left\|x^{1 / 2} w\right\|_{L_{2}}^{2}
\end{aligned}
$$

Finally, for the case $q(x)=0$, we provide error estimates for $\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L_{2}}$ and $\left\|u-u_{h}\right\|_{L_{2}}$.

Theorem 2.3. If $q(x)=0$, then

$$
\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L_{2}} \leq C h\left\|u^{\prime \prime}\right\|_{L_{2}} .
$$

Proof. $u_{h}$ satisfies the following equation,

$$
\begin{equation*}
\left\langle x\left(u^{\prime}-u_{h}^{\prime}\right), v^{\prime}\right\rangle=0, \quad \forall v \in V_{h} . \tag{2.30}
\end{equation*}
$$

In fact, $V_{h}^{\prime}:=\left\{v^{\prime}: v \in V_{h}\right\}$ is $\left\{w=\sum_{j=2}^{M} c_{j} \chi_{I_{j}}:\left\{c_{j}\right\} \in \mathbb{R}^{M-1}\right\}$. In other words, $V_{h}^{\prime}$ is the linear span of the piecewise constants on each interval $I_{j}, j>1$ with 0 on the interval $I_{1}$.

Let $u_{h}^{\prime}=\sum_{j=2}^{M} c_{j} \chi_{I_{j}}$. Then (2.30) is equivalent to

$$
\int_{x_{i-1}}^{x_{i}}\left(u^{\prime}(x)-c_{i}\right) x d x=0, \quad i>1
$$

Therefore,

$$
c_{i}=\int_{x_{i-1}}^{x_{i}} u^{\prime}(t) t d t / \int_{x_{i-1}}^{x_{i}} t d t
$$

Setting $\Gamma:=\int_{x_{i-1}}^{x_{i}} t d t$, we obtain
$(2.31)\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L_{2}\left(I_{i}\right)}^{2}=\int_{x_{i-1}}^{x_{i}}\left(u^{\prime}(x)-u_{h}^{\prime}(x)\right)^{2} d x$

$$
=\int_{x_{i-1}}^{x_{i}}\left(\int_{x_{i-1}}^{x_{i}}\left(u^{\prime}(x)-u^{\prime}(t)\right) t d t\right)^{2} d x / \Gamma^{2}, \quad i>1
$$

Estimating the last term in the above equation gives
(2.32)
$\int_{x_{i-1}}^{x_{i}}\left(\int_{x_{i-1}}^{x_{i}}\left(u^{\prime}(x)-u^{\prime}(t)\right) t d t\right)^{2} d x \leq \int_{x_{i-1}}^{x_{i}}\left(\int_{x_{i-1}}^{x_{i}}\left|u^{\prime}(x)-u^{\prime}(t)\right|^{2} d t \int_{x_{i-1}}^{x_{i}} t^{2} d t\right) d x$.
Now we have

$$
\begin{align*}
\int_{x_{i-1}}^{x_{i}} & \left(u^{\prime}(x)-u^{\prime}(t)\right)^{2} d t=\int_{x_{i-1}}^{x_{i}}\left|\int_{x}^{t} u^{\prime \prime}(s) d s\right|^{2} d t  \tag{2.33}\\
& \leq\left.\int_{x_{i-1}}^{x_{i}}\left|\int_{x}^{t}\right| u^{\prime \prime}(s)\right|^{2} d s| | \int_{x}^{t} 1^{2} d s \mid d t \leq\left\|u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)}^{2} h_{i}^{2}
\end{align*}
$$

Hence, after plugging (2.34) into (2.32), we estimate the term $\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L_{2}\left(I_{i}\right)}^{2}$ in (2.32) by

$$
\begin{aligned}
\int_{x_{i-1}}^{x_{i}}\left|u^{\prime}(x)-u_{h}^{\prime}(x)\right|^{2} d x & \leq h_{i}^{3}\left\|u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)}^{2} \frac{\int_{x_{i-1}}^{x_{i}} t^{2} d t}{\left(\int_{x_{i-1}}^{x_{i}} t d t\right)^{2}} \\
& \leq h_{i}^{3}\left\|u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)}^{2} \frac{x_{i}^{2} \int_{x_{i-1}}^{x_{i}} d t}{\left(x_{i-1} \int_{x_{i-1}}^{x_{i}} d t\right)^{2}} \leq\left(\frac{x_{i}}{x_{i-1}}\right)^{2} h^{2}\left\|u^{\prime \prime}\right\|_{L_{2}\left(I_{i}\right)}^{2} .
\end{aligned}
$$

On the interval $I_{1}$,

$$
\begin{aligned}
\int_{0}^{x_{1}}\left|u^{\prime}(x)\right|^{2} d x & =\int_{0}^{x_{1}}\left|\int_{0}^{x} u^{\prime \prime}(s) d s\right|^{2} d x \\
& \leq\left.\int_{0}^{x_{1}}\left|\int_{0}^{x}\right| u^{\prime \prime}(s)\right|^{2} d s \int_{0}^{x} 1^{2} d s \mid d x \leq\left\|u^{\prime \prime}\right\|_{L_{2}\left(0, x_{1}\right)}^{2} h^{2}
\end{aligned}
$$

Thus the proof is completed.

Theorem 2.4. If $q(x)=0$, there exists a constant $h_{0}$ such that for any $h<h_{0}$,

$$
\left\|u-u_{h}\right\|_{L_{2}(0,1)} \leq C h\left\|u^{\prime}-u_{h}^{\prime}\right\|_{L_{2}(0,1)} .
$$

Proof. Let $e:=u-u_{h}$, and let $w$ solve the problem,

$$
-\left(x w^{\prime}\right)^{\prime}=e, \quad \text { with } \quad w(1)=w^{\prime}(0)=0
$$

Then

$$
\begin{equation*}
w(x)=\int_{x}^{1} \frac{1}{t} \int_{0}^{t} e(s) d s d t \tag{2.34}
\end{equation*}
$$

Differentiating (2.34) twice, it follows that

$$
w^{\prime \prime}(x)=\frac{1}{x^{2}} \int_{0}^{x} e(s) d s-\frac{1}{x} e(x)
$$

By using Hardy's inequality, we have

$$
\begin{equation*}
\left\|x w^{\prime \prime}\right\|_{L_{2}} \leq 2\|e\|_{L_{2}}+\|e\|_{L_{2}}=3\|e\|_{L_{2}} . \tag{2.35}
\end{equation*}
$$

On the other hand,

$$
\|e\|_{L_{2}}^{2}=\langle e, e\rangle=-\int_{0}^{1}\left(x w^{\prime}\right)^{\prime} e
$$

Integrating it by parts yields

$$
\|e\|_{L_{2}}^{2}=\int_{0}^{1} x w^{\prime} e^{\prime} d x=\int_{0}^{1} x\left(w^{\prime}-v^{\prime}\right) e^{\prime} d x, \quad \forall v \in V_{h}
$$

which implies

$$
\|e\|_{L_{2}}^{2} \leq\left\|e^{\prime}\right\|_{L_{2}} \inf _{v \in V_{h}}\left\|x\left(w^{\prime}-v^{\prime}\right)\right\|_{L_{2}}
$$

Suppose we have the result $\inf _{v \in V_{h}}\left\|x\left(w^{\prime}-v^{\prime}\right)\right\|_{L_{2}} \leq C h\left\|x w^{\prime \prime}\right\|_{L_{2}}$ for the moment. Then together with (2.35), we have

$$
\|e\|_{L_{2}}^{2} \leq C h\|e\|_{L_{2}}\left\|e^{\prime}\right\|_{L_{2}}
$$

Thus completes the proof.

It remains to prove

Lemma 2.6. For an element $w \in \dot{H}_{0}^{1} \cap H^{2}, w^{\prime}(0)=0$, there exists a constant $h_{0}$ such that for any $h<h_{0}$,

$$
\inf _{v \in V_{h}}\left\|x\left(w^{\prime}-v^{\prime}\right)\right\|_{L_{2}} \leq C h\left\|x w^{\prime \prime}\right\|_{L_{2}}
$$

Proof. Let $v$ be an element in $V_{h}$. Then there exist $c_{2}, \ldots, c_{M}$ such that $v=$ $\sum_{j=2}^{M} c_{j} \chi_{I_{j}}$. Let

$$
c_{j}=\int_{I_{j}} x w^{\prime}(x) d x / \int_{I_{j}} x d x
$$

Then a similar approach in Theorem 2.3 can be taken to obtain

$$
\left\|x\left(w^{\prime}-v^{\prime}\right)\right\|_{L_{2}}^{2} \leq C h\left\|x w^{\prime \prime}\right\|_{L_{2}}^{2}
$$

Detail steps are as follows:
On $I_{j}, j>1$, let $\Gamma:=\int_{I_{j}} x d x$. It follows that

$$
\begin{aligned}
\int_{I_{j}}\left|x\left(w^{\prime}-v^{\prime}\right)\right|^{2} d x & =\int_{I_{j}}\left|x \int_{x_{j-1}}^{x_{j}}\left(w^{\prime}(x)-w^{\prime}(t)\right) t d t\right|^{2} d x /\left(\Gamma^{2}\right) \\
& \leq \frac{1}{\Gamma^{2}} \int_{I_{j}}\left|x \int_{x_{j-1}}^{x_{j}} t \int_{t}^{x} w^{\prime \prime} d s d t\right|^{2} d x \\
& \leq C_{1} \frac{1}{\Gamma^{2}} \int_{I_{j}}\left|x \int_{x_{j-1}}^{x_{j}}\right| \int_{t}^{x}\left|w^{\prime \prime}(s) s\right| d s|d t|^{2} d x \\
& \left.\leq C_{1} \frac{1}{\Gamma^{2}} \int_{I_{j}} \right\rvert\, x\left(\left.\left\|s w^{\prime \prime}\right\|_{L_{2}} \sqrt{h_{j}} h_{j}\right|^{2} d x\right. \\
& =C_{1} h_{j}^{3}\left\|s w^{\prime \prime}\right\|_{L_{2}}^{2} \frac{\int_{I_{j}} x^{2} d x}{\Gamma^{2}} \\
& \leq C h^{2}\left\|s w^{\prime \prime}\right\|_{L_{2}}^{2}
\end{aligned}
$$

where $C_{1}=x_{j} / x_{j-1}\left(\right.$ since $\left.\frac{s}{t} \leq \frac{x_{j}}{x_{j-1}}\right)$.
On the interval $I_{1}$, we use the trick in Lemma 2.5 again, and we obtain

$$
\begin{aligned}
\int_{0}^{x_{1}}\left|x\left(w^{\prime}(x)-w^{\prime}\left(x_{1}\right)\right)\right|^{2} d x & =\int_{0}^{x_{1}}\left|x \int_{x}^{x_{1}} w^{\prime \prime}(t) d t\right|^{2} d x \\
& \leq \int_{0}^{x_{1}}\left|x \int_{x}^{x_{1}}\right| w^{\prime \prime}(t)(t / x)|d t|^{2} d x \\
& \leq \int_{0}^{x_{1}}\left|\int_{x}^{x_{1}}\right| w^{\prime \prime}(t) t|d t|^{2} d x \\
& \leq h^{2}| | w^{\prime \prime}(t) t \|_{L_{2}\left(I_{1}\right)}^{2}
\end{aligned}
$$

Let $\Gamma_{1}:=\int_{0}^{x_{1}}\left|x w^{\prime}(x)\right|^{2} d x$, and $\Gamma_{2}:=\int_{0}^{x_{1}}\left|x w^{\prime}\left(x_{1}\right)\right|^{2} d x=\frac{x_{1}^{3}}{3}\left|w^{\prime}\left(x_{1}\right)\right|^{2}$. Then we have

$$
\Gamma_{1} \leq 2 h^{2}\left\|t w^{\prime \prime}(t)\right\|_{L_{2}\left(I_{1}\right)}^{2}+2 \Gamma_{2} .
$$

Estimating $2 \Gamma_{2}$ gives

$$
2 \Gamma_{2}=\frac{2}{3} x_{1}^{3}\left|w^{\prime}\left(x_{1}\right)\right|^{2} \leq \frac{2}{3} h^{3}\left(\int_{0}^{x_{1}} w^{\prime \prime}(t) d t\right)^{2} \leq \frac{2}{3} h^{2}\left(\left\|w^{\prime \prime}\right\|_{L_{2}\left(I_{1}\right)}^{2} h^{2}\right) .
$$

Hence there exists $h_{0}$ depending on $w$, such that for any $h<h_{0}$,

$$
\left\|w^{\prime \prime}\right\|_{L_{2}\left(I_{1}\right)}^{2} h^{2} \leq\left\|x w^{\prime \prime}\right\|_{L_{2}(0,1)}^{2}
$$

This completes the proof.

### 2.4 Examples

Now we shall provide the results of numerical examples to illustrate the theory developed earlier.

For the first example, let $q(x)=0$ and

$$
f(x)=-\frac{\pi}{2} x^{-1}\left(\sin \frac{\pi}{2} x+x \frac{\pi}{2} \cos \frac{\pi}{2} x\right)
$$

Then the solution is

$$
u=-\cos \frac{\pi}{2} x
$$

To define the coarsest grid, we split the domain $(0,1)$ into $2^{1}=2$ pieces, and then divide each piece into 2 pieces of equal length. We keep on splitting until there are $2^{n}$ pieces. Therefore, we have $n$ levels of nested subspaces.

Concerning the singular boundary value problem (2.4), once the finite dimensional subspace $V_{n}$ is fixed, the stiffness matrix $A_{n}$ is also fixed, as well as the preconditioner $P H$. Therefore, we will demonstrate the performance of the preconditioner first by comparing the condition numbers of $A_{n}$ with those of $(P H) A_{n}(P H)^{T}$ with different $n$. In Table 2.1, we display the maximum eigenvalues, minimum eigenvalues and the condition numbers of the two matrices $A_{n}$ and $(P H) A_{n}(P H)^{T}$ for the different $n$.

| $n$ | 5 | 6 | 7 | 8 | 9 | 10 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max , A_{n}}$ | $1.10+2$ | $2.33+2$ | $4.83+2$ | $9.87+2$ | $2.00+3$ | $4.04+3$ |
| $\lambda_{\min , A_{n}}$ | 0.0477 | 0.0232 | 0.0115 | 0.0057 | 0.0028 | 0.0014 |
| $\kappa\left(A_{n}\right)$ | $2.30+3$ | $1.00+4$ | $4.21+4$ | $1.74+5$ | $7.06+5$ | $2.85+6$ |
| $\lambda_{\max , P H A_{n}(P H)^{T}}$ | 1.9688 | 1.984 | 1.992 | 1.996 | 1.998 | 1.999 |
| $\lambda_{\min , P H A_{n}(P H)^{T}}$ | 1.031 | 1.016 | 1.008 | 1.004 | 1.003 | 1.001 |
| $\kappa\left(P H A_{n}(P H)^{T}\right)$ | 1.909 | 1.954 | 1.977 | 1.988 | 1.994 | 1.997 |

Table 2.1: Condition numbers of the matrix $A_{n}$ and $(P H) A_{n}(P H)^{T}$

Computing results in Table 2.1 show that condition numbers of the preconditioned stiffness matrices are uniformly bounded by 2 , which verifies the result of the Corollary in section 2.2.

Now, we are ready to implement our preconditioning method to solve the first example. We use the Galerkin method to solve the problem with mesh size $1 / 2^{n}$ and let $u$ and $u_{h}$ denote the solution and the Galerkin solution of the singular problem, respectively. $\left|\left(u-u_{h}\right)\right|_{H^{1}},\left\|\left(u-u_{h}\right)\right\|_{L_{2}}$ with different $n$ are listed in Table 2.2. As predicted by Theorem 2.3, 2.4, the Galerkin method with the piecewise linear nodal basis preserves $O(h), O\left(h^{2}\right)$ convergence rates for $\left|\left(u-u_{h}\right)\right|_{H^{1}},\left\|\left(u-u_{h}\right)\right\|_{L_{2}}$, respectively.

| n | $\left\|u-u_{h}\right\|_{H^{1}}$ | $\left\\|u-u_{h}\right\\|_{L_{2}}$ |
| :---: | :---: | :---: |
| 5 | 0.0307 | $1.73 \times 10^{-4}$ |
| 6 | 0.0134 | $4.53 \times 10^{-5}$ |
| 7 | 0.0062 | $1.20 \times 10^{-5}$ |
| 8 | 0.0029 | $3.16 \times 10^{-6}$ |
| 9 | 0.0014 | $8.17 \times 10^{-7}$ |
| 10 | $7.05 \times 10^{-4}$ | $2.12 \times 10^{-7}$ |

Table 2.2: Estimates of $\left|u-u_{h}\right|_{H^{1}},\left|\left|u-u_{h}\right|_{L_{2}}\right.$

After we apply the weighted Jacob iterative method with $\omega=2 / 3$ (see, e.g., [8]), and an initial guess of 0 , to the preconditioned linear system for $10,15,20$ iterative times, we obtain the numerical solutions $u_{i t}$. The errors between the numerical solution and the Galerkin solution, estimated in $H^{1}$ semi-norm and $L_{2}$ norm, are given in Table 2.3.

|  | $\left\|u_{i t}-u_{h}\right\|_{H^{1}}$ |  |  | $\left\\|u_{i t}-u_{h}\right\\|_{L_{2}}$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| n | itno=10 | itno=15 | itno=20 | itno=10 | itno=15 | itno $=20$ |
| 5 | $4.00-4$ | $1.02-5$ | $2.60-7$ | $1.31-4$ | $3.28-6$ | $8.34-8$ |
| 6 | $4.05-4$ | $1.04-5$ | $2.67-7$ | $1.38-4$ | $3.48-6$ | $8.91-8$ |
| 7 | $4.06-4$ | $1.05-5$ | $2.69-7$ | $1.41-4$ | $3.57-6$ | $9.14-8$ |
| 8 | $4.07-4$ | $1.05-5$ | $2.70-7$ | $1.43-4$ | $3.61-6$ | $9.25-8$ |
| 9 | $4.07-4$ | $1.05-5$ | $2.70-7$ | $1.43-4$ | $3.63-6$ | $9.30-8$ |
| 10 | $4.07-4$ | $1.05-5$ | $2.70-7$ | $1.43-4$ | $3.64-6$ | $9.32-8$ |

Table 2.3: Estimates of $\left|u_{i t}-u_{h}\right|_{H^{1}},\left\|u_{i t}-u_{h}\right\|_{L_{2}}$ with the different iterative numbers

Since the condition number is strictly bounded by 2 , it can be expected that the convergence rate of $u_{i t}-u_{h}$ shall be like $O\left(\rho^{m}\right)$, where iterative number is denoted by $m$, and $\rho<1$. Numerical results in Table 2.3 indicate that $\rho \approx \frac{1}{2}$. In fact, a careful analysis of iterative methods with a given condition number bound may give an estimate of the convergence rate, which shall not be discussed in detail here.

Similar to the example shown in [24], let $q(x)=1-x^{2}, f(x)=\left(1-x^{2}\right)^{2}+4$ in our second example. Subspace level is up to $n=10$. In this situation, $u(x)=1-x^{2}$, and $\left|\left(u-u_{h}\right)\right|_{H^{1}},\left\|\left(u-u_{h}\right)\right\|_{L_{2}}$ are computed in Table 2.4 for different $n$. Condition numbers of the preconditioned system are shown in Table 2.5. Similar computing results to example one are obtained. These numerical results confirm the performance of our preconditioning method.

| n | $\left\|u-u_{h}\right\|_{H^{1}}$ | $\left\\|u-u_{h}\right\\|_{L_{2}}$ |
| :---: | :---: | :---: |
| 5 | 0.0306 | $2.07 \times 10^{-4}$ |
| 6 | 0.0141 | $5.49 \times 10^{-5}$ |
| 7 | 0.0067 | $1.44 \times 10^{-5}$ |
| 8 | 0.0033 | $3.70 \times 10^{-6}$ |
| 9 | 0.0016 | $9.45 \times 10^{-7}$ |
| 10 | $8.03 \times 10^{-4}$ | $2.39 \times 10^{-7}$ |

Table 2.4: Estimates of $\left|u-u_{h}\right|_{H^{1}},\left\|u-u_{h}\right\|_{L_{2}}$

| $n$ | 5 | 6 | 7 | 8 | 9 | 10 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\lambda_{\max , A_{n}}$ | $1.10+2$ | $2.33+2$ | $4.83+2$ | $9.87+2$ | $2.00+3$ | $4.04+3$ |
| $\lambda_{\min , A_{n}}$ | $5.38-2$ | $2.61-2$ | $1.29-2$ | $6.4-3$ | $3.2-3$ | $1.6-3$ |
| $\kappa\left(A_{n}\right)$ | $2.04+3$ | $8.91+3$ | $3.75+4$ | $1.54+5$ | $6.28+5$ | $2.54+6$ |
| $\lambda_{\max , P H A_{n}(P H)^{T}}$ | 2.016 | 2.021 | 2.023 | 2.024 | 2.024 | 2.024 |
| $\lambda_{\min , P H A_{n}(P H)^{T}}$ | 1.036 | 1.018 | 1.009 | 1.004 | 1.002 | 1.001 |
| $\kappa\left(P H A_{n}(P H)^{T}\right)$ | 1.946 | 1.986 | 2.005 | 2.015 | 2.019 | 2.022 |

Table 2.5: Condition numbers of the matrix $A_{n}$ and $(P H) A_{n}(P H)^{T}$

|  | $\left\|u_{i t}-u_{h}\right\|_{H^{1}}$ |  |  | $\left\\|u_{i t}-u_{h}\right\\|_{L_{2}}$ |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| n | itno $=10$ | itno=15 | itno=20 | itno=10 | itno=15 | itno $=20$ |
| 5 | $3.89-4$ | $1.04-5$ | $2.81-7$ | $1.10-4$ | $2.78-6$ | $7.19-8$ |
| 6 | $3.92-4$ | $1.05-5$ | $2.86-7$ | $1.14-4$ | $2.89-6$ | $7.48-8$ |
| 7 | $3.93-4$ | $1.05-5$ | $2.87-7$ | $1.16-4$ | $2.93-6$ | $7.59-8$ |
| 8 | $3.93-4$ | $1.06-5$ | $2.87-7$ | $1.16-4$ | $2.95-6$ | $7.64-8$ |
| 9 | $3.93-4$ | $1.06-5$ | $2.87-7$ | $1.17-4$ | $2.96-6$ | $7.66-8$ |
| 10 | $3.93-4$ | $1.06-5$ | $2.87-7$ | $1.17-4$ | $2.97-6$ | $7.67-8$ |

Table 2.6: Estimates of $\left|u_{i t}-u_{h}\right|_{H^{1}}, \| u_{i t}-\left.u_{h}\right|_{L_{2}}$ with the different iterative numbers

## Chapter 3

## $C^{1}$ wavelets on two dimensional triangular meshes

### 3.1 Introduction

We are interested in the multilevel analysis on a sequence of nested subspaces

$$
V_{0} \subset V_{1} \subset V_{2} \cdots
$$

on triangulations $\left\{\mathcal{T}_{k}\right\}_{k=0}^{\infty}$ of a polygonal domain $\Omega$ in $\mathbb{R}^{2}$. The multilevel decomposition of $\left\{V_{k}\right\}_{k=0}^{\infty}$ is to seek the proper subspaces that

$$
V_{J}=V_{0}+\widetilde{W}_{1}+\widetilde{W}_{2}+\cdots+\widetilde{W}_{J-1}+\widetilde{W}_{J}
$$

$\widetilde{W}_{k} \subset V_{k}$ is chosen to be orthogonal to $V_{k-1}$ with respect to some kind of inner product. The basis functions for $\widetilde{W}_{k}$ are generally called wavelets. If the usual $L_{2}$ inner product $\langle\cdot, \cdot\rangle_{L_{2}}$ on $\Omega$ is applied, then basis functions, say $\left\{\psi_{j}^{k}\right\}_{j \in I_{\psi, k}}$ for $\widetilde{W}_{k}$, are traditionally called semi-wavelets $[14,16]$, where $I_{\psi, k}$ denotes the index set of wavelets. Let $\left\{\phi_{j}^{k}\right\}_{j \in I_{k}}$ be basis functions for $V_{k}$. Then

$$
\left\langle\psi_{j}^{k}, \phi_{j^{\prime}}^{k^{\prime}}\right\rangle_{L_{2}}=0, \quad k^{\prime}<k, j \in I_{k}, j^{\prime} \in I_{\psi, k^{\prime}}
$$

A basis $\left\{\phi_{j}\right\}_{j=1}^{\infty}$ is stable if we have

$$
\left\|\sum_{j=1}^{\infty} \alpha_{j} \phi_{j}\right\|^{2} \simeq \sum_{j=1}^{\infty} \alpha_{j}^{2}
$$

where $\|\cdot\|$ is the norm of interest. Here,$\simeq$ refers to that two terms can be bounded by some constant multiple of each other, with the constant independent of the parameters on which these two terms may depend. Similarly, We let $\lesssim(\gtrsim)$ denote that the first (second) term can be bounded by a constant multiple of the second (first) term.

Moreover, $\left\{\phi_{j}\right\}_{j=1}^{\infty}$ is a Bessel sequence if and only if [16]

$$
\left\|\sum_{j=1}^{\infty} \alpha_{j} \phi_{j}\right\|^{2} \lesssim \sum_{j=1}^{\infty} \alpha_{j}^{2}
$$

A well-known norm equivalence theory in literature (for example, see [21, 54]) reads

$$
\begin{equation*}
\|u\|_{H^{s}}^{2} \simeq \sum_{k=0}^{J} 2^{2 k s}\left\|u_{k}\right\|_{L_{2}}^{2} \tag{3.1}
\end{equation*}
$$

where

$$
u=\sum_{k=0}^{J} u_{k}, u_{k} \in \widetilde{W}_{k}, \widetilde{W}_{0} \equiv V_{0}
$$

Consequently, if the (semi-)wavelet basis $\left\{\psi_{j}^{k}\right\}_{j \in I_{\psi, k}}$ for $\widetilde{W}_{k}$ is stable in the $L_{2}$ space, then wavelet system $\left\{2^{-k s} \psi_{j}^{k}\right\}_{k=0,1, \ldots, j \in I_{\psi, k}}$ forms a stable basis in the Sobolev space $H^{s}[1,10]$, where $s$ is some positive real number.

Stable bases in Sobolev spaces have broad applications in numerically solving partial differential equations as well as integral equations [12, 13, 11, 19, $21,34,38,45,52]$. For higher order problems, wavelets with higher order smoothness, such as $C^{1}$ wavelets, are required. This motivates people to construct the stable wavelet systems in Sobolev spaces. In particular, some stable wavelets in Sobolev spaces have been constructed and discussed on the uniform meshes in [39, 41, 50, 51].

Validation of (3.1) requires some mild conditions on the subspaces $\left\{V_{k}\right\}_{k=0}^{\infty}$, i.e., the Jackson inequality and the Bernstein inequality. In general, the Jackson inequality implies the suitable approximation capability provided by the subspaces $\left\{V_{k}\right\}_{k=0}^{\infty}$, and the Bernstein inequality, known as the inverse inequality for a long time in finite element literature, usually holds true for the underlying subspaces $\left\{V_{k}\right\}_{k=0}^{\infty}$.

Floater and Quak have constructed the semi-wavelets from the piecewise linear nodal basis with small supports on the irregular meshes in [25, 26, 27]. Their wavelets are orthogonal among different levels with respect to the $L_{2}$ inner product. By (3.1), it is clear that their wavelets also form a stable basis in the Sobolev space after properly scaled. To reduce the support of wavelet, Stevenson introduced the discrete $L_{2}$ inner product in [46, 47] for each subspace $V_{k}$, which is also generated from the piecewise linear nodal basis. $\widetilde{W}_{k}$ is then the orthogonal complement of $V_{k-1}$ in $V_{k}$ with respect to the discrete inner product for $V_{k}$. His wavelet for the linear nodal basis for two dimensional case is simple and has three coefficients. Numerical results on the regular mesh (see [39]) show the potential for numerically solving partial differential equations. The idea of constructing wavelets which are not $L_{2}$ orthogonal among levels can also be found in several other papers (see [48, 49]).
$C^{1}$ wavelets on general meshes are of particular importance for fulfilling the smoothness condition required by numerically solving high order problems. However, the difficulty in choosing the proper $C^{1}$ scaling functions on general meshes as well as the lack of general theory to verify the stability of wavelets make it intensely difficult to construct stable $C^{1}$ wavelets in Sobolev spaces. Main goal of this chapter is to investigate the general theory on the construction of $C^{1}$ continuous wavelets on the general meshes. We wish to develop the theory independent on the Fourier analysis, which is considered as the fundamental theoretical tool for the wavelet analysis on uniform meshes. On the other hand, to further generalize Stevenson's idea to more complicated finite element spaces may not be suitable because one key estimate of the projection operator (Theorem 4.3 in [47]) involves three levels of subspaces. In
this chapter, we shall overcome this weakness by estimating the same operator only involving two levels. Our method greatly reduces the complexity in the estimate and thus makes the estimate for more complicated basis functions possible. Moreover, our wavelets are extremely short supported, and this leads to the fast algorithm for applications.

We will take Powell-Sabin elements (PS element, see [42]) as basis functions for subspaces $\left\{V_{k}\right\}_{k=0}^{\infty}$. On the three-direction meshes, where grid lines run only in three directions (see [36] and [3], p.294), PS element is $C^{1}$ continuous on the meshes and the subspaces $\left\{V_{k}\right\}_{k=0}^{\infty}$ are nested. We design the discrete $L_{2}$ inner product $\langle\cdot, \cdot\rangle_{D_{2}(k)}$ for each subspace $V_{k}$ and our wavelet subspace $\widetilde{W}_{k}$ is orthogonal to $V_{k-1}$ with respect to the discrete inner product $\langle\cdot, \cdot\rangle_{D_{2}(k)}$. One improvement in our estimate is that the proof is proceeded on the standard equilateral triangle in stead of an arbitrary triangle due to the affine invariance property of the basis functions. The stability of the wavelet basis in the Sobolev space $H^{2}$ is verified.

Although the three-direction meshes are a type of regular triangular meshes, our theory does not employ any property of the regular structure of meshes. The obstacle preventing us from applying our theory to irregular meshes is the difficulty in seeking the nested subspaces from suitable $C^{1}$ continuous refinable functions. However, our theory is sufficiently flexible to be extendable to irregular meshes once proper $C^{1}$ continuous refinable basis functions on irregular meshes are obtained.

This chapter is divided into four parts. In Section 3.2, we introduce the knowledge of Powell-Sabin $C^{1}$ element which is employed later as scaling functions to produce the sequence of the nested subspaces. We construct the wavelets in Section 3.3, and some basic properties of the wavelets are presented. More precisely, we define one discrete $L_{2}$ inner product for each subspace, and on the basis of such inner products, we construct the wavelets to form a Riesz basis in Sobolev spaces. In Section 3.4, we focus on the general theory of Riesz basis in the Sobolev space. We shall verify that our wavelets satisfy all required conditions to form a Riesz basis in the Sobolev space $H^{2}$.

The last section is devoted to the examples of the computation for wavelets on the regular triangular meshes.

### 3.2 PS element and $C^{1}$ basis functions

Powell-Sabin (PS) element has been traced back to 1977 and was widely used in CAGD because of its simple structure (compared with other $C^{1}$ elements). On a three-direction mesh, it is refinable and thus the sequence of subspaces from PS element is nested. On a general mesh, PS element is no longer guaranteed to be $C^{1}$ continuous globally. Applications utilizing these nonconforming PS element were discussed in [40]. In this chapter, we use PS element as the scaling functions to generate the sequence of nested subspaces on two dimensional three-direction meshes.

The structure of this section is organized as follows. In Section 3.2.1, we give the definition of PS element and study the local basis functions on an arbitrary triangle. In connection with the local basis functions on the standard equilateral triangle of edge length one, we discuss the affine map and affine invariance properties of the local basis in Section 3.2.2. By (3.10), we shows that affine invariance property allows us computing independently on the shapes of the triangles and thus it provides an efficient way to carry out computing and constructing wavelets. In Section 3.2.3, we combine local basis functions together to get $C^{1}$ continuous basis functions on a three-direction mesh. Associated with such basis, discrete $L_{2}$ inner products and norms are introduced.

### 3.2.1 Powell-Sabin element

PS element is composed of piecewise quadratic polynomials on an arbitrary triangle $\triangle P_{1} P_{2} P_{3}$ with 9 degrees of freedom [42]. We concern with one type of subdivision of a triangle shown in Figure 3.1. Let $O$ be the centroid of the triangle $T:=\triangle P_{1} P_{2} P_{3}$, and PS element used throughout this chapter
is defined to be a piecewise quadratic function on 6 small sub-triangles and $C^{1}$ continuous on $T$. $O$ is chosen to be the centroid because we wish the multilevel subspaces from PS element are nested. For the PS element on a non-degenerated triangle, it has 9 degrees of freedom. In other words, we have the following property for PS element.

Proposition 3.1. For a given piecewise quadratic function on the splitting sub-triangles as shown in Figure 3.1, if it is $C^{1}$ on the triangle $\triangle P_{1} P_{2} P_{3}$ with given values, as well as derivatives at $P_{1}, P_{2}$, and $P_{3}$, then this function is uniquely determined.


Figure 3.1: PS element is composed of 6 piecewise quadratic polynomials $p_{1}$, $p_{2}, \ldots, p_{6}$, respectively on a triangle

Existence of such PS element is a consequence of the affine invariance property of PS element, and this will be shown in Section 3.2.2. Uniqueness has been proved in [42]. The interpolation data we discussed here have 9 degrees of freedom, i.e., function values and the derivatives at 3 vertices. For derivatives at a vertex, there are 2 degrees of freedom if we notice that any two different directional derivatives of a $C^{1}$ continuous function at a point determine its tangent plane at that point from fundamental calculus.

We may construct 9 basis functions associated with these 9 degrees of freedom for PS elements on a triangle. Such basis functions are so called local
basis functions, compared to the basis functions defined on the whole mesh. We may drop local if the context is clear. For a given triangle $\triangle P_{1} P_{2} P_{3}$ in Figure 3.2, we may define 9 basis functions $\phi_{T, P, 0}, \phi_{T, P, 1}$ and $\phi_{T, P, 2}$ centered at the vertex $P \in\left\{P_{1}, P_{2}, P_{3}\right\}$. For $\phi_{T, P_{1}, 0}$, we have

$$
\begin{equation*}
\left.\phi_{T, P_{1}, 0}\right|_{P_{1}}=1,\left.\quad \frac{\partial \phi_{T, P_{1}, 0}}{\partial d_{1}\left(P_{1}\right)}\right|_{P_{1}}=0,\left.\quad \frac{\partial \phi_{T, P_{1}, 0}}{\partial d_{2}\left(P_{1}\right)}\right|_{P_{1}}=0 \tag{3.2}
\end{equation*}
$$

and

$$
\left.\phi_{T, P_{1}, 0}\right|_{P_{j}}=0,\left.\quad \frac{\partial \phi_{T, P_{1}, 0}}{\partial d_{1}\left(P_{j}\right)}\right|_{P_{j}}=0,\left.\quad \frac{\partial \phi_{T, P_{1}, 0}}{\partial d_{2}\left(P_{j}\right)}\right|_{P_{j}}=0, \quad j=2,3 .
$$



Figure 2a


Figure 2b

Figure 3.2: Notation for the triangle $T$
Recall that directional derivative of a function $f$ in $\vec{d}$ direction is $\frac{\partial f}{\partial \vec{d}}=$ $\cos (\theta) \frac{\partial f}{\partial \vec{x}}+\sin (\theta) \frac{\partial f}{\partial \vec{y}}$, where $\vec{d}$ is a direction with the unit length and the angle between $\vec{d}$ and $x$ axis is $\theta$. Here, by $d_{1}(P), d_{2}(P)$ we denote the directions of the unit length along the edges starting from the vertex $P$ following the right hand rule. For example, $d_{1}\left(P_{1}\right)$ is the direction in $\overrightarrow{P_{1} P_{2}}$, and $d_{2}\left(P_{1}\right)$ is in $\overrightarrow{P_{1} P_{3}}$. We drop $P$ from $d_{1}(P), d_{2}(P)$ if the context is clear. Let $\left|d_{j}(P)\right|, j=1,2$ be the length of the edge starting from the vertex $P$ and in the direction $d_{j}(P)$. For instance, $\left|d_{1}\left(P_{1}\right)\right|=\left|P_{1} P_{2}\right| . T$ may come into the notation to suggest that the relevant symbols be based on the triangle $T$, such as $d_{j}(T, P),\left|d_{j}(T, P)\right|$ in Section 3.2.3. Subscript $T$ is used to imply that $\phi_{T, P_{1}, 0}$ is defined on the
triangle $T$ locally. If the context is clear, the subscript $T$ in $\phi_{T, P_{1}, 0}$ is omitted. Clearly, $\phi_{P_{1}, 0}$ interpolates the value at the vertex $P_{1}$ while keeps all other interpolation data such as derivatives at $P_{1}, P_{2}, P_{3}$ and values at $P_{2}, P_{3}$ zeros. This is similar to the Hermite interpolants in one dimensional case [34]. $\phi_{P_{1}, 1}$ interpolates the derivative at $P_{1}$ in $d_{1}$ direction, i.e.,

$$
\begin{equation*}
\left.\frac{\partial \phi_{P_{1}, 1}}{\partial d_{1}}\right|_{P_{1}}=1 /\left|P_{1} P_{2}\right|,\left.\quad \phi_{P_{1}, 1}\right|_{P_{1}}=0,\left.\quad \frac{\partial \phi_{P_{1}, 1}}{\partial d_{2}}\right|_{P_{1}}=0 \tag{3.3}
\end{equation*}
$$

and

$$
\left.\phi_{P_{1}, 1}\right|_{P_{j}}=0,\left.\quad \frac{\partial \phi_{P_{1}, 1}}{\partial d_{1}}\right|_{P_{j}}=0,\left.\quad \frac{\partial \phi_{P_{1}, 1}}{\partial d_{2}}\right|_{P_{j}}=0 \quad j=2,3,
$$

where scale parameter $1 /\left|P_{1} P_{2}\right|$ is used for the affine invariance purpose. In particular, if $\left|P_{1} P_{2}\right|=1$, then $\phi_{P_{1}, 1}$ has the unit derivative in $\overrightarrow{P_{1} P_{2}}$ direction at $P_{1}$. Function $\phi_{P_{1}, 2}$ which interpolates the directional derivative at $P_{1}$ in the direction $\overrightarrow{P_{1} P_{3}}$ (i.e., $d_{2}\left(P_{1}\right)$ ) is defined similarly. Likewise, we may define the basis functions centered at $P_{2}, P_{3}$.

These 9 basis functions are PS elements with prescribed interpolation data. By Proposition 3.1, we claim such basis functions are uniquely determined. They are linear independent. If not, then we have a PS element which is a linear combination of them such that this PS element is identically zero on the triangle. If we notice that each basis function interpolates the different data of values and derivatives at 3 vertices, then this PS element must have non zero value or directional derivative in some direction at one vertex. This contradicts the fact that this PS element is identically zero on the triangle. Thus we have

Proposition 3.2. $\left\{\phi_{T, P, j}\right\}_{P=P_{1}, P_{2}, P_{3}, j=0,1,2}$ is a local basis for PS elements on the triangle $\triangle P_{1} P_{2} P_{3}$.

Now we consider when two PS elements on two neighboring triangles join $C^{1}$ continuously. First, we give a proposition on the property of PS element (see Figure 3.1).

Proposition 3.3. If quadratic functions $p 1, p 2$ are $C^{1}$ continuous across the common boundary $O q_{3}$ (see Figure 3.1), then

$$
\begin{equation*}
p 1=p 2+\lambda l^{2} \tag{3.4}
\end{equation*}
$$

where $\lambda$ is an arbitrary constant, and $l=0$ refers to the function of the line through $O$ and $q 3$.

Let's consider the directional derivative of $p 1, p 2$ in Figure 3.1 in $\overrightarrow{O q_{3}}$ direction. By (3.4), we have

$$
\frac{\partial p 1}{\partial \overrightarrow{O q_{3}}}=\frac{\partial p 2}{\partial \overrightarrow{O q_{3}}}, \quad\left(\frac{\partial l}{\partial \overrightarrow{O q_{3}}}=0\right)
$$

Then directional derivative of the PS element in $\overrightarrow{O q_{3}}$ direction on $P_{1} P_{2}$ is linear. Moreover, it interpolates two corresponding directional derivatives at $P_{1}$ $P_{2}$. However, the directional derivative in $\overrightarrow{P_{1} P_{2}}$ direction of the PS element are respectively two linear functions on $P_{1} q_{3}$ and $q_{3} P_{2}$, and they join $C^{1}$ continuously. The same analysis is applied to $T^{\prime}:=\triangle P_{1} P_{3} P_{4}$ in Figure 3.3. Let $q$ be the midpoint of $P_{1} P_{3}$, and $O, O^{\prime}$ be the centroids of $\triangle P_{1} P_{2} P_{3}, \triangle P_{1} P_{3} P_{4}$, respectively. Then we may anticipate that two PS elements on two neighboring triangles sharing the same values and derivatives at vertices $P_{1}$ and $P_{3}$ could fail to join $C^{1}$ continuously across the common boundary $P_{1} P_{3}$ unless $O, q, O^{\prime}$ are co-linear.


Figure 3.3: Two neighboring triangles sharing the common edge $P_{1} P_{3}$
However, on a three-direction mesh, $P_{1} P_{2} P_{3} P_{4}$ forms a parallelogram, and then $O, q, O^{\prime}$ fall on the line through $P_{2}$ and $P_{4}$ automatically. Hence, two PS
elements join $C^{1}$ continuously in this case. This is the reason why we restrict ourselves to three-direction meshes.

### 3.2.2 Affine map and the local basis on the standard equilateral triangle

In this section, we illustrate how to compute basis functions efficiently in a uniform way. Affine map plays an important role in the connection of a PS element on an arbitrary triangle with a PS element on the standard triangle. The following proposition is a basic property of the affine map.

Proposition 3.4. There is a unique affine map mapping from $T$ to $T_{1}$ or from $T_{1}$ to $T$, where $T, T_{1}$ are two non-degenerated triangles.

Proof. First we prove there is a unique affine map which maps three vertices of $T$ to the corresponding vertices of $T_{1}$. Let $P_{1}\left(x_{1}, y_{1}\right), P_{2}\left(x_{2}, y_{2}\right), P_{3}\left(x_{3}, y_{3}\right)$ be three vertices of $T$ and $P_{1}^{\prime}\left(x_{1}^{\prime}, y_{1}^{\prime}\right), P_{2}^{\prime}\left(x_{2}^{\prime}, y_{2}^{\prime}\right), P_{3}^{\prime}\left(x_{3}^{\prime}, y_{3}^{\prime}\right)$ be the vertices of $T_{1}$, respectively. The affine map $\mathcal{A}: T \rightarrow T_{1}$ satisfies

$$
\binom{x_{i}^{\prime}}{y_{i}^{\prime}}:=\left(\begin{array}{ll}
a_{11} & a_{12} \\
a_{21} & a_{22}
\end{array}\right)\binom{x_{i}}{y_{i}}+\binom{b_{1}}{b_{2}}, \quad i=1,2,3 .
$$

Rearranging the above equations yields

$$
\left(\begin{array}{c}
x_{1}^{\prime} \\
x_{2}^{\prime} \\
x_{3}^{\prime} \\
y_{1}^{\prime} \\
y_{2}^{\prime} \\
y_{3}^{\prime}
\end{array}\right)=\left(\begin{array}{cccccc}
x_{1} & y_{1} & 1 & 0 & 0 & 0 \\
x_{2} & y_{2} & 1 & 0 & 0 & 0 \\
x_{3} & y_{3} & 1 & 0 & 0 & 0 \\
0 & 0 & 0 & x_{1} & y_{1} & 1 \\
0 & 0 & 0 & x_{2} & y_{2} & 1 \\
0 & 0 & 0 & x_{3} & y_{3} & 1
\end{array}\right)\left(\begin{array}{c}
a_{11} \\
a_{12} \\
b_{1} \\
a_{21} \\
a_{22} \\
b_{2}
\end{array}\right)
$$

Obviously, unique existence of the affine map is equivalent to the non-singularity of the matrix

$$
\left(\begin{array}{lll}
x_{1} & y_{1} & 1 \\
x_{2} & y_{2} & 1 \\
x_{3} & y_{3} & 1
\end{array}\right)
$$

which implies that $P_{1}, P_{2}, P_{3}$ are non-collinear. This uniquely fixed $a_{i j}, b_{j}$, $i=1,2, j=1,2$ and gives the affine map. Moreover, the map is injective. If not, then there exist two different points $x_{1}, x_{2}$ in $T$ with the expression

$$
x_{1}=\alpha_{1} P_{1}+\alpha_{2} P_{2}+\alpha_{3} P_{3}, \quad x_{2}=\beta_{1} P_{1}+\beta_{2} P_{2}+\beta_{3} P_{3}
$$

such that $\mathcal{A} x_{1}=\mathcal{A} x_{2}$, that is

$$
\left(\alpha_{1}-\beta_{1}\right) P_{1}^{\prime}+\left(\alpha_{2}-\beta_{2}\right) P_{2}^{\prime}+\left(\alpha_{3}-\beta_{3}\right) P_{3}^{\prime}=0 .
$$

This shows that $\alpha_{j}=\beta_{j}, j=1,2,3$ and thus $x_{1}=x_{2}$. This is a contradiction. The map is surjective. For any point $x^{\prime}=\sum_{i=1}^{3} \alpha_{i} P_{i}^{\prime} \in T_{1}$, it is clear that $\mathcal{A} x=x^{\prime}$, with $x=\sum_{i=1}^{3} \alpha_{i} P_{i} \in T$.

It shows that this affine map is a one to one correspondence from $T$ to $T_{1}$. Existence and uniqueness of the affine map from $T_{1}$ to $T$ can be verified in an analogue manner. This completes the proof.

Proposition 3.5. Let the affine map $\mathcal{A}: T \rightarrow T_{1}$, with $\mathcal{A}\left(P_{1}\right)=P_{1}^{\prime}, \mathcal{A}\left(P_{2}\right)=$ $P_{2}^{\prime}$. If $x=P_{1}+t\left(P_{2}-P_{1}\right)$, then $x^{\prime}=\mathcal{A}(x)=P_{1}^{\prime}+t\left(P_{2}^{\prime}-P_{1}^{\prime}\right)$ is on the line $P_{1}^{\prime} P_{2}^{\prime}$.

Proof.

$$
\begin{aligned}
x^{\prime} & =\mathcal{A}(x)=\mathcal{A}\left(P_{1}\right)+t \mathcal{A}\left(P_{2}-P_{1}\right)=t \mathcal{A}\left(P_{2}\right)+(1-t) \mathcal{A}\left(P_{1}\right) \\
& =t P_{2}^{\prime}+(1-t) P_{1}^{\prime}=P_{1}^{\prime}+t\left(P_{2}^{\prime}-P_{1}^{\prime}\right)
\end{aligned}
$$

This completes the proof.
Proposition 3.5 tells that if a function $f\left(x^{\prime}\right)$ is defined on the line segment $P_{1}^{\prime} P_{2}^{\prime}$, and the affine map satisfies $\mathcal{A}\left(P_{1}\right)=P_{1}^{\prime}, \mathcal{A}\left(P_{2}\right)=P_{2}^{\prime}$, then $f(\mathcal{A}(x))$ is defined on the line segment $P_{1} P_{2}$. Based on Proposition 3.5, next proposition explores the relationship of the derivatives of the two functions $f\left(x^{\prime}\right)$ and $g(x)=f(\mathcal{A}(x))$.

Proposition 3.6. If the affine map $\mathcal{A}$ satisfies $\mathcal{A}\left(P_{1}\right)=P_{1}^{\prime}, \mathcal{A}\left(P_{2}\right)=P_{2}^{\prime}$, and $g(x)=f(\mathcal{A}(x))$, then

$$
\left.\left|P_{1} P_{2}\right| \frac{\partial g}{\partial\left(\overrightarrow{P_{1} P_{2}}\right)}\right|_{P_{1}}=\left.\left|P_{1}^{\prime} P_{2}^{\prime}\right| \frac{\partial f}{\partial\left(\overline{P_{1}^{\prime} P_{2}^{\prime}}\right)}\right|_{P_{1}^{\prime}}
$$

where $\left.\frac{\partial g}{\partial\left(\overrightarrow{\left.P_{1} P_{2}\right)}\right.}\right|_{P_{1}}=\nabla g \cdot \frac{\overrightarrow{P_{1} P_{2}}}{\left|\overrightarrow{P_{1} P_{2}}\right|}$ is the directional derivative of $g$ at point $P_{1}$.
By Proposition 3.6 we have

$$
\begin{equation*}
\phi_{T, P_{i}, j}(x)=\phi_{\tilde{T}, \tilde{P}_{i}, j}\left(\mathcal{A}_{T}(x)\right) /\left|d_{j}\left(T, P_{i}\right)\right|, \quad i=1,2,3, j=1,2 \tag{3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\phi_{T, P_{i}, 0}(x)=\phi_{\tilde{T}, \tilde{P}_{i}, 0}\left(\mathcal{A}_{T}(x)\right), \quad i=1,2,3 . \tag{3.6}
\end{equation*}
$$

Here, $\tilde{T}:=\triangle \tilde{P}_{1} \tilde{P}_{2} \tilde{P}_{3}$ is the standard equilateral triangle with unit edge lengths and affine $\operatorname{map} \mathcal{A}_{T}: T \rightarrow \tilde{T}$. By $\left|d_{j}\left(T, P_{i}\right)\right|$ we denote the length of the edge of the triangle $T$ starting from $P_{i}$ in $d_{j}$ direction.

Proposition 3.6 provide a general way to construct the PS element on an arbitrary triangle. More precisely, We construct the basis functions on the standard triangle $\tilde{T}$ first, and find the affine map mapping a given general triangle $T$ to $\tilde{T}$ next. Finally, $(3.5,3.6)$ determines the basis functions on $T$.

If $\mathcal{A}(T)=\tilde{T}, g(x)=f(\mathcal{A}(x)), x \in T, \tilde{x} \in \tilde{T}$, and $f(\tilde{x})$ is a PS element on $\tilde{T}$, then $g(x)$ is also a PS element on $T$. Furthermore, $\left\{\phi_{\tilde{T}, \tilde{P}, j}(\mathcal{A}(x))\right\}_{\tilde{P}=\tilde{P}_{1}, \tilde{P}_{2}, \tilde{P}_{3}, j=0,1,2}$ is a basis of PS elements on $T .\left\{\phi_{\tilde{T}, \tilde{P}, j}(\tilde{x})\right\}_{\tilde{P}=\tilde{P}_{1}, \tilde{P}_{2}, \tilde{P}_{3}, j=0,1,2}$ are given explicitly on $\tilde{T}$ in the end of this section, and this, on the other hand, verifies the existence of the PS element in Proposition 3.1.

The representation of a PS element on $T$ and its correspondence after being affinely mapped onto $\tilde{T}$ have the following affine invariance property.

Proposition 3.7. Let the PS element $f$ be defined on $T$ with the expression

$$
\begin{equation*}
f(x)=\sum_{P_{i} \in \mathcal{N}(T), j=0,1,2} \alpha_{T, P_{i}, j} \phi_{P, j}(x) \tag{3.7}
\end{equation*}
$$

where $\mathcal{N}(T):=\left\{P_{1}, P_{2}, P_{3}\right\}$ is the set of three vertices of triangle $T$. By the affine map, we map the function $f$ on the triangle $T$ to its correspondence $\tilde{f}$ on the standard triangle $\tilde{T}$ by

$$
\begin{equation*}
\tilde{f}(\tilde{x}):=f(\tilde{A} \tilde{x}+\tilde{b}) \tag{3.8}
\end{equation*}
$$

where $P_{j}=\tilde{A} \tilde{P}_{j}+\tilde{b}, j=1,2,3$ and $\left\{\tilde{P}_{j}, j=1,2,3\right\}$ are three vertices of $\tilde{T} . \tilde{f}$ is also a PS element on $\tilde{T}$, and it has the expression

$$
\begin{equation*}
\tilde{f}(\tilde{x})=\sum_{i=1,2,3, j=0,1,2} \alpha_{\tilde{T}, \tilde{P}_{i}, j} \phi_{\tilde{T}, \tilde{P}_{i}, j}(\tilde{x}) . \tag{3.9}
\end{equation*}
$$

Then, we have

$$
\begin{equation*}
\alpha_{T, P_{i}, j}=\alpha_{\tilde{T}, \tilde{P}_{i}, j}, \quad i=1,2,3, j=0,1,2 \tag{3.10}
\end{equation*}
$$

Proof. Taking directional derivative of (3.7), we have

$$
\left.\frac{\partial f}{\partial d_{j}}\right|_{P_{i}}=\left.\alpha_{T, P_{i}, j} \frac{\partial \phi_{T, P_{i}, j}}{\partial d_{j}}\right|_{P_{i}},
$$

which implies

$$
\alpha_{T, P_{i}, j}=\left.\frac{\partial f}{\partial d_{j}\left(T, P_{i}\right)}\right|_{P_{i}}\left|d_{j}\left(T, P_{i}\right)\right|, \quad j=1,2
$$

Likewise, we have

$$
\tilde{\alpha}_{\tilde{T}, \tilde{P}_{i}, j}=\left.\frac{\partial \tilde{f}}{\partial d_{j}\left(\tilde{T}, \tilde{P}_{i}\right)}\right|_{\tilde{P}_{i}}\left|d_{j}\left(\tilde{T}, \tilde{P}_{i}\right)\right|, \quad j=1,2 .
$$

Using the property of affine map in Proposition 3.6 and (3.8), we obtain

$$
\left.\frac{\partial \tilde{f}}{\partial d_{j}\left(\tilde{T}, \tilde{P}_{i}\right)}\right|_{\tilde{P}_{i}}\left|d_{j}\left(\tilde{T}, \tilde{P}_{i}\right)\right|=\left.\frac{\partial f}{\partial d_{j}\left(T, P_{i}\right)}\right|_{P_{i}}\left|d_{j}\left(T, P_{i}\right)\right|
$$

This implies (3.10) and thus completes the proof.
Equation (3.10) motivates us to introduce the scale parameter $1 /\left|P_{1} P_{2}\right|$ in the definition of PS local basis in (3.3). The representations of the PS element therfore are invariant on different triangles. This enables us to carry
out computing on the standard reference triangle instead of a general triangle. Furthermore, (3.10) provide a quick way to map a PS element on a general triangle to a PS element on $\tilde{T}$.

To end this subsection, we shall construct the basis functions on $\tilde{T}$ with explicit expressions due to their importance. As shown in Figure 3.2b, $\tilde{T}$ is divided into 6 sub-triangles. We set up the coordinate system as follows. Let $\tilde{P}_{1}$ be the origin, and $\overrightarrow{\tilde{P}_{1} \tilde{P}_{2}}$ be in the positive $x$-axis direction. On each subtriangle, a basis function is a quadratic function. Let $p 1, \ldots, p 6$ be polynomials on sub-triangles $\triangle O \tilde{P}_{1} q 3, \triangle O q 3 \tilde{P} 2, \ldots, \triangle O q 2 \tilde{P}_{1}$, respectively.

Using Proposition 3.3, together with the prescribed values and derivatives at three vertices of $\tilde{T}$, we compute the basis functions on $\tilde{T}$. For the basis function $\phi_{P_{1}, 1}$, we have

$$
\begin{align*}
& p 1=p 5=p 6=x-\frac{1}{\sqrt{3}} y-\frac{3}{2} x^{2}+\frac{1}{\sqrt{3}} x y+\frac{1}{6} y^{2}  \tag{3.11}\\
& p 2=p 3=p 4=p 1+2(x-1 / 2)^{2} .
\end{align*}
$$

$\phi_{P_{1}, 2}$ shall be symmetric to $\phi_{P_{1}, 1}$ about the line $y=\left(\tan 30^{\circ}\right) x$, and $\phi_{P_{1}, 0}$ has the expression
(3.12) $p 1=1-2 x^{2}-2 y^{2}$,

$$
\begin{aligned}
p 2 & =1-2 x^{2}-2 y^{2}+4(x-1 / 2)^{2} \\
p 3 & =p 4=1-2 x^{2}-2 y^{2}+4(x-1 / 2)^{2}+3\left(y+\frac{1}{\sqrt{3}} x-\frac{1}{\sqrt{3}}\right)^{2} \\
p 5 & =1-2 x^{2}-2 y^{2}+3\left(y+\frac{1}{\sqrt{3}} x-\frac{1}{\sqrt{3}}\right)^{2}, \\
p 6 & =p 1 .
\end{aligned}
$$

Other six basis functions on the vertices $\tilde{P}_{2}, \tilde{P}_{3}$ can be described in an analogue manner, or be obtained by using the symmetric property of the basis functions.

As a conclusion of this sub-section, we derive the basic properties of the affine map, and compute the basis functions on the standard triangle. Combining the affine map with the basis functions on $\tilde{T}$, we get the basis functions on a general triangle. Moreover, we also give an explicit method to transform a PS element on a general triangle to its correspondence on $\tilde{T}$ by (3.10). In
general, we provide the necessary basic knowledge of the local basis functions, which are the fundamental bricks to build the basis functions on three-direction meshes.

### 3.2.3 Basis functions on the meshes

Once we have the basis functions on an arbitrary triangle $\triangle P_{1} P_{2} P_{3}$, we shall define the basis functions on three-direction meshes. Let $\mathcal{T}$ be the triangulation, and $\mathcal{N}$ be the vertex set. There are 3 basis functions denoted by $\left\{\phi_{P, j}\right\}_{j=0, x, y}$ associated with each vertex $P$. Here, in order to define the basis function uniformly, let $\phi_{P, x}, \phi_{P, y}$ be the functions interpolating the derivatives at the node $P$ in the directions of the positive $x$-axis, $y$-axis, respectively. The supports of $\phi_{P, x}, \phi_{P, y}$ consist of the neighboring triangles of the vertex $P$. On each neighboring triangle $T, \phi_{P, x}, \phi_{P, y}$ can be written as the linear combination of the local basis functions $\phi_{T, P, 1}, \phi_{T, P, 2}$, if we notice that the tangent plane at a vertex of a $C^{1}$ function is uniquely determined by any two different directional derivatives. Let the set $S$ be $S:=\left\{\phi_{P, j}, P \in \mathcal{N}, j=0, x, y\right\}$. Then the space $V$ is the linear span of the set $S . V$ reproduces the polynomial of degree 2 on each triangle in the given mesh.

We assume that the refinement procedure is to subdivide each triangle by connecting the middle point of each edge. After each refinement, one triangle is subdivided into 4 similar sub-triangles. Furthermore, we assume that for a triangle $T \in \mathcal{T}$, three interior angles are bounded below and above by some constants. This assumption implies that the ratio $l_{i} / l_{j}$ is bounded up and below for $i, j=1,2,3$ by some constants, where $l_{1}, l_{2}, l_{3}$ are the lengthes of a given triangle $T$. The mesh size thus can be measured by any one of $l_{1}, l_{2}$, and $l_{3}$. If the mesh size is $h$, which is defined by $h:=\max _{T \in \mathcal{T}}\{$ diameter of $T\}$, then $\left\|\phi_{P, j}\right\|_{L_{2}} \simeq h$.

On a triangle, say $\triangle P_{1} P_{2} P_{3}$, we have two representations of a function in $V$, i.e., in terms of local basis functions $\left\{\phi_{T, P_{i}, j}\right\}_{i=1,2,3, j=0,1,2}$ or $\left\{\phi_{P_{i}, j}\right\}_{i=1,2,3, j=0, x, y}$. In the following, we wish to discuss the basis $\left\{\phi_{P, j}\right\}_{P \in \mathcal{N}, j=0, x, y}$ in detail. $\phi_{P, j}$
is defined on the six neighboring triangles (denoted by $\mathcal{T}(P)$ ) of the vertex $P$ in $\mathcal{T}$. On its supported triangle, $\phi_{P, j}$ is a PS element, and

$$
\left.\phi_{P, 0}\right|_{P}=1,\left.\phi_{P, 0}\right|_{q}=0, \mathcal{N} \ni q \neq P,\left.\frac{\partial \phi_{P, 0}}{\partial j}\right|_{q}=0, q \in \mathcal{N}, j=x, y
$$

and

$$
\left.\frac{\partial \phi_{P, x}}{\partial x}\right|_{P}=1 / h,\left.\frac{\partial \phi_{P, x}}{\partial y}\right|_{q}=0,\left.\phi_{P, x}\right|_{q}=0, q \in \mathcal{N},\left.\frac{\partial \phi_{P, x}}{\partial x}\right|_{q}=0, \mathcal{N} \ni q \neq P
$$

$\phi_{P, y}$ is defined exactly in the same way. For any $T \in \mathcal{T}\left(P_{1}\right)$, say $T$ in Figure 3.1, $\phi_{P_{1}, x}$ ( or $\phi_{P_{1}, y}$ ) on $T$ can be represented in terms of local basis $\phi_{T, P_{1}, 1}$ and $\phi_{T, P_{1}, 2}$. If the angle between $P_{1} P_{2}$ and $x$ axis is $\theta_{1}$ and the angle between $P_{1} P_{3}$ and $x$ axis is $\theta_{2}$, then

$$
\left.\frac{\partial \phi_{P_{1}, x}}{\partial d_{1}\left(T, P_{1}\right)}\right|_{P_{1}}=\left.\cos \left(\theta_{1}\right) \frac{\partial \phi_{P_{1}, x}}{\partial x}\right|_{P_{1}},\left.\quad \frac{\partial \phi_{P_{1}, x}}{\partial d_{2}\left(T, P_{1}\right)}\right|_{P_{1}}=\left.\cos \left(\theta_{2}\right) \frac{\partial \phi_{P_{1}, x}}{\partial x}\right|_{P_{1}}
$$

Recall that $d_{i}\left(T, P_{1}\right), i=1,2$, are two directions along two edges of $T$ starting from the vertex $P_{1} . \phi_{P_{1}, x}$ restricted on $T$ can be written as

$$
\begin{equation*}
\left.\phi_{P_{1}, x}\right|_{T}=\cos \left(\theta_{1}\right) \frac{\left|d_{1}\left(T, P_{1}\right)\right|}{h} \phi_{T, P_{1}, 1}+\cos \left(\theta_{2}\right) \frac{\left|d_{2}\left(T, P_{1}\right)\right|}{h} \phi_{T, P_{1}, 2} . \tag{3.13}
\end{equation*}
$$

$\phi_{P_{1}, 0}$ is simply the same as $\phi_{T, P_{1}, 0}$ on T, i.e.,

$$
\begin{equation*}
\left.\phi_{P_{1}, 0}\right|_{T}=\phi_{T, P_{1}, 0} \tag{3.14}
\end{equation*}
$$

(3.13-3.14) gives the local representations of the (global) basis functions. It is clear that $\phi_{P, j}, j=0, x, y$ can be locally represented by the local basis function $\phi_{T, P, j} T \in \mathcal{T}(P), P \in \mathcal{N}(T), j=0,1,2$. In other words, basis function $\phi_{P, j}, j=0, x, y$ are PS elements on each neighboring triangle $T$ of $P$. Now we propose the definition of discrete $L_{2}$ inner product and its associated norm. For any function $f \in V$, it is a PS element on a triangle $T \in \mathcal{T} . f$ can be represented in terms of the (global) basis functions or in terms of local basis
functions triangle by triangle. We use the latter representation. Suppose $f$ is represented as (3.7) on a triangle $T \in \mathcal{T}$. By (3.8), we obtain

$$
\begin{equation*}
\|f\|_{L_{2}(T)}^{2}=\|\tilde{f}\|_{L_{2}(\tilde{T})}^{2} \operatorname{vol}(T) \tag{3.15}
\end{equation*}
$$

Note that

$$
\left\|\phi_{\tilde{T}, \tilde{P}, j}\right\|_{L_{2}}^{2} \simeq 1, \quad \tilde{P} \in \mathcal{N}(\tilde{T}), j=0,1,2
$$

Then

$$
\|\tilde{f}\|_{L_{2}(\tilde{T})}^{2} \simeq \sum_{i=1,2,3, j=0,1,2} \tilde{\alpha}_{\tilde{T}, \tilde{P}_{i}, j}^{2}
$$

Combining with (3.15), we have

$$
\|f\|_{L_{2}(T)}^{2} \simeq \sum_{i=1,2,3, j=0,1,2} \alpha_{T, P_{i}, j}^{2} \operatorname{vol}(T)
$$

where $\operatorname{vol}(T) \simeq h^{2}$ if the mesh size is $h$.
For two given functions $f(x), g(x) \in V$ which have the representations like (3.7) with the coefficients $\left\{\alpha_{T, P_{i}, j}\right\}$ and $\left\{\beta_{T, P_{i}, j}\right\}$, respectively, their discrete $L_{2}$ inner product can be defined as follows,

$$
\begin{equation*}
\langle f, g\rangle_{D_{2}}:=\sum_{T \in \mathcal{T}} \operatorname{vol}(T) \sum_{P \in \mathcal{N}(T), j=0,1,2} \alpha_{T, P, j} \beta_{T, P, j} \tag{3.16}
\end{equation*}
$$

Obviously, the discrete $L_{2}$ norm induced by the discrete $L_{2}$ inner product is equivalent to the $L_{2}$ norm for any function in $V$, i.e.,

## Proposition 3.8.

$$
\langle f, f\rangle_{D_{2}} \simeq\|f\|_{L_{2}}^{2}, \quad f \in V
$$

By (3.10), we may replace $\alpha_{P_{i}, j}, \beta_{P_{i}, j}$ with $\widetilde{\alpha}_{\tilde{P}_{i}, j}, \widetilde{\beta}_{\tilde{P}_{i}, j}$ in the definition of discrete inner product (3.16). It suggests that we use affine map to map the function on $T$ to the function on the standard triangle $\tilde{T}$ first with the representation (3.9) on $\tilde{T}$. The discrete $L_{2}$ inner product or norm then can be computed on $\tilde{T}$, respectively. Let $T$ run over all triangles in $\mathcal{T}$ and sum the discrete inner products or norms up, then we get the corresponding discrete
$L_{2}$ inner product or norm. Next, we introduce the new notation for multilevel analysis purpose. Recall that the uniform subdivision procedure is utilized to refine the mesh. We have 3 basis functions sitting on each vertex in the mesh for each level. We add the superscript or the subscript $k$ to denote the triangulations, set of nodes, discrete inner products, etc, in the corresponding level $k$. For instance, for the $k$-th level mesh, we may use the following notation: $\mathcal{T}_{k}, \mathcal{N}_{k},\langle\cdot, \cdot\rangle_{D_{2}(k)}, \phi_{P, j}^{k}, \alpha_{P, j}^{k}, S_{k}, V_{k}$ etc.

Since (global) basis functions are PS elements on each triangle and $C^{1}$ continuous, we may claim the refinability of the basis functions.

## Proposition 3.9. Basis functions $\left\{\phi_{P, j}^{k}\right\}_{P \in \mathcal{N}, j=0, x, y}$ are refinable.

Proof. We prove that $\phi_{P_{1}, x}^{0}$ can be represented as the linear combination of the higher level basis functions. Refinability of other basis functions shall be proved similarly. We take the triangle $T$ in Figure 3.1. Other neighboring triangles of $P_{1}$ on which $\phi_{P_{1}, x}^{0}$ is supported can be dealt with in the same way.

It is easily seen that $\phi_{P_{1}, x}^{0}$ determines the values and derivatives in $x, y$ directions at $P_{1}, q_{2}, q_{3}$. Let these data be ( $\alpha_{P_{1}, 0}=0, \alpha_{P_{1}, x}=1, \alpha_{P_{1}, y}=0$ ), $\left(\alpha_{q_{2}, 0}, \alpha_{q_{2}, x}, \alpha_{q_{2}, y}\right)$, and ( $\alpha_{q_{3}, 0}, \alpha_{q_{3}, x}, \alpha_{q_{3}, y}$ ). Then we prove that

$$
f=\sum_{P \in\left\{P_{1}, q_{2}, q_{3}\right\}} \alpha_{P, 0} \phi_{P, 0}^{1}+\frac{1}{2} \sum_{P \in\left\{P_{1}, q_{2}, q_{3}\right\}, j=x, y} \alpha_{P, j} \phi_{P, j}^{1}
$$

is identical to $\phi_{P_{1}, x}^{0}$ on $\triangle P_{1} q 2 q 3$.
On $\triangle P_{1} q_{2} q_{3}, f$ shares the same values and derivatives at $P_{1}, q_{2}, q_{3}$ with $\phi_{P_{1}, x}^{0}$. Furthermore, $f$ and $\phi_{P_{1}, x}^{0}$ are both PS elements on $\triangle P_{1} q_{2} q_{3}$ in $V_{1}$. By Proposition 3.1, they are identical. $f \equiv \phi_{P_{1}, x}^{0}$ can be proved on the rest subtriangles of $T$ similarly. This completes the proof.

Refinability of the basis functions implies that the subspaces $\left\{V_{k}\right\}_{k=0}^{\infty}$ are nested, i.e.,

$$
V_{0} \subset V_{1} \cdots
$$

Here, $S_{k}:=\left\{\phi_{P, j}^{k}\right\}_{P \in \mathcal{N}_{k} ; j=0, x, y}$ and $V_{k}:=\operatorname{span} S_{k}$.

The wavelet subspace $\widetilde{W}_{k}$ is defined to be the orthogonal complement of $V_{k-1}$ in $V_{k}$ with respect to the discrete inner product $\langle\cdot, \cdot\rangle_{D_{2}(k)}$ for $V_{k}$. More precisely, for any $w \in \widetilde{W}_{k}$, we have

$$
\langle w, v\rangle_{D_{2}(k)}=0, \quad \forall v \in V_{k-1} .
$$

$\widetilde{W}_{k}$ is the counterpart of $W_{k}$, which denotes the $L_{2}$ orthogonal complement of $V_{k-1}$ in $V_{k}$. Here, we define the $L_{2}$ orthogonal projection to be $Q_{k}: L_{2} \rightarrow V_{k}$ with $\left\langle Q_{k} u, v\right\rangle_{L_{2}}=\langle u, v\rangle_{L_{2}}, u \in L_{2}, v \in V_{k}$. Then we have

$$
W_{k}:=P_{k}\left(\bigcup_{l=0}^{\infty} V_{l}\right)
$$

where $P_{k}:=Q_{k}-Q_{k-1}, k=0,1, \ldots,\left(Q_{-1}:=0\right)$.
For $\left\{V_{k}\right\}_{k=0}^{\infty}$, it is easily seen that the following Jackson inequality

$$
\begin{equation*}
\left\|v_{l}-Q_{k} v_{l}\right\|_{L_{2}} \lesssim 2^{-2(l-k)}\left|v_{l}\right|_{H^{2}}, v_{k} \in V_{k}, 0 \leq k<l \tag{3.17}
\end{equation*}
$$

and the Bernstein inequality

$$
\begin{equation*}
\left\|v_{k}\right\|_{H^{2}} \lesssim 4^{k}\left\|v_{k}\right\|_{L_{2}}, v_{k} \in V_{k} \tag{3.18}
\end{equation*}
$$

hold true.
In the next section, the basis functions for $\widetilde{W}_{k}$ shall be constructed, and we call these functions wavelets.

### 3.3 Construction of wavelets

One favorite feature for the wavelets is the simple structure, which means the local (or small) supports of the wavelets. For each refinement, we get more newly created vertices. These vertices are middle points of edges in the old mesh. There are wavelet functions associated with these new vertices in the finer level mesh. We shall construct the wavelets for $\widetilde{W}_{1}$. Wavelet functions for other wavelet subspaces can be constructed similarly.

Let two triangles $\triangle P_{1} P_{2} P_{3}$ and $\triangle P_{1} P_{3} P_{4}$ in the mesh $\mathcal{T}_{0}$ share the common edge $P_{1} P_{3}$ (see Figure 3.3). After the refinement, new vertex $q$, the middle point of the edge $P_{1} P_{3}$, is created. We define 3 wavelets associated with $q$ in the following form,

$$
\begin{equation*}
\psi_{q, j^{\prime}}^{1}:=\beta_{q, j^{\prime}}^{1} \phi_{q, j^{\prime}}^{1}(x)+\sum_{i=1,3, j=0, x, y} \beta_{P_{i}, j}^{1} \phi_{P_{i}, j}^{1}(x), \quad j^{\prime}=0, x, y . \tag{3.19}
\end{equation*}
$$

$\psi_{q, j^{\prime}}^{1}$ is automatically orthogonal to all basis functions of $V_{0}$, except $\left\{\phi_{P, j}^{0}, P=\right.$ $\left.P_{1}, P_{3}, j=0, x, y\right\}$ with respect to the discrete $L_{2}$ inner product $\langle\cdot, \cdot\rangle_{D_{2}(1)}$. Since we have 7 coefficients in the expression of the wavelet with only 6 constraints (which are the orthogonality conditions between the wavelet and 6 basis functions sitting at $P_{1}, P_{3}$ ), there is at least one nontrivial solution for the coefficients in (3.19). We want further that $\beta_{q, j^{\prime}}^{1} \neq 0$, and thus we may normalize the coefficients to make $\beta_{q, j^{\prime}}^{1}=1$. Suppose that there is only one nontrivial solution with $\beta_{q, j^{\prime}}^{1}=0$. Then, we prove that the rest coefficients must be zeros, too. For a fixed $j^{\prime}$, suppose we have

$$
\begin{equation*}
\psi_{q, j^{\prime}}^{1}=\sum_{i=1,3, j=0, x, y} \beta_{P_{i}, j}^{1} \phi_{P_{i}, j}^{1}(x) . \tag{3.20}
\end{equation*}
$$

Considering the support of $\phi_{P_{i}, j}^{1}(x), i=1,3, j=0, x, y$, we find that $\sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ shall be orthogonal to $\phi_{P_{1}, j}^{0}(x), j=0, x, y$, because $\phi_{P_{3}, j}^{1}(x), j=0, x, y$ are orthogonal to $\phi_{P_{1}, j}^{0}(x), j=0, x, y$ automatically. Here, orthogonality is in the sense of discrete $L_{2}$ inner product.

First, we shall note that $\sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ has zero values and zero derivatives at vertices other than $P_{1}$ in $\mathcal{N}_{1} . \sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ orthogonal to $\phi_{P_{1}, j}^{0}(x), j=0, x, y$ implies that $\sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ must have zero value, as well as zero derivatives at vertex $P_{1}$. Thus $\sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ has zero values as well as zero derivatives at all vertices in $\mathcal{N}_{1}$. Since on each triangle in $\mathcal{T}_{1}$, $\sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ is a PS element, $\sum_{j=0, x, y} \beta_{P_{1}, j}^{1} \phi_{P_{1}, j}^{1}(x)$ is then identically zero on all triangles in $\mathcal{T}_{1}$ by Proposition 3.1. Hence $\beta_{P_{1}, j}^{1}=0, j=0, x, y$. Likewise, we obtain $\beta_{P_{3}, j}^{1}=0, j=0, x, y$.

This contradicts the fact that there is at least one non-trivial solution for
coefficients and thus excludes the solution with $\beta_{q, j^{\prime}}^{1}=0$. Consequently, we may write the wavelet in the form,

$$
\begin{equation*}
\psi_{q, j^{\prime}}^{1}:=\phi_{q, j^{\prime}}^{1}+\sum_{i=1,3, j=0, x, y} \beta_{P_{i}, j}^{1} \phi_{P_{i}, j}^{1}(x), \quad j^{\prime}=0, x, y . \tag{3.21}
\end{equation*}
$$

Clearly, our short supported wavelet has at most 7 non-zero coefficients. For convenience of stating the next lemma, we write $\psi_{q, j^{\prime}}^{1}$ in a clearer expression:

$$
\begin{equation*}
\psi_{q, j^{\prime}}^{1}:=\phi_{q, j^{\prime}}^{1}+\sum_{P \in \mathcal{N}_{q}, j=0, x, y} \beta_{q, j^{\prime}, P, j}^{1} \phi_{P, j}^{1}(x), \quad j^{\prime}=0, x, y . \tag{3.22}
\end{equation*}
$$

Here, $\mathcal{N}_{q}$ is defined to be the set of two vertices in $\mathcal{N}_{0}$, whose midpoint is $q$. Accordingly, $\mathcal{N}_{q}^{k}, k=0,1,2, \ldots$ is defined for the meshes in different levels.

Note that for two different wavelets $\psi_{q, j}^{1}$ and $\psi_{q^{\prime}, j^{\prime}}^{1}$ with $\{q, j\} \neq\left\{q^{\prime}, j^{\prime}\right\}$, they have the different components $\phi_{q, j}^{1}$ and $\phi_{q^{\prime}, j^{\prime}}^{1}$. This leads to the following lemma.

Lemma 3.1. $\left\{\psi_{q, j}^{1}\right\}_{q \in \mathcal{N}_{1} \backslash \mathcal{N}_{0}, j=0, x, y}$ is a stable basis for $\widetilde{W}_{1}$ in $L_{2}$ space.
Proof. Let

$$
g=\sum_{q \in \mathcal{N}_{1} \backslash \mathcal{N}_{0}} \sum_{j=0, x, y} \gamma_{q, j} \psi_{q, j}^{1} .
$$

Recalling that

$$
\psi_{q, j}^{1}=\phi_{q, j}^{1}+\sum_{P \in \mathcal{N}_{q}} \sum_{j_{1}=0, x, y} \beta_{q, j, P, j_{1}}^{1} \phi_{P, j_{1}}^{1},
$$

we have

$$
\begin{align*}
g & =\sum_{q \in \mathcal{N}_{1} \backslash \mathcal{N}_{0}, j=0, x, y}\left[\gamma_{q, j} \phi_{q, j}^{1}+\gamma_{q, j} \sum_{P \in \mathcal{N}_{q}, j_{1}=0, x, y} \beta_{q, j, P, j_{1}}^{1} \phi_{P, j_{1}}^{1}\right] \\
& =\sum_{q \in \mathcal{N}_{1} \backslash \mathcal{N}_{0}, j=0, x, y} \gamma_{q, j} \phi_{q, j}^{1}+\sum_{P \in \mathcal{N}_{0}, j=0, x, y} \alpha_{P, j} \phi_{P, j}^{1}, \tag{3.23}
\end{align*}
$$

where

$$
\alpha_{P, j}=\sum_{P \in \mathcal{N}_{q}, j_{1}=0, x, y} \gamma_{q, j_{1}} \beta_{q, j_{1}, P, j}^{1}
$$

Since there are six neighboring triangles around the vertex $P \in \mathcal{N}_{0}$ in a three-direction mesh, we have

$$
\begin{equation*}
\alpha_{P, j}^{2} \leq C \sum_{P \in \mathcal{N}_{q}, j_{1}=0, x, y} \gamma_{q, j_{1}}^{2} \tag{3.24}
\end{equation*}
$$

where $C$ is a fixed constant only dependent on the number of the neighboring triangles and coefficients $\left\{\beta_{q, j_{1}, P, j}^{1}\right\}$.

Note that $\left\{\phi_{P, j}^{1}\right\}_{P \in \mathcal{N}_{1}, j=0, x, y}$ is a Riesz basis of $V_{1}$ in $L_{2}$ space. By (3.23), we have

$$
\|g\|_{L_{2}}^{2} \simeq \sum_{q \in \mathcal{N}_{1} \backslash \mathcal{N}_{0}, j=0, x, y} \gamma_{q, j}^{2}\left\|\phi_{q, j}^{1}\right\|_{L_{2}}^{2}+\sum_{P \in \mathcal{N}_{0}, j=0, x, y} \alpha_{P, j}^{2}\left\|\phi_{P, j}^{1}\right\|_{L_{2}}^{2}
$$

Assume the mesh size of $\mathcal{T}_{1}$ is $h$. Then

$$
\left\|\phi_{P, j}^{1}\right\|_{L_{2}}^{2} \simeq h^{2}, \quad P \in \mathcal{N}_{1}, j=0, x, y
$$

By the estimate of (3.24), we have

$$
\|g\|_{L_{2}}^{2} h^{-2} \simeq \sum_{q \in \mathcal{N}_{1} \backslash \mathcal{N}_{0}, j=0, x, y} \gamma_{q, j}^{2}
$$

This completes the proof.
As a direct consequence of Lemma 3.1, we have

## Corollary 3.1.

$$
\left\|\sum_{P \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}, j=0, x, y} \alpha_{P, j}^{k} \psi_{P, j}^{k}\right\|_{L_{2}}^{2} \simeq\left(2^{-k}\right)^{2} \sum_{P \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}, j=0, x, y}\left(\alpha_{P, j}^{k}\right)^{2}, k=0,1, \ldots,
$$

provided that the initial mesh size is $\simeq 1$.
Counting the number of basis functions, we claim that $\left\{\psi_{q, j}^{k}\right\}_{q \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}}$ is a $L_{2}$ stable basis for $\widetilde{W}_{k}$. To keep the consistence of the notation, we let $\psi_{q, j}^{0}=\phi_{q, j}^{0}$ and $\mathcal{N}_{-1}$ be the empty set in what follows. Then we have $\widetilde{W}_{0}=V_{0}$.

Let $Y_{k}$ be the orthogonal projection from $V_{k+1}$ to $V_{k}$ with respect to $\langle\cdot, \cdot\rangle_{D_{2}(k+1)}$, i.e., $\left\langle Y_{k} v_{k+1}, v^{\prime}\right\rangle_{D_{2}(k+1)}=\left\langle v_{k+1}, v^{\prime}\right\rangle_{D_{2}(k+1)}, v_{k} \in V_{k+1}, v^{\prime} \in V_{k}$. The following proposition exhibits the stability of the wavelet basis for $\widetilde{W}_{k}$ in the Sobolev space $H^{2}$.

## Proposition 3.10.

$$
4^{k}| | u_{k} \|_{L_{2}} \simeq\left|u_{k}\right|_{H^{2}}, \quad u_{k} \in \widetilde{W}_{k}, \quad k=1,2, \cdots
$$

Proof. We have

$$
\begin{aligned}
\left\|u_{k}\right\|_{L_{2}} & \simeq\left\|u_{k}\right\|_{D_{2}(k)}=\left\|\left(I-Y_{k-1}\right) u_{k}\right\|_{D_{2}(k)} \\
& \leq\left\|\left(I-Q_{k-1}\right) u_{k}\right\|\left\|_{D_{2}(k)} \simeq\right\|\left(I-Q_{k-1}\right) u_{k} \|_{L_{2}} \lesssim 2^{-2 k}\left|u_{k}\right|_{H^{2}}
\end{aligned}
$$

In the last inequality, we use the approximation property (the Jackson inequality (3.17)) of the subspace $V_{k}$. Hence,

$$
\left|\left|u_{k} \|_{L_{2}} \lesssim 4^{-k}\right| u_{k}\right|_{H^{2}}
$$

By the inverse inequality, we obtain

$$
\left|u_{k}\right|_{H^{2}} \lesssim 4^{k}| | u_{k} \|_{L_{2}}
$$

Combining the previous two inequalities together completes the proof.

### 3.4 Stability of the wavelets in Sobolev space $H^{2}$

In this section, we consider the stability of the wavelets in Sobolev space $H^{2}$. We go through several necessary lemmas before we finally reach the result. The first lemma is a type of strengthened Cauchy-Schwarz inequality, which is of particular importance in ensuring that the wavelets form a Bessel sequence in $H^{2}$. Similar lemmas appeared in [2, 47]. In [47], a particular short proof is presented on the basis of the advanced knowledge of spaces interpolation. For self-contain purpose, we shall give here a quite basic proof based on some cancellation properties employed in [2].

Lemma 3.2. (Strengthened Cauchy-Schwarz inequality)

$$
\begin{equation*}
\left\langle\triangle u_{k}, \triangle v_{l}\right\rangle_{L_{2}} \leq C 2^{-(l-k)}\left(2^{2 k}\left\|u_{k}\right\|_{L_{2}}\right)\left(2^{2 l}\left\|v_{l}\right\|_{L_{2}}\right), \quad u_{k} \in V_{k}, v_{l} \in V_{l}, k \leq l \tag{3.25}
\end{equation*}
$$

where $C$ is a constant independent of $k, l$. Here, $\triangle$ is the Laplacian operator defined by $\Delta u:=u_{x x}+u_{y y}$.

Proof. Without loss of generality, we take $k=0$. We consider the case $l>$ $k+2$, and (3.25) holds true for the case $l \leq k+2$ automatically by CauchySchwarz inequality and inverse inequality.

We concern with $L_{2}$ inner product of $\triangle u_{0}$ and $\triangle v_{l}$ on a given triangle $T \in \mathcal{T}_{0}$, as shown in Figure 3.1. By the rule of the refinement procedure, the sub-triangles in the mesh with level number $l$ are all similar to $T$. Moreover, $v_{l}$ can be written in terms of basis functions in $V_{l}$ sitting on the vertices in $\mathcal{N}_{l} \cap \bar{T}$.

Note that $\triangle u_{0}$ are piecewise constants on 6 sub-triangles of $T$, and $\triangle v_{l}$ are piecewise constants on all sub-triangles of $T^{\prime} \in \mathcal{T}_{l}$. We focus on the basis functions in level $l$ with their support in one of 6 sub-triangles of $T$ in level 0 , say $\triangle O P_{3} q_{2}$. The shaded area in Figure 3.4 is defined in such a way that every basis function $v_{l} \in V_{l}$ centered in the shaded area has the support totally within $\triangle O P_{3} q_{2}$. We choose the largest possible triangle in $\triangle O P_{3} q_{2}$ as the shaded area. For any basis function $v_{l} \in V_{l}$ sitting on the vertex within the shaded area in Figure 3.4, because its first derivatives vanish on the boundary of $\triangle O P_{3} q_{2}$, it is orthogonal to $\triangle u_{0}$ because
$\int_{\triangle O P_{3} q_{2}}\left(\triangle u_{0}\right)\left(\nabla \cdot \nabla v_{l}\right)=\left(\triangle u_{0}\right) \int_{\triangle O P_{3} q_{2}}\left(\nabla \cdot \nabla v_{l}\right)=\left(\triangle u_{0}\right) \int_{\partial \triangle O P_{3} q_{2}} \nabla v_{l} \cdot \vec{n} d s=0$, where $\partial \triangle O P_{3} q_{2}$ is the boundary of the triangle $\triangle O P_{3} q_{2}$.

Roughly speaking, the distance between the shaded area and the boundary $\partial \triangle O P_{3} q_{2}$ is about the same size of the support of basis function in $V_{l}$, i.e., $2^{-l}$ multiple of the length of the edge of $\triangle O P_{3} q_{2} .2^{-l}$ comes from the uniform refinement procedure. This leads to the estimate

$$
\begin{equation*}
\frac{\operatorname{area}\left(\triangle O P_{3} q_{2}\right)-\text { shaded area }}{\operatorname{area}\left(\triangle O P_{3} q_{2}\right)} \leq C 2^{-l}, \quad(k=0) \tag{3.26}
\end{equation*}
$$

Note that the basis functions centered on the boundary of $\triangle O P_{3} q_{2}$ has the support outside the shaded area and $\triangle u_{0}$ is a constant on $\triangle O P_{3} q_{2}$. Then, for


Figure 3.4: Supports of two PS elements in two different levels
a function $v_{l} \in V_{l}$, we have

$$
\begin{aligned}
\int_{\Delta O P_{3} q_{2}}\left(\Delta u_{0}\right)\left(\triangle v_{l}\right) & =\int_{\Delta O P_{3} q_{2} \backslash \text { ShadedArea }}\left(\triangle u_{0}\right)\left(\Delta v_{l}\right) \\
& \leq\left\|\Delta u_{0}\right\|_{L_{2}\left(\Delta O P_{3} q_{2} \backslash \text { ShadedArea) }\right)}\left\|\Delta v_{l}\right\|_{L_{2}\left(\Delta O P_{3} q_{2} \backslash \text { ShadedArea }\right)} \\
& \leq\left\|\Delta u_{0}\right\|_{L_{2}\left(\Delta O P_{3} q_{2} \backslash \text { ShadedArea }\right)}\left\|\Delta v_{l}\right\|_{L_{2}\left(\Delta O P_{3} q_{2}\right)} \\
& \leq C 2^{-l}\left\|\Delta u_{k}\right\|_{L_{2}\left(\Delta O P_{3} q_{2}\right)}\left\|\Delta v_{l}\right\|_{L_{2}\left(\Delta O P_{3} q_{2}\right)}
\end{aligned}
$$

where $2^{-l}$ comes from (3.26).
Inequality (3.25) on the sub-triangle $\triangle O P_{3} q_{2}$ can be obtained by the use of inverse inequality. After we sum the estimate on each sub-triangle up, the result follows.

Lemma 3.2, together with Corollary 3.1 shows that the wavelet system $\left\{2^{-k} \psi_{q, j}^{k}\right\}_{k, q, j}$ is a Bessel sequence in $H^{2}$.

Theorem 3.1. Let $u=\sum_{k=0}^{J} u_{k}$ with $u_{k}=\sum_{q \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}} \sum_{j=0, x, y} \beta_{q, j}^{k}\left(2^{-k} \psi_{q, j}^{k}\right) \in$ $\widetilde{W}_{k}, k=0,1, \ldots, J$. Then

$$
\begin{equation*}
\|u\|_{H^{2}}^{2} \lesssim \sum_{k=0}^{J} \sum_{q \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}} \sum_{j=0, x, y}\left(\beta_{q, j}^{k}\right)^{2} . \tag{3.27}
\end{equation*}
$$

Proof. It follows that

$$
\begin{align*}
|u|_{H^{2}}^{2} & \simeq \sum_{l, m=0}^{J}\left\langle\triangle u_{l}, \Delta u_{k}\right\rangle_{L_{2}}  \tag{3.28}\\
& \lesssim \sum_{l, m=0}^{J} 2^{-|l-m|}\left(4^{l}\left\|u_{l}\right\|_{L_{2}}\right)\left(4^{m}\left\|u_{m}\right\|_{L_{2}}\right)(\text { Lemma } \\
& \lesssim \sum_{l=0}^{J} 4^{2 l}\left\|u_{l}\right\|_{L_{2}}^{2}
\end{align*}
$$

Moreover,

$$
\begin{aligned}
\|u\|_{L_{2}}^{2} & =\left\|\sum_{l=0}^{J} u_{l}\right\|^{2} \leq \sum_{l, m=0}^{J} 4^{-(l+m)}\left(4^{l}\left\|u_{l}\right\|_{L_{2}}^{2}\right)\left(4^{m}\left\|u_{k}\right\|_{L_{2}}^{2}\right) \\
& \lesssim \sum_{l=0}^{J} 4^{2 l}\left\|u_{l}\right\|_{L_{2}}^{2} .
\end{aligned}
$$

Combining above two inequalities, we have

$$
\|u\|_{H^{2}}^{2} \lesssim \sum_{k=0}^{J} 4^{2 k}\left\|u_{k}\right\|_{L_{2}}^{2} \lesssim \sum_{k=0}^{J} 4^{2 k}\left(2^{-2 k} \sum_{q \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}} \sum_{j=0, x, y}\left(2^{-k} \beta_{q, j}^{k}\right)^{2}\right)
$$

where we used Corollary 3.1 in the last inequality.
This completes the proof.

Bessel sequence property of the wavelet basis guarantees the upper bound for $\|u\|_{H^{2}}$ in (3.27). If we can prove the lower bound for $\|u\|_{H^{2}}$, then the wavelet basis form a Riesz basis in $H^{2}$. This is the task of the remaining section.

Let the projection operator $Z_{k}^{l}$ be defined by,

$$
Z_{k}^{l}:=\left(1-Y_{k-1}\right) Y_{k} \cdots Y_{l-1}, \quad k=1,2, \ldots, l-1 ; \quad Z_{0}^{l}:=Y_{0} Y_{1} \cdots Y_{l-1} .
$$

$Z_{k}^{l}$ projects a function in $V_{l}$ into $\widetilde{W}_{k}(k<l)$ and the estimate of this operator is the key in verifying the stability of the wavelet system in Sobolev space.

## Lemma 3.3.

$$
\left\|Z_{k}^{l}\right\|_{L_{2} \rightarrow L_{2}} \leq C 2^{\lambda(l-k)}, \quad k<l
$$

where $\lambda<2$, and $C$ is a constant independent of $l, k$.
Suppose Lemma 3.3 holds true for a moment, then we claim
Theorem 3.2. $\left\{2^{-k} \psi_{q, j}^{k}\right\}_{k, q, j}$ is a Riesz basis for $H^{2}$.
Proof. Let $u=\sum_{k=0}^{J} u_{k}$ with $u_{k}=\sum_{q \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}} \sum_{j=0, x, y} \beta_{q, j}^{k}\left(2^{-k} \psi_{q, j}^{k}\right) \in \widetilde{W}_{k}$, $k=0, \ldots, J$. First, we verify that

$$
\|u\|_{H^{2}}^{2} \gtrsim \sum_{k=0}^{J} 4^{2 k}\left\|u_{k}\right\|_{L_{2}}^{2} .
$$

For simplicity, by $Z_{k}$ we denote $Z_{k}^{J}$ in the proof. Note that $u_{k}=Z_{k} u, k=$ $0, \ldots, J$. Let $s=2$. Then we have

$$
\begin{aligned}
\mathrm{RHS} & \equiv \sum_{k=0}^{J} 4^{k s}\left\|Z_{k} u\right\|_{L_{2}}^{2} \\
& =\sum_{k=0}^{J} 4^{k s}\left\langle Z_{k} u, Z_{k} u\right\rangle_{L_{2}}=\sum_{k=0}^{J} 4^{k s}\left\langle Z_{k} \sum_{l=0}^{J} P_{l} u, Z_{k} \sum_{m=0}^{J} P_{m} u\right\rangle_{L_{2}} \\
& =\sum_{l, m=0}^{J} \sum_{k=0}^{\min (l, m)} 4^{k s}\left\langle Z_{k} P_{l} u, Z_{k} P_{m} u\right\rangle_{L_{2}} \\
& \leq \sum_{l, m=0}^{J} \sum_{k=0}^{\min (l, m)} 4^{k s}\left\|Z_{k} P_{l} u\right\|_{L_{2}}\left\|Z_{k} P_{m} u\right\|_{L_{2}} \quad\left(\left\|Z_{k}^{l}\right\|_{L_{2} \rightarrow L_{2}} \lesssim 2^{\lambda(l-k)}, k<l\right) \\
& \lesssim \sum_{l, m=0}^{J} \sum_{k=0}^{\min (l, m)} 4^{k s}\left\|P_{l} u\right\|_{L_{2}} 2^{\lambda(l-k)}\left\|P_{m} u\right\|_{L_{2}} 2^{\lambda(m-k)} \\
& =\sum_{l, m=0}^{J} \sum_{k=0}^{\min (l, m)} 2^{-(l+m-2 k)(s-\lambda)}\left(2^{l s}\left\|P_{l} u\right\|_{L_{2}}\right)\left(2^{m s}\left\|P_{m} u\right\|_{L_{2}}\right) \\
& \lesssim \sum_{l, m=0}^{J} 2^{-(l+m-2 \min (l, m))(s-\lambda)}\left(2^{l s}\left\|P_{l} u\right\|_{L_{2}}\right)\left(2^{m s}\left\|P_{m} u\right\|_{L_{2}}\right) \\
& \lesssim \sum_{l=0}^{J}\left(2^{l s}\left\|P_{l} u\right\|_{L_{2}}\right)^{2} .
\end{aligned}
$$

Therefore, we obtain

$$
\mathrm{RHS} \lesssim\|u\|_{H^{2}}^{2}
$$

By Corollary 3.1,

$$
\|u\|_{H^{2}}^{2} \gtrsim \mathrm{RHS} \simeq \sum_{k=0}^{J} \sum_{q \in \mathcal{N}_{k} \backslash \mathcal{N}_{k-1}} \sum_{j=0, x, y}\left(\beta_{q, j}^{k}\right)^{2}
$$

Note that $\left\{2^{-k} \psi_{q, j}^{k}\right\}_{k, q, j}$ is a Bessel sequence in $H^{2}$ by Theorem 3.1 and $\cup_{k=0}^{\infty} V_{k}$ is dense in $H^{2}$. It follows that $\left\{2^{-k} \psi_{q, j}^{k}\right\}_{k, q, j}$ is a Riesz basis in $H^{2}$. This completes the proof.

It remains to prove Lemma 3.3.

Proof.

$$
\begin{aligned}
\left\|Z_{k}^{l}\right\|_{L_{2} \leftarrow L_{2}} & :=\sup _{u_{l} \in V_{l}} \frac{\left\|Z_{k}^{l} u_{l}\right\|_{L_{2}}}{\left\|u_{l}\right\|_{L_{2}}} \\
& =\sup _{u_{l} \in V_{l}} \frac{\left\|Z_{k}^{l} u_{l}\right\|_{L_{2}}}{\left\|v_{k}\right\|_{D_{2}(k)}} \frac{\left\|v_{k}\right\|_{D_{2}(k)}}{\left\|v_{k+1}\right\|_{D_{2}(k+1)}} \cdots \frac{\left\|v_{l-1}\right\|_{D_{2}(l-1)}}{\left\|v_{l}\right\|_{D_{2}(l)}} \frac{\left\|v_{l}\right\|_{D_{2}(l)}}{\left\|v_{l}\right\|_{L_{2}}},
\end{aligned}
$$

where $v_{j}:=Y_{j}^{l} u_{l}, j<l, v_{l}:=u_{l}$ and $Y_{j}^{l}:=Y_{j} \cdots Y_{l-1}$.
By Proposition 3.8, we have for $k=1,2, \ldots, l-1$,

$$
\frac{\left\|Z_{k}^{l} u_{l}\right\|_{L_{2}}}{\left\|v_{k}\right\|_{D_{2}(k)}} \simeq \frac{\left\|\left(1-Y_{k-1}\right) v_{k}\right\|_{D_{2}(k)}}{\left\|v_{k}\right\|_{D_{2}(k)}} \lesssim 1, \quad \frac{\left\|v_{l}\right\|_{D_{2}(l)}}{\left\|v_{l}\right\|_{L_{2}}} \simeq 1
$$

For $k=0, Z_{0}^{l} u_{l}=v_{0}$, and thus

$$
\frac{\left\|Z_{0}^{l} u_{l}\right\|_{L_{2}}}{\left\|v_{0}\right\|_{D_{2}(0)}} \simeq 1
$$

Now we consider $\frac{\left\|v_{j}\right\|_{D_{2}(j)}}{\left\|v_{j+1}\right\| D_{2}(j+1)}, k \leq j<l$. Note that $v_{j}=Y_{j} v_{j+1}$. Then

$$
\frac{\left\|v_{j}\right\|_{D_{2}(j)}}{\left\|v_{j+1}\right\|_{D_{2}(j+1)}}=\frac{\left\|v_{j}\right\|_{D_{2}(j)}}{\left\|v_{j}\right\|_{D_{2}(j+1)}} \frac{\left\|Y_{j} v_{j+1}\right\|_{D_{2}(j+1)}}{\left\|v_{j+1}\right\|_{D_{2}(j+1)}} \leq \frac{\left\|v_{j}\right\|_{D_{2}(j)}}{\left\|v_{j}\right\|_{D_{2}(j+1)}} .
$$

$\frac{\left\|Y_{j} v_{j+1}\right\|_{D_{2}(j+1)}}{\left\|v_{j+1}\right\|_{D_{2}(j+1)}} \leq 1$ because $Y_{j}$ is the orthogonal projection from $V_{j+1}$ to $V_{j}$ with respect to the discrete $L_{2}$ inner product $\langle\cdot, \cdot\rangle_{D_{2}(j+1)}$.

If we assume $\frac{\left\|v_{j}\right\|_{D_{2}(j)}}{\left\|v_{j}\right\| D_{D_{2}(j+1)}} \leq 2^{\lambda}(\lambda<2), k \leq j<l$ for a moment and postpone its lengthy and tedious proof to the next lemma, then we have $\frac{\left\|v_{j}\right\|_{D_{2}(j)}}{\left\|v_{j+1}\right\|_{D_{2}(j+1)}} \leq$ $2^{\lambda}$. Hence,

$$
\left\|Z_{k}^{l}\right\|_{L_{2} \rightarrow L_{2}} \leq C 2^{\lambda(l-k)}
$$

This completes the proof.

As pointed out in the previous lemma, the key step to verify $\left\|Z_{k}^{l}\right\|_{L_{2} \rightarrow L_{2}} \leq$ $C 2^{\lambda(l-k)}$ is that the ratio $\|f\|_{D_{2}(j)} /\|f\|_{D_{2}(j+1)} \leq 2^{\lambda}$ for any $f \in V_{j}$.

Lemma 3.4. There exists a positive number $\lambda<2$ that

$$
\|f\|_{D_{2}(k)}^{2} /\|f\|_{D_{2}(k+1)}^{2} \leq 2^{2 \lambda}, \quad f \in V_{k} .
$$

Proof. Without loss of generality, we take $k=0$ and verify that for some $\lambda<2$,

$$
\left\|\left.f\right|_{T}\right\|_{D_{2}(0)}^{2} /\left\|\left.f\right|_{T}\right\|_{D_{2}(1)}^{2} \leq 2^{2 \lambda}
$$

is true for any $f \in V_{0}$ on any triangle $T$ in $\mathcal{T}_{0}$.
For a function $f(x) \in V_{0}$ on $T$ as shown in Figure 3.1, let

$$
\begin{equation*}
f(x)=\sum_{i=1,2,3, j=0,1,2} \alpha_{P_{i}, j} \phi_{P_{i}, j}(x), \quad x \in T \tag{3.29}
\end{equation*}
$$

By the unique affine $\operatorname{map} \mathcal{A}: \tilde{T} \rightarrow T$, we map $f(x)$ to $\tilde{f}(\tilde{x})$ by letting $\tilde{f}(\tilde{x}):=f(\mathcal{A}(\tilde{x}))$, where $\tilde{T}$ is a standard equilateral triangle. By (3.10) $\left(\alpha_{P_{i}, j}=\right.$ $\left.\widetilde{\alpha}_{\tilde{P}_{i}, j}\right)$, we have

$$
\begin{equation*}
\tilde{f}(\tilde{x})=\sum_{i=1,2,3, j=0,1,2} \alpha_{P_{i}, j} \phi_{\tilde{P}_{i}, j}(\tilde{x}), \quad \tilde{x} \in \widetilde{T} \tag{3.30}
\end{equation*}
$$

This indicates that we may compute the discrete $L_{2}$ norm of a given function after it is affine mapped onto the standard triangle. Therefore, without loss of generality, we assume that $f$ is a function defined on the standard equilateral triangle $T$ with unit edge lengths in the following proof.

Let $f(x)$ be the PS element on $T$ in the form of (3.29). To compute $\left\|\left.f\right|_{T}\right\|_{D_{2}(0)}$, we simply have

$$
\left\|\left.f\right|_{T}\right\|_{D_{2}(0)}^{2}=\operatorname{vol}(T) \sum_{i=1,2,3, j=0,1,2} \alpha_{P_{i}, j}^{2}
$$

It is clear that $\left||f|_{T} \|_{D_{2}(0)}^{2}\right.$ is the sum of squares of 9 variables $\left\{\alpha_{P_{i}, j}\right\}_{i=1,2,3, j=0,1,2}$. On the other hand, these variables uniquely determine the PS element on $T$. To compute $\left\|\left.f\right|_{T}\right\|_{D_{2}(1)}$ of the discrete $L_{2}$ norm in higher level, we need values and derivatives of $f$ at the midpoints of three edges of $T$. In other words, we shall write $f(x)$ in terms of higher level PS elements with local basis on each sub-triangle $T_{1}:=\triangle P_{1} q 3 q 2, T_{2}:=\triangle q 2 q 3 q 1, T_{3}:=\triangle q 3 P_{2} q 1, T_{4}:=\triangle q 2 q 1 P_{3}$, respectively. Since $T$ is symmetric, the data (values and derivatives) of two basis functions $\left\{\phi_{P_{1}, 0}, \phi_{P_{1}, 1}\right\}$ at $q_{1}, q_{2}$, and $q_{3}$ are sufficient for the computation. Data of other 7 basis functions at $q_{1}, q_{2}$, and $q_{3}$ can be obtained similarly due to the symmetry of the basis functions.

We concern with the values and derivatives of $f(x)$ at the midpoint $q 3$. Note that $\phi_{P_{3}, j}, j=0,1,2$ has zeros values and derivatives at any point on the line $P_{1} P_{2}$. Thus, $\alpha_{P_{3}, j}, j=0,1,2$ have no contribution to the value or derivatives at $q 3$. Only $\left\{\alpha_{P, j}, P=P_{1}, P_{2}, j=0,1,2\right\}$ have contribution to $q 3$. Such contribution can be obtained by studying the values and derivatives of $\left\{\phi_{P_{1}, 0}, \phi_{P_{1}, 1}\right\}$ at $q 3, q 2$. First, we derive the value and derivatives of $\phi_{P_{1}, 0}$ at $q 3$. $\phi_{P_{1}, 0}$ is a piecewise quadratic polynomials on $T$, and its expression is given explicitly in (3.12). By simple computing, we have

$$
\begin{align*}
\phi_{P_{1}, 0}(q 3) & =1 / 2,  \tag{3.31}\\
\frac{\partial \phi_{P_{1}, 0}}{\partial d_{2}\left(T_{1}, q 3\right)} & =2 \\
\frac{\partial \phi_{P_{1}, 0}}{\partial d_{1}\left(T_{1}, q 3\right)} & =1 \\
\frac{\partial \phi_{P_{1}, 0}}{\partial d_{2}\left(T_{3}, q 3\right)} & =-1 \\
\frac{\partial \phi_{P_{1}, 0}}{\partial d_{1}\left(T_{3}, q 3\right)} & =-2,
\end{align*}
$$

where $d_{j}(T, P)$ type definition of direction is given in Section 3.2.1. Since $\phi_{P_{1}, 0}$ is symmetric about the line $P_{1} q 1$, its value and derivatives at $q 2$ can be obtained by (3.31) (see Figure 3.5). Consequently, contribution of $\alpha_{P_{1}, 0}$ to $q 3, q 2$ can be computed from (3.31). Because of the symmetry, contribution of $\alpha_{P_{2}, 0}$ to $q 3$ can be computed from (3.31), too.


Figure 3.5: Value and derivatives of $\phi_{P_{1}, 0}$ on $T$
Second, we study the contribution of $\phi_{P_{1}, 1}$ at $q 3, q 2$. By the expression of $\phi_{P_{1}, 1}$ as in (3.11), we have (see Figure 3.6)

$$
\begin{align*}
\phi_{P_{1}, 1}(q 3) & =1 / 8  \tag{3.32}\\
\frac{\partial \phi_{P_{1}, 1}}{\partial d_{2}\left(T_{1}, q 3\right)} & =1 / 2 \\
\frac{\partial \phi_{P_{1}, 1}}{\partial d_{1}\left(T_{1}, q 3\right)} & =0 \\
\frac{\partial \phi_{P_{1}, 0}}{\partial d_{2}\left(T_{3}, q 3\right)} & =-1 / 2 \\
\frac{\partial \phi_{P_{1}, 1}}{\partial d_{1}\left(T_{3}, q 3\right)} & =-1 / 2 .
\end{align*}
$$

With (3.31-3.32) in hand, we compute the value and the derivatives of $f(x)$ at $q 3$ as follows,


Figure 3.6: Value and derivatives of $\phi_{P_{1}, 1}$ on $T$

$$
\begin{aligned}
f(q 3) & =\sum_{j=0,1,2, P=P_{1}, P_{2}} \alpha_{P, j} \phi_{P, j}(q 3) \\
\frac{\partial f}{\partial d}(q 3) & =\sum_{j=0,1,2, P=P_{1}, P_{2}} \alpha_{P, j} \frac{\partial \phi_{P, j}(q 3)}{\partial d}(q 3),
\end{aligned}
$$

where $d$ refers to one of four directions $\left\{d_{2}\left(T_{1}, q 3\right), d_{1}\left(T_{1}, q 3\right), d_{2}\left(T_{3}, q 3\right)\right.$, $\left.d_{1}\left(T_{3}, q 3\right)\right\}$ at $q 3$ shown in Figure 3.5 or Figure 3.6. Hence, the value and the derivatives of $f(x)$ at $q 3$ is a linear combination of $\left\{\alpha_{P_{1}, 0}, \alpha_{P_{1}, 1}, \alpha_{P_{1}, 2}, \alpha_{P_{2}, 0}, \alpha_{P_{2}, 1} \alpha_{P_{2}, 2}\right\}$. We give the details in the following,

$$
\begin{aligned}
f(q 3) & =\frac{1}{2}\left(\alpha_{P_{1}, 0}+\alpha_{P_{2}, 0}\right)+\frac{1}{8}\left(\alpha_{P_{1}, 1}+\alpha_{P_{2}, 2}\right), \\
\frac{\partial f}{\partial d_{2}\left(T_{1}, q 3\right)}(q 3) & =2\left(\alpha_{P_{1}, 0}-\alpha_{P_{2}, 0}\right)+\frac{1}{2}\left(\alpha_{P_{1}, 1}-\alpha_{P_{2}, 2}\right), \\
\frac{\partial f}{\partial d_{1}\left(T_{1}, q 3\right)}(q 3) & =\left(\alpha_{P_{1}, 0}-\alpha_{P_{2}, 0}\right)-\frac{1}{2} \alpha_{P_{2}, 2}+\frac{1}{2}\left(\alpha_{P_{1}, 2}+\alpha_{P_{2}, 1}\right), \\
\frac{\partial f}{\partial d_{2}\left(T_{3}, q 3\right)}(q 3) & =\left(\alpha_{P_{2}, 0}-\alpha_{P_{1}, 0}\right)-\frac{1}{2} \alpha_{P_{1}, 1}+\frac{1}{2}\left(\alpha_{P_{1}, 2}+\alpha_{P_{2}, 1}\right), \\
\frac{\partial f}{\partial d_{1}\left(T_{3}, q 3\right)}(q 3) & =2\left(\alpha_{P_{2}, 0}-\alpha_{P_{1}, 0}\right)+\frac{1}{2}\left(\alpha_{P_{2}, 2}-\alpha_{P_{1}, 1}\right) .
\end{aligned}
$$

Using the symmetric property of $T$, we get other values and derivatives at $q 2, q 1$.

To compute $\left\|\left.f\right|_{T_{1}}\right\|_{D_{2}(1)}$, we have

$$
\alpha_{T_{1}, q, 0}^{1}=f(q), \quad q=P_{1}, q 2, q 3
$$

and

$$
\alpha_{T_{1}, q, j}^{1}=\frac{1}{2} \frac{\partial f}{\partial d_{j}\left(T_{1}, q\right)}(q), \quad q=P_{1}, q 2, q 3, j=1,2
$$

where $\frac{1}{2}$ comes from the definition of $\frac{\partial \phi_{T_{1}, P, j}^{1}}{\partial d_{j}\left(T_{1}, P\right)}=2, \quad j=1,2$. It follows that

$$
\left\|\left.f\right|_{T_{1}}\right\|_{D_{2}(1)}^{2}=\frac{\operatorname{vol}(T)}{4} \sum_{P=P_{1}, q 2, q 3} \sum_{j=0,1,2}\left(\alpha_{T_{1}, P, j}^{1}\right)^{2}
$$

It is easily seen that $\|\left.\left. f\right|_{T_{1}}\right|_{D_{2}(1)} ^{2}$ is a quadratic form of the variables $\left\{\alpha_{P, j}\right\}_{P=P_{1}, P_{2}, P_{3}, j=0,1,2}$. Let $\alpha$ be the column vector form of these variables. Then in a similar way, we may compute $\left\|\left.f\right|_{T_{j}}\right\|_{D_{2}(1)}^{2}, j=2,3,4$, which are all quadratic forms of $\alpha$. Hence, we may write

$$
\left||f|_{T}\left\|_{D_{2}(1)}^{2}=\sum_{i=1}^{4}\right\| f\right|_{T_{i}} \|_{D_{2}(1)}^{2}=\operatorname{vol}(T) \alpha^{T} D \alpha
$$

where $D$ is a $9 \times 9$ symmetric positive definite matrix, and $\alpha^{T}$ is the transpose of the vector $\alpha$. It follows that

$$
1 / \lambda_{D, \max } \leq \frac{\left\|\left.f\right|_{T}\right\|_{D_{2}(0)}^{2}}{\left\|\left.f\right|_{T}\right\|_{D_{2}(1)}^{2}}=\frac{\alpha^{T} \alpha}{\alpha^{T} D \alpha} \leq 1 / \lambda_{D, \min }
$$

where $\lambda_{D, \text { min }}, \lambda_{D, \text { max }}$ are respectively the minimal, maximal eigenvalues of the matrix $D$. From the computation, we have

$$
\lambda_{D, \min }>0.0679>4^{(-2)}
$$

Let $\lambda=\log _{4}(1 / 0.0679)$. Then we have

$$
\frac{\left\|\left.f\right|_{T}\right\|_{D_{2}(0)}^{2}}{\left\|\left.f\right|_{T}\right\|_{D_{2}(1)}^{2}}<2^{2 \lambda}
$$

It is clear that on an arbitrary triangle $T$, we have the same estimate. After we sum up the above inequality over all triangles in $\mathcal{T}_{0}$, we have

$$
\frac{\|f\|_{D_{2}(0)}^{2}}{\|f\|_{D_{2}(1)}^{2}}<2^{2 \lambda}
$$

This completes the proof.

### 3.5 Examples

In this section, we shall demonstrate the computation of basis functions, as well as wavelet functions, on the regular triangular mesh. The mesh is constructed by connecting lower left vertex with upper right vertex of each square in the uniform tensor product mesh. The support of each basis function contains 6 triangles, as shown in Figure 3.7. On each triangle, the basis function is a PS element which is a piecewise continuous quadratic function on its six pieces of sub-triangles. Therefore, we have to compute each PS element one by one, and put them together to obtain the basis functions.


Figure 3.7: Supports of the basis functions $\phi_{O, j}^{0}, j=0, x, y$

For the basis function $\phi_{0,0}^{0}$, it satisfies the following conditions,

$$
\begin{aligned}
\phi_{O, 0}^{0}(O) & =1, \quad \phi_{O, 0}^{0}(P)=0, P \in \mathcal{N} e(O) \backslash\{O\} \\
\frac{\partial \phi_{O, 0}^{0}}{\partial x}(P) & =0, \quad \frac{\partial \phi_{O, 0}^{0}}{\partial y}(P)=0, P \in \mathcal{N} e(O)
\end{aligned}
$$

where $\mathcal{N} e(O):=\left\{O, P_{1}, P_{2}, \ldots, P_{6}\right\}$ is the set of vertices neighboring $O$ plus $O$ itself in Figure 3.7. $\phi_{O, x}^{0}, \phi_{O, y}^{0}$ satisfy the similar conditions at the vertex $O$ in $\mathcal{N} e(O)$.

For one triangle in the support of $\phi_{O, 0}^{0}$, say $\triangle O P_{1} P_{2}$, the PS element of the basis function $\phi_{O, 0}^{0}$ satisfies the above conditions at three vertices $O, P_{1}, P_{2}$, too. By the properties of PS element, unique PS element is determined.

There are two ways to carry out the computation. The first method is based on the affine transformation. Because we have the basis functions on a standard equilateral triangle, we may compute the basis functions on an arbitrary triangle by finding the affine map between these two triangles. The interpolation data on the vertices, such as directional derivatives, shall be computed accordingly. The second method uses the properties of PS element and computes the quadratic functions on each sub-triangles directly with the given conditions on three vertices.

In the following, we give the basis functions for $\left\{\phi_{O, j}^{0}\right\}_{j=0, x, y}$ (see Figure 3.7).

$$
\begin{aligned}
& \phi_{\mathbf{0}, \mathbf{0}}^{0} \\
& P_{11}= P_{21}=P_{32}=P_{42}=P_{53}=P_{63}=1-2 x^{2}+2 x y-2 y^{2} ; \\
& P_{12}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(y-2 x+1)^{2} ; \\
& P_{13}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(y-2 x+1)^{2}+(x+y-1)^{2} ; \\
& P_{14}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(x+y-1)^{2} ; \\
& P_{22}= P_{14} ; \\
& P_{23}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(x+y-1)^{2}+4(y-x / 2-1 / 2)^{2} ; \\
& P_{24}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+4(y-x / 2-1 / 2)^{2} ; \\
& P_{31}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(y-2 x-1)^{2} ; \\
& P_{33}=P_{24} ; \\
& P_{34}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(y-2 x-1)^{2}+4(y-x / 2-1 / 2)^{2} ; \\
& P_{41}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(x+y+1)^{2} ; \\
& P_{43}=P_{31} ; \\
& P_{44}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(x+y+1)^{2}+(y-2 x-1)^{2} ; \\
& P_{51}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+(x+y+1)^{2}+4(y-x / 2+1 / 2)^{2} ; \\
& P_{52}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+4(y-x / 2+1 / 2)^{2} ; \\
& P_{54}=P_{41} \\
& P_{61}=\left(1-2 x^{2}+2 x y-2 y^{2}\right)+4(y-x / 2+1 / 2)^{2}+(y-2 x+1)^{2} ; \\
& P_{62}=P_{12} ; \\
& P_{64}=P_{52} ;
\end{aligned}
$$

$$
\begin{aligned}
& \phi_{\mathbf{O}, \mathbf{x}}^{0} \\
& P_{11}=-\frac{3}{2} x^{2}+x y-y^{2}+x ; \\
& P_{12}=P_{11}+\frac{1}{2}(y-2 x+1)^{2} ; \\
& P_{13}=P_{11}+\frac{1}{2}(y-2 x+1)^{2}+\frac{1}{2}(x+y-1)^{2} ; \\
& P_{14}=P_{11}+\frac{1}{2}(x+y-1)^{2} ; \\
& P_{21}=P_{11}+(x-y)^{2} ; \\
& P_{22}=P_{21}+\frac{1}{2}(x+y-1)^{2} ; \\
& P_{23}=P_{22} ; \\
& P_{24}=P_{21} ; \\
& P_{31}=P_{21}+2 x^{2}-\frac{1}{2}(y-2 x-1)^{2} ; \\
& P_{32}=P_{21}+2 x^{2} ; \\
& P_{33}=P_{32} ; \\
& P_{34}=P_{31} ; \\
& P_{41}=\left(-\frac{3}{2} x^{2}-x y+y^{2}+x\right)-\frac{1}{2}(x+y+1)^{2} ; \\
& P_{42}=\left(-\frac{3}{2} x^{2}-x y+y^{2}+x\right) ; \\
& P_{43}=\left(-\frac{3}{2} x^{2}-x y+y^{2}+x\right)-\frac{1}{2}(y-2 x-1)^{2} ; \\
& P_{44}=\left(-\frac{3}{2} x^{2}-x y+y^{2}+x\right)-\frac{1}{2}(y-2 x-1)^{2}-\frac{1}{2}(x+y+1)^{2} ; \\
& P_{51}=\left(\frac{1}{2} x^{2}+x y+x\right)-\frac{1}{2}(x+y+1)^{2} ; \\
& P_{52}=\left(\frac{1}{2} x^{2}+x y+x\right) ; \\
& P_{53}=P_{52} ; \\
& P_{54}=P_{51} ; \\
& P_{61}=\left(-\frac{3}{2} x^{2}+x y+x\right)+\frac{1}{2}(y-2 x+1)^{2} ; \\
& P_{62}=P_{61} ; \\
& P_{63}=\left(-\frac{3}{2} x^{2}+x y+x\right) ; \\
& P_{64}=P_{63} ; \\
& \\
&
\end{aligned}
$$

and

$$
\begin{aligned}
& \phi_{O, y}^{0} \\
& P_{11}=-x y-\frac{1}{2} y^{2}+y ; \\
& P_{12}=P_{11} ; \\
& P_{13}=\left(-x y-\frac{1}{2} y^{2}+y\right)+\frac{1}{2}(x+y-1)^{2} ; \\
& P_{14}=P_{13} ; \\
& P_{21}=-x^{2}+x y-\frac{3}{2} y^{2}+y ; \\
& P_{22}=\left(-x^{2}+x y-\frac{3}{2} y^{2}+y\right)+\frac{1}{2}(x+y-1)^{2} ; \\
& P_{23}=\left(-x^{2}+x y-\frac{3}{2} y^{2}+y\right)+\frac{1}{2}(x+y-1)^{2}+2(y-x / 2-1 / 2)^{2} ; \\
& P_{24}=\left(-x^{2}+x y-\frac{3}{2} y^{2}+y\right)+2(y-x / 2-1 / 2)^{2} ; \\
& P_{31}=x y-\frac{3}{2} y^{2}+y ; \\
& P_{32}=P_{31} ; \\
& P_{33}=\left(x y-\frac{3}{2} y^{2}+y\right)+2(y-x / 2-1 / 2)^{2} ; \\
& P_{34}=P_{33} ; \\
& P_{41}=\left(x y+\frac{1}{2} y^{2}+y\right)-\frac{1}{2}(x+y+1)^{2} ; \\
& P_{42}=\left(x y+\frac{1}{2} y^{2}+y\right) ; \\
& P_{43}=P_{42} ; \\
& P_{44}=P_{41} ; \\
& P_{51}=\left(x^{2}-x y+\frac{3}{2} y^{2}+y\right)-\frac{1}{2}(x+y+1)^{2}-2(y-x / 2+1 / 2)^{2} ; \\
& P_{52}=\left(x^{2}-x y+\frac{3}{2} y^{2}+y\right)-2(y-x / 2+1 / 2)^{2} ; \\
& P_{53}=\left(x^{2}-x y+\frac{3}{2} y^{2}+y\right) ; \\
& P_{54}=\left(x^{2}-x y+\frac{3}{2} y^{2}+y\right)-\frac{1}{2}(x+y+1)^{2} ; \\
& P_{61}=\left(-x y+\frac{3}{2} y^{2}+y\right)-2(y-x / 2+1 / 2)^{2} ; \\
& P_{62}=\left(-x y+\frac{3}{2} y^{2}+y\right) ; \\
& P_{63}=P_{62} ; \\
& P_{64}=P_{61} \text {. }
\end{aligned}
$$

To compute the wavelets, we shall use the discrete $L_{2}$ inner product defined by (3.16) to determine the coefficients in (3.19). Let's recall the definition of the wavelet in (3.22)

$$
\begin{equation*}
\psi_{q, j^{\prime}}^{1}=\beta_{q, j^{\prime}}^{1} \phi_{q, j^{\prime}}^{1}(x)+\sum_{P \in \mathcal{N}_{q}, j=0, x, y} \beta_{q, j^{\prime}, P, j}^{1} \phi_{P, j}^{1}(x), j^{\prime}=0, x, y \tag{3.33}
\end{equation*}
$$

where vertex $q$ is the mid-point of two vertices in $\mathcal{N}_{q}$ in level 0 mesh. Recall that for a regular triangular mesh, level 0 mesh has the mesh size 1 , and after the refinement, the level 0 mesh becomes level 1 mesh with the mesh size $h_{1}=1 / 2$.
$\psi_{q, j^{\prime}}^{1}$ is required to be orthogonal to $\left\{\phi_{P, j}^{0}\right\}_{\left\{P \in \mathcal{N}_{q}, j=0, x, y\right\}}$ with respect to the discrete inner product $\langle\cdot, \cdot\rangle_{D_{2}(1)}$ for $V_{1}$. To find the suitable coefficients in (3.33) for wavelets, we need the following inner products,

$$
\left\langle\phi_{P^{\prime}, j^{\prime}}^{1}, \phi_{P, j}^{0}\right\rangle_{D_{2}(1)}, \quad P^{\prime} \in\left\{q, \mathcal{N}_{q}\right\}, P \in \mathcal{N}_{q}, j^{\prime}, j=0, x, y
$$

Let's recall the definition of the discrete $L_{2}$ inner product $\langle\cdot, \cdot\rangle_{D_{2}(1)}$ in level 1 before we carry out the computation.

$$
\begin{equation*}
\langle f, g\rangle_{D_{2}(1)}:=h_{1}^{2} \sum_{T \in \mathcal{T}_{1}} \sum_{P \in \mathcal{N}(T), j=0,1,2} \alpha_{T, P, j} \beta_{T, P, j}, \tag{3.34}
\end{equation*}
$$

where

$$
f(x)=\sum_{T \in \mathcal{T}_{1}} \sum_{P \in \mathcal{N}(T), j=0,1,2} \alpha_{T, P, j} \phi_{T, P, j}^{1}(x)
$$

and

$$
g(x)=\sum_{T \in \mathcal{T}_{1}} \sum_{P \in \mathcal{N}(T), j=0,1,2} \beta_{T, P, j} \phi_{T, P, j}^{1}(x) .
$$

Here, $\mathcal{T}_{1}$ is the triangulation of level 1 , and $\mathcal{N}(T)$ is the set of three vertices of the triangle $T$. Therefore, we shall write the basis functions in (3.33) and $\left\{\phi_{P, j}^{0}\right\}_{P \in \mathcal{N}_{q}, j=0, x, y}$ in terms of the local basis functions in $V_{1}$, and use their coefficients to compute the discrete inner product $\langle\cdot, \cdot\rangle_{D_{2}(1)}$.

In Figure 3.8, we shall compute 9 wavelets sitting on $q 1, q 2, q 3$, and their dilations and shifts form the wavelet basis. First, we give the explicit forms of nine wavelets in the following,


Figure 3.8: Computation of the wavelets

$$
\begin{aligned}
\psi_{q 3,0}^{1} & =\frac{-1}{8}\left(-8 \phi_{q 3,0}^{1}+4 \phi_{O, 0}^{1}-\phi_{O, x}^{1}+2 \phi_{O, y}^{1}+4 \phi_{P_{3}, 0}^{1}+\phi_{P_{3}, x}^{1}-2 \phi_{P_{3}, y}^{1}\right) \\
\psi_{q 3, x}^{1} & =\frac{-1}{2}\left(-2 \phi_{q 3, x}^{1}+\phi_{O, x}^{1}+\phi_{P_{3}, x}^{1}\right) \\
\psi_{q 3, y}^{1} & =\frac{1}{2}\left(2 \phi_{q 3, y}^{1}+2 \phi_{O, 0}^{1}-\phi_{O, x}^{1}+\phi_{O, y}^{1}-2 \phi_{P_{3}, 0}^{1}-\phi_{P_{3}, x}^{1}+\phi_{P_{3}, y}^{1}\right) \\
\psi_{q 2,0}^{1} & =\frac{-1}{8}\left(-8 \phi_{q 2,0}^{1}+4 \phi_{O, 0}^{1}+\phi_{O, x}^{1}+\phi_{O, y}^{1}+4 \phi_{P_{2}, 0}^{1}-\phi_{P_{2}, x}^{1}-\phi_{P_{2}, y}^{1}\right) \\
\psi_{q 2, x}^{1} & =\frac{1}{2}\left(2 \phi_{q 2, x}^{1}+2 \phi_{O, 0}^{1}+\phi_{O, y}^{1}-2 \phi_{P_{2}, 0}^{1}+\phi_{P_{2}, y}^{1}\right) \\
\psi_{q 2, y}^{1} & =\frac{1}{2}\left(2 \phi_{q 2, y}^{1}+2 \phi_{O, 0}^{1}+\phi_{O, x}^{1}-2 \phi_{P_{2}, 0}^{1}+\phi_{P_{2}, x}^{1}\right), \\
\psi_{q 1,0}^{1} & =\frac{-1}{8}\left(-8 \phi_{q 1,0}^{1}+4 \phi_{O, 0}^{1}+2 \phi_{O, x}^{1}-\phi_{O, y}^{1}+4 \phi_{P_{1}, 0}^{1}-2 \phi_{P_{1}, x}^{1}+\phi_{P_{1}, y}^{1}\right), \\
\psi_{q 1, x}^{1} & =\frac{-1}{2}\left(-2 \phi_{q 1, x}^{1}-2 \phi_{O, 0}^{1}-\phi_{O, x}^{1}+\phi_{O, y}^{1}+2 \phi_{P_{1}, 0}^{1}-\phi_{P_{1}, x}^{1}+\phi_{P_{1}, y}^{1}\right), \\
\psi_{q 1, y}^{1} & =\frac{-1}{2}\left(-2 \phi_{q 1, y}^{1}+\phi_{O, y}^{1}+\phi_{P_{1}, y}^{1}\right) .
\end{aligned}
$$

Next, we illustrate how to obtain the wavelets by using an example for computing $\psi_{q 3, x}^{1}$ in Figure 3.8.

Let

$$
\psi_{q 3, x}^{1}(x)=\beta_{q 3, x}^{1} \phi_{q 3, x}^{1}+\sum_{P \in\left\{O, P_{3}\right\}, j=0, x, y} \beta_{P, j}^{1} \phi_{P, j}^{1}(x) .
$$

Then $\psi_{q 3, x}^{1}$ shall be orthogonal to six basis functions in level 0 , i.e.,

$$
V:=\left\{\phi_{O, 0}^{0}, \phi_{O, x}^{0}, \phi_{O, y}^{0}, \phi_{P_{3}, 0}^{0}, \phi_{P_{3}, x}^{0}, \phi_{P_{3}, y}^{0}\right\} .
$$

Let $V^{1}:=\left\{\phi_{q 3, x}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{3}, 0}^{1}, \phi_{P_{3}, x}^{1}, \phi_{P_{3}, y}^{1}\right\}$, then the orthogonality between $\psi_{q 3, x}^{1}$ and $V$ can be written in the matrix form

$$
\left(\left\langle V^{T}, V^{1}\right\rangle_{D_{2}(1)}\right) \beta^{T}=0
$$

Where vector $\beta:=\left\{\beta_{q 3, x}^{1}, \beta_{O, 0}^{1}, \beta_{O, x}^{1}, \beta_{O, y}^{1}, \beta_{P_{3}, 0}^{1}, \beta_{P_{3}, x}^{1}, \beta_{P_{3}, y}^{1}\right\}$.
$\left(\left\langle V^{T}, V^{1}\right\rangle_{D_{2}(1)}\right)$ is a 6 by 7 matrix with $(i, j)$ element defined by $\langle f, g\rangle_{D_{2}(1)}$, where $f, g$ are $i$-th and $j$-th elements of the vectors $V, V^{1}$, respectively.

By the definition of discrete inner product, we shall write all involved functions in terms of local basis functions to obtain the matrix $\left(\left\langle V^{T}, V^{1}\right\rangle_{D_{2}(1)}\right)$. We shall give an example to compute $\left\langle\phi_{q 3, x}^{1}, \phi_{O, 0}^{0}\right\rangle_{D_{2}(1)}$ to illustrate how to get the required discrete inner products.

First, the support of $\phi_{q 3, x}^{1}$ is composed of six small triangles around the vertex $q 3$, such as $\triangle q 3 P_{3} w 1 . \phi_{O, 0}^{0}$ 's support are six triangles around $O$ with $P_{1}, P_{2}, \ldots, P_{6}$ as vertices. It is clear that the overlap of the supports of two functions is the support of $\psi_{q 3, x}^{1}$, i.e., six small triangles around vertex $q 3$.

Second, we shall represent two functions $\phi_{q 3, x}^{1}, \phi_{O, 0}^{0}$ in terms of local basis functions. For the derivatives, we focus on the six directional derivatives in six directions, i.e., $0, \pi / 4, \pi / 2$ and their opposite directions. It's easy to find that

$$
\frac{\partial \phi_{q 3, x}^{1}}{\partial x}(q 3)=1 / h_{1}=2, \frac{\partial \phi_{q 3, x}^{1}}{\partial y}(q 3)=0
$$

and

$$
\frac{\partial \phi_{q 3, x}^{1}}{\partial d_{\pi / 4}}=\frac{\partial \phi_{q 3, x}^{1}}{\partial x}(q 3) / \sqrt{2}=2 / \sqrt{2}
$$

where $d_{\pi / 4}$ denotes the direction with an anti-clockwise angle $\pi / 4$ and $h_{1}=$ $1 / 2$. Note that $\phi_{q 3, x}^{1}$ has all zeros data for its values and derivatives at vertices
other than $q 3$ and

$$
\frac{\partial \phi_{\triangle q 2 q 3 w 1, q 3,2}^{1}}{\partial d_{\triangle q 2 q 3 w 1, q 3,2}}=1 /\left|d_{\triangle q 2 q 3 w 1, q 3,2}\right|=1 /|q 3 w 1|=2 / \sqrt{2} .
$$

Thus we have the local representation for $\phi_{q 3, x}^{1}$ on its support. We list the representation for $\phi_{q 3, x}^{1}$ on $\triangle q 2 q 3 w 1$ as follows,

$$
\begin{equation*}
\left.\phi_{q 3, x}^{1}\right|_{\Delta q 2 q 3 w 1}=\phi_{\Delta q 2 q 3 w 1, q 3,1}^{1}+\phi_{\Delta q 2 q 3 w 1, q 3,2}^{1} . \tag{3.35}
\end{equation*}
$$

Next, we shall compute the local representation for $\phi_{0,0}^{0}$. Note that $\phi_{q 3, x}^{1}$ has all zeros data for its values and derivatives at all vertices other than $q 3$. We thus only interested in the directional derivatives of $\phi_{O, 0}^{0}$ at $q 3$ in six directions.

Since we have the explicit expression of $\phi_{O, 0}^{0}$, we may calculate the value and directional derivatives of $\phi_{O, 0}^{0}$ at $q 3$ by

$$
\phi_{O, 0}^{0}(q 3)=1 / 2, \frac{\partial \phi_{O, 0}^{0}}{\partial x}(q 3)=1, \frac{\partial \phi_{O, 0}^{0}}{\partial y}(q 3)=-2,
$$

and

$$
\frac{\partial \phi_{O, 0}^{0}}{\partial d_{\pi / 4}}=(1-2) / \sqrt{2}=-1 / \sqrt{2} .
$$

Therefore, we write the local representation of $\phi_{O, 0}^{0}$ on $\triangle q 2 q 3 w 1$ by

$$
\begin{equation*}
\left.\phi_{O, 0}^{0}\right|_{\triangle q 2 q 3 w 1}=\frac{1}{2} \phi_{\triangle q 2 q 3 w 1, q 3,0}^{1}+\frac{1}{2} \phi_{\triangle q 2 q 3 w 1, q 3,1}^{1}-\frac{1}{2} \phi_{\triangle q 2 q 3 w 1, q 3,2}^{1}+I, \tag{3.36}
\end{equation*}
$$

where $I$ includes the basis functions centered not at $q 3$.
By (3.35) and (3.36), we get the discrete inner product of $\phi_{q 3, x}^{1}$ and $\phi_{O, 0}^{0}$ on $\triangle q 2 q 3 w 1$ as follows

$$
\left.\left\langle\phi_{q 3, x}^{1}, \phi_{O, 0}^{0}\right\rangle_{D_{2}(1)}\right|_{\Delta q 2 q 3 w 1}=h_{1}^{2}\left(0 \times \frac{1}{2}+1 \times \frac{1}{2}+1 \times \frac{-1}{2}\right)=0 .
$$

Likewise, discrete inner products on other 5 sub-triangles in the support of $\phi_{q 3, x}^{1}$ can be done. In fact, we only need the information of directional derivatives of $\phi_{q 3, x}^{1}$ and $\phi_{O, 0}^{0}$ at $q 3$ in the computing. Taking the summation of inner products on all sub-triangles, we have

$$
\left\langle\phi_{q 3, x}^{1}, \phi_{O, 0}^{0}\right\rangle_{D_{2}(1)}=0 .
$$

In a similar way, we get the matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ and list it in the following

$$
h_{1}^{2}\left(\begin{array}{ccccccc}
0 & 1 & 0 & 0 & 0 & 0 & 0 \\
2 / 6 & 0 & 4 / 6 & 2 / 6 & 0 & 0 & 0 \\
1 / 6 & 0 & 2 / 6 & 4 / 6 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 & 1 & 0 & 0 \\
2 / 6 & 0 & 0 & 0 & 0 & 4 / 6 & 2 / 6 \\
1 / 6 & 0 & 0 & 0 & 0 & 2 / 6 & 4 / 6
\end{array}\right) .
$$

Null space of the column vectors is the solution for the wavelet $\psi_{q 3, x}^{1}$, and existence of the solution is proved previously in Section 3.3. From computation, we have

$$
\beta=\frac{-1}{2}\{-2,0,1,0,0,1,0\}^{T} .
$$

If we change $V^{1}$ to be $\left\{\phi_{q 3,0}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{3}, 0}^{1}, \phi_{P_{3}, x}^{1}, \phi_{P_{3}, y}^{1}\right\}$, then the associated matrix $\left(\left\langle V^{T}, V^{1}\right\rangle_{D_{2}(1)}\right)$ becomes

$$
h_{1}^{2}\left(\begin{array}{ccccccc}
1 / 2 & 1 & 0 & 0 & 0 & 0 & 0 \\
0 & 0 & 4 / 6 & 2 / 6 & 0 & 0 & 0 \\
1 / 8 & 0 & 2 / 6 & 4 / 6 & 0 & 0 & 0 \\
1 / 2 & 0 & 0 & 0 & 1 & 0 & 0 \\
0 & 0 & 0 & 0 & 0 & 4 / 6 & 2 / 6 \\
-1 / 8 & 0 & 0 & 0 & 0 & 2 / 6 & 4 / 6
\end{array}\right) .
$$

For this case, $\psi_{q 3,0}^{1}=V^{1} \beta^{T}$, where $\beta=\frac{-1}{8}\{-8,4,-1,2,4,1,-2\}$.
It's worth to point out that the above two matrices for $\left(\left\langle V^{T}, V^{1}\right\rangle_{D_{2}(1)}\right)$ are the same except for the first column. Therefore, in the following, we shall only list the first column of the associated matrix $\left(\left\langle V^{T}, V^{1}\right\rangle_{D_{2}(1)}\right)$ to compute other wavelets.

Let $V^{1}:=\left\{\phi_{q 3, y}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{3}, 0}^{1}, \phi_{P_{3}, x}^{1}, \phi_{P_{3}, y}^{1}\right\}$, and the first column of the matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{-1,1 / 6,-1 / 6,1,1 / 6,-1 / 6\}^{T}$. Thus, $\psi_{q 3, y}^{1}=$ $V^{1} \beta^{T}$, with $\beta=\frac{1}{2}\{2,2,-1,1,-2,-1,1\}$.

In the following, we give $\left\{\psi_{q 2, j}^{1}, \psi_{q 1, j}^{1}\right\}_{j=0, x, y}$.
1.) $\psi_{q, 2}^{1}$ :

Let $V^{1}=\left\{\phi_{q 2,0}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{2}, 0}^{1}, \phi_{P_{2}, x}^{1}, \phi_{P_{2}, y}^{1}\right\}$, and
$V=\left\{\phi_{O, 0}^{0}, \phi_{O, x}^{0}, \phi_{O, y}^{0}, \phi_{P_{2}, 0}^{0}, \phi_{P_{2, x}, x}^{0}, \phi_{P_{2, y}, y}^{0}\right\}$. The first column of the associated matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{1 / 2,1 / 8,1 / 8,1 / 2,-1 / 8,-1 / 8\}^{T}$. Then

$$
\psi_{q 2,0}^{1}=V^{1} \beta^{T},
$$

where

$$
\beta=\frac{-1}{8}\{-8,4,1,1,4,-1,-1\} .
$$

2.) $\psi_{q, x}^{1}$ :

Let $V^{1}=\left\{\phi_{q 2, x}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{2}, 0}^{1}, \phi_{P_{2}, x}^{1}, \phi_{P_{2}, y}^{1},\right\}$, and the first column of the associated matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{-1,-1 / 6,-2 / 6,1,-1 / 6,-2 / 6\}^{T}$. Then

$$
\psi_{q 2, x}^{1}=V^{1} \beta^{T},
$$

where

$$
\beta=\frac{1}{2}\{2,2,0,1,-2,0,1\} .
$$

3.) $\psi_{q 2, y}^{1}$ :

Let $V^{1}=\left\{\phi_{q 2, y}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{2}, 0}^{1}, \phi_{P_{2}, x}^{1}, \phi_{P_{2, y}}^{1}\right\}$, and the first column of the associated matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{-1,-2 / 6,-1 / 6,1,-2 / 6,-1 / 6\}^{T}$. Then

$$
\psi_{q 2, y}^{1}=V^{1} \beta^{T},
$$

where

$$
\beta=\frac{1}{2}\{2,2,1,0,-2,1,0\} .
$$

4.) $\psi_{q, 0}^{1}$ :

Let $V^{1}=\left\{\phi_{q, 0}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{1}, 0}^{1}, \phi_{P_{1}, x}^{1}, \phi_{P_{1}, y}^{1}\right\}$, and $V=\left\{\phi_{O, 0}^{0}, \phi_{O, x}^{0}, \phi_{O, y}^{0}, \phi_{P_{1}, 0}^{0}, \phi_{P_{1}, x}^{0}, \phi_{P_{1}, y}^{0}\right\}$. The first column of the associated matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{1 / 2,1 / 8,0,1 / 2,-1 / 8,0\}^{T}$. Then

$$
\psi_{q 1,0}^{1}=V^{1} \beta^{T},
$$

where

$$
\beta=\frac{-1}{8}\{-8,4,2,-1,4,-2,1\} .
$$

5.) $\psi_{q 1, x}^{1}$ :

Let $V^{1}=\left\{\phi_{q 1, x}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{1}, 0}^{1}, \phi_{P_{1}, x}^{1}, \phi_{P_{1}, y}^{1}\right\}$, and the first column of the associated matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{-1,-1 / 6,1 / 6,1,-1 / 6,1 / 6\}^{T}$. Then

$$
\psi_{q 1, x}^{1}=V^{1} \beta^{T}
$$

where

$$
\beta=\frac{-1}{2}\{-2,-2,-1,1,2,-1,1\} .
$$

6.) $\psi_{q 1, y}^{1}$ :

Let $V^{1}=\left\{\phi_{q 1, y}^{1}, \phi_{O, 0}^{1}, \phi_{O, x}^{1}, \phi_{O, y}^{1}, \phi_{P_{1}, 0}^{1}, \phi_{P_{1}, x}^{1}, \phi_{P_{1}, y}^{1},\right\}$, and the first column of the associated matrix $\left(\left\langle V^{T}, V 1\right\rangle_{D_{2}(1)}\right)$ is $\{0,1 / 6,2 / 6,0,1 / 6,2 / 6\}^{T}$. Then

$$
\psi_{q 1, y}^{1}=V^{1} \beta^{T}
$$

where

$$
\beta=\frac{-1}{2}\{-2,0,0,1,0,0,1\} .
$$

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## Appendix A

## Additive Schwarz-Type Preconditioner for Hermite Cubic Splines

Due to its built in parallelism as well as simple implementation, additive Schwarz type preconditioner has been received more and more attention recently [6, 29, 53, 57]. In Appendix A, we shall construct the additive Schwarz preconditioner for the Hermite cubic splines and prove that the preconditioned system has the uniformly bounded condition number. Hermite cubic splines are well known in the field of the approximation [20,34], and their $C^{1}$ continuity and high order approximation property make them attractive in practice.

This Appendix is divided into three parts. In section A.2, we sketch the basic framework of the additive Schwarz preconditioner. Hermite cubic splines and their properties shall be briefly reviewed in section A.3. Finally, we construct the nested finite element spaces with Hermite cubic splines, and show that the condition number of the preconditioned system by the additive Schwarz preconditioner for the Hermite cubic splines is uniformly bounded in section A.4.

## A. 1 Abstract additive Schwarz preconditioner

In this section, we introduce the notation, and the basic concepts of the additive Schwarz precndtioner we may use later. We are following the setting introduced in [7, 29].

Let $S_{0} \subset S_{1} \subset \cdots \subset S_{n}=S$ be a nested sequence of finite dimensional Hilbert spaces and

$$
S=\sum_{j=0}^{n} S_{j}
$$

where $n$ is a positive integer.
Let $a(\cdot, \cdot): S \times S \rightarrow \mathbb{R}$ be a positive definite and symmetric bilinear form with the properties

$$
a(v, w)=a(w, v) \quad \forall v, w \in S
$$

and

$$
a(v, v)>0
$$

Define $A: S \rightarrow S$, and

$$
a(v, w)=(A v, w) \quad \forall v, w \in S
$$

where $(\cdot, \cdot)$ is the scalar product in $S$.
Let each subspace $S_{j}, j=1, \ldots, n$, equipped with a positive definite and symmetric form $b_{j}(v, w)=\left(B_{j} v, w\right), \quad v, w \in S_{j}$ with $B_{j}: S_{j} \rightarrow S_{j}$. Finally, we define the operator $I_{j}: S_{j} \rightarrow S$ to be the nature injection operator, and its transpose is denoted by $I_{j}^{t}$.

The abstract additive Schwarz preconditioner can be written as

$$
\begin{equation*}
B=\sum_{j=0}^{n} I_{j} B_{j}^{-1} I_{j}^{t} \tag{A.1}
\end{equation*}
$$

Remark 1. In practice, $a(\cdot, \cdot)$ usually is the bilinear form introduced from the given second order elliptic problem, and $A$ corresponds to the stiffness matrix. $I_{j}$ is referred to as the transformation matrix for the basis in $S_{j}$ and the basis in $S . I_{j}^{t}$ is the transpose of $I_{j} . b_{j}(\cdot, \cdot)$ is closely related to the scalar product
in $S_{j}$, and $B_{j}$ usually can be represented as a diagonal matrix. Therefore, the additive Schwarz preconditioner (A.1) can be easily implemented.

Once we have a sequence of nested subspaces, we may estimate the maximum and minimum eigenvalues by the following theorem

Theorem A.1. The maximum and minimum eigenvalues of $B A$ can be characterized by

$$
\lambda_{\max }(B A):=\max _{0 \neq v \in S} \frac{a(v, v)}{\min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} b_{j}\left(v_{j}, v_{j}\right)},
$$

and

$$
\lambda_{\min }(B A):=\min _{0 \neq v \in S} \frac{a(v, v)}{\min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} b_{j}\left(v_{j}, v_{j}\right)} .
$$

Proofs of the theorem can be found in several sources [7, 29, 41].
We may find that whether the additive Schwarz preconditioner works well or not is depending on the ratio (the condition number) of $\lambda_{\max }$ to $\lambda_{\text {min }}$ in the Theorem A.1.

## A. 2 Hermite cubic splines

Recall that the Hermite cubic splines $\phi_{1}$ and $\phi_{2}$ be given by

$$
\phi_{1}(x):=\left\{\begin{array}{cll}
\phi_{10}:=(x+1)^{2}(1-2 x) & \text { for } & x \in[-1,0] \\
\phi_{11}:=(1-x)^{2}(2 x+1) & \text { for } & x \in[0,1] \\
0 & \text { for } & x \in \mathbb{R} \backslash[-1,1]
\end{array}\right.
$$

and

$$
\phi_{2}(x):=\left\{\begin{array}{cl}
\phi_{20}:=x(x+1)^{2} & \text { for } x \in[-1,0] \\
\phi_{21}:=x(x-1)^{2} & \text { for } x \in[0,1] \\
0 & \text { for } x \in \mathbb{R} \backslash[-1,1]
\end{array}\right.
$$

Then

$$
\phi_{1}(j)=\delta(j), \quad \phi_{1}^{\prime}(j)=0, \quad \phi_{2}(j)=0, \quad \phi_{1}^{\prime}(j)=\delta(j), j \in \mathbb{Z}
$$

where

$$
\delta(j):=\left\{\begin{array}{ll}
0, & 0 \neq j \in \mathbb{Z} \\
1, & j=0
\end{array} .\right.
$$

Applications of the Hermite cubic splines in solving ODE numerically can be found in [34].

On a mesh with $0<x_{0}<x_{1}<\cdots<x_{N}=1$, we may scale and shift $\phi_{1}, \phi_{2}$ to construct the basis functions for the finite dimensional space as follows

$$
\phi_{1, i}(x):=\left\{\begin{array}{cll}
\phi_{10}\left(\frac{x-x_{i}}{x_{i}-x_{i-1}}\right) & \text { for } & x \in\left[x_{i-1}, x_{i}\right] \\
\phi_{11}\left(\frac{x-x_{i}}{x_{i+1}-x_{i}}\right) & \text { for } & x \in\left[x_{i}, x_{i+1}\right] \\
0 & \text { for } & x \in \mathbb{R} \backslash\left[x_{i-1}, x_{i+1}\right]
\end{array}\right.
$$

and

$$
\phi_{2, i}(x):=\left\{\begin{array}{cll}
\phi_{20}\left(\frac{x-x_{i}}{x_{i}-x_{i-1}}\right) & \text { for } & x \in\left[x_{i-1}, x_{i}\right] \\
\phi_{21}\left(\frac{x-x_{i}}{x_{i+1}-x_{i}}\right) \frac{\left(x_{i+1}-x_{i}\right)}{\left(x_{i}-x_{i-1}\right)} & \text { for } & x \in\left[x_{i}, x_{i+1}\right], \\
0 & \text { for } & x \in \mathbb{R} \backslash\left[x_{i-1}, x_{i+1}\right],
\end{array}\right.
$$

where $i=0, \ldots, N$.
We may verify that $\phi_{1, i}\left(x_{j}\right)=\delta(i-j), \phi_{1, i}^{\prime}\left(x_{j}\right)=0, \phi_{2, i}\left(x_{j}\right)=0, \phi_{2, i}^{\prime}\left(x_{j}\right)=$ $\delta(i-j) /\left(x_{i}-x_{i-1}\right), \quad i, j=0, \ldots, N$.

Now we introduce one lemma on the Hermite cubic splines.

Lemma A.1. Let $v=\alpha_{0} \phi_{11}(x)+\alpha_{1} \phi_{10}(x-1)+\beta_{0} \phi_{20}(x)+\beta_{1} \phi_{21}(x-1)$, then we have

$$
\begin{equation*}
C_{1}\|v\|_{L_{2}(0,1)}^{2} \leq \alpha_{0}^{2}+\alpha_{1}^{2}+\beta_{0}^{2}+\beta_{1}^{2} \leq C_{2}\|v\|_{L_{2}(0,1)}^{2}, \tag{A.2}
\end{equation*}
$$

$$
\begin{equation*}
C_{3}\left(\beta_{0}^{2}+\beta_{1}^{2}\right) \leq\left\|v^{\prime}\right\|_{L_{2}(0,1)}^{2}, \tag{A.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|v^{\prime}\right\|_{L_{2}(0,1)} \leq C_{4}\|v\|_{L_{2}(0,1)} \tag{A.4}
\end{equation*}
$$

where $C_{1}, C_{2}, C_{3}$ and $C_{4}$ are four constants independent of $v$, and $L_{2}$ norm on the interval $I \subset[0,1]$ is defined to be $\|v\|_{L_{2}(I)}:=\left(\int_{I}|v(x)|^{2} d x\right)^{1 / 2}$.

Proof. Let the vector $\alpha:=\left(\alpha_{0}, \alpha_{1}, \beta_{0}, \beta_{1}\right)^{t}$, then

$$
\|v\|_{L_{2}(0,1)}^{2}=\alpha^{t} D \alpha
$$

where

$$
D=\left(\begin{array}{cccc}
\frac{13}{35} & \frac{9}{70} & \frac{-11}{210} & \frac{13}{420} \\
\frac{9}{70} & \frac{13}{35} & \frac{-13}{420} & \frac{11}{210} \\
\frac{-11}{210} & \frac{-13}{420} & \frac{1}{105} & \frac{-1}{110} \\
\frac{13}{420} & \frac{11}{210} & \frac{-1}{140} & \frac{1}{105}
\end{array}\right) .
$$

Note that $D$ is symmetric and positive definite (i.e., its eigenvalues are strictly positive and bounded). Then (A.2) holds true.

Likewise, we have

$$
\left\|v^{\prime}\right\|_{L_{2}(0,1)}^{2}=\alpha^{t} D_{1} \alpha
$$

where

$$
D_{1}=\left(\begin{array}{cccc}
\frac{6}{5} & \frac{-6}{5} & \frac{-1}{10} & \frac{-1}{10} \\
\frac{-6}{5} & \frac{6}{5} & \frac{1}{10} & \frac{1}{10} \\
\frac{-1}{10} & \frac{1}{10} & \frac{2}{15} & \frac{-1}{30} \\
\frac{-1}{10} & \frac{1}{10} & \frac{-1}{30} & \frac{2}{15}
\end{array}\right)
$$

Note that $D_{1}-0.082 D_{2}$ is symmetric and with the eigenvalues nonnegative, where $D_{2}$ is a diagonal matrix with the diagonal entries $(0,0,1,1)$. Then (A.3) follows if we set $C_{3}=0.082$.

Last inequality (A.4) is true because

$$
\left\|v^{\prime}\right\|_{L_{2}(0,1)} \leq C_{4}\left(\alpha_{0}^{2}+\alpha_{1}^{2}+\beta_{0}^{2}+\beta_{1}^{2}\right)
$$

If we note that $D_{1}$ 's maximum eigenvalue is bounded, then, by (A.2), (A.4) holds.

This completes the proof.

## A. 3 Additive Schwarz preconditioner for the Hermite cubic splines

For the given second order elliptic model two points boundary value problem

$$
-\left(p(x) u(x)^{\prime}\right)^{\prime}+q(x) u(x)=f(x), \quad x \in(0,1)
$$

with the boundary conditions

$$
u(0)=u(1)=0
$$

we may define the bilinear form $a(\cdot, \cdot)$ to be

$$
a(v, w)=\int_{0}^{1} p(x) v^{\prime}(x) w^{\prime}(x) d x+\int_{0}^{1} q(x) v(x) w(x) d x, v, w \in H_{0}^{1}(0,1)
$$

where $p(x)>0, q(x) \geq 0$ for $x \in(0,1)$, and $H_{0}^{1}(0,1)$ is the usual Sobolev space with the norm, semi-norm denoted by $\|\cdot\|_{1},|\cdot|_{1}$, respectively. It's well known that

$$
a(v, v) \simeq\|v\|_{1}^{2}
$$

Here and in what follows, we use $X \simeq Y$ to denote the equivalence of the two terms $X$ and $Y(X, Y$ can be bounded each other by multiply by some constants independent of the mesh.), and let $C, C_{i}(i=1,2 \ldots)$ denote the generic constants independent of the mesh.

For $H_{0}^{1}(0,1), H^{1}$ semi-norm is an equivalent norm to $H^{1}$ norm.
The Galerkin method is to seek an element $u$ in $H_{0}^{1}(0,1)$, such that

$$
\begin{equation*}
a(u, v)=(f, v) \quad \forall v \in H_{0}^{1}(0,1) \tag{A.5}
\end{equation*}
$$

where $(v, w):=\int_{0}^{1} v(x) w(x) d x$ is the traditional inner product for the realvalued function space $L_{2}(0,1)$.

If we have a finite dimensional subspace $S \subset H_{0}^{1}(0,1)$, then the finite element method is to seek an element $u_{n} \in S$ such that

$$
\begin{equation*}
a\left(u_{n}, v\right)=(f, v) \quad \forall v \in S \tag{A.6}
\end{equation*}
$$

Now we construct the finite element space based on the Hermite cubic splines. For the convenience of statement, we focus on the uniform mesh, although quasi-uniform mesh also admits the additive Schwarz preconditioner for the Hermite cubic splines.

Let $\phi_{j, k}^{e}(x):=\phi_{e}\left(2^{j} x-k\right), e=1,2$, where $j, k$, known as scales and shifts, are both nonnegative integers. Then we may check that $\operatorname{supp}\left(\phi_{j, k}^{e}\right)=[k-1, k+$ $1] / 2^{j}$.

Let the finite dimensional space $S_{j}$ be the linear span of the basis functions $\left\{\phi_{j, k}^{e}\right\}, k=1, \ldots, 2^{j}-1$ for $e=1$, and $k=0, \ldots, 2^{k}$ for $e=2$. Then

$$
S=S_{n}=\sum_{j=0}^{n} S_{j}, \quad S \subset H_{0}^{1}(0,1)
$$

We may write the basis function in a vector by $\Phi_{j}:=\left\{\varphi_{j, l}\right\}$, where $\varphi_{j, 2 k}=$ $\phi_{j, k}^{2}, k=0, \ldots, 2^{j}-1, \varphi_{j, 2 k-1}=\phi_{j, k}^{1}, k=1, \ldots, 2^{j}-1$, and $\varphi_{j, 2^{j+1}-1}=\phi_{j, 2^{j}}^{2}$.

The corresponding stiffness matrix is generated by $A:=\left(a\left(\varphi_{n, k 1}, \varphi_{n, k 2}\right)\right)$, and the transformation matrix $I_{j}$ can be obtained from the so called refinement equations

$$
\Phi(x)=\sum_{k=-1}^{1} R_{k} \Phi(2 x-k)
$$

where the vector of functions $\Phi$ is defined to be $\left(\phi^{0}(x), \phi^{1}(x)\right)^{T}$ and the matrices

$$
R_{-1}=\left[\begin{array}{cc}
\frac{1}{2} & \frac{3}{4} \\
\frac{-1}{8} & \frac{-1}{8}
\end{array}\right], R_{0}=\left[\begin{array}{ll}
1 & 0 \\
0 & \frac{1}{2}
\end{array}\right], R_{1}=\left[\begin{array}{cc}
\frac{1}{2} & \frac{-3}{4} \\
\frac{1}{8} & \frac{-1}{8}
\end{array}\right]
$$

In other words, every basis function $\varphi_{j, k}$ can be written as a linear combination of no more than 6 basis functions in $\left\{\varphi_{j+1, k}\right\}$. Therefore, $S_{j} \subset S_{j+1}$. Moreover, we have the $2^{j+2}$ by $2^{j+1}$ transformation matrix $T_{j}$, such that

$$
\Phi_{j}^{t}=\Phi_{j+1}^{t} T_{j}
$$

Thus, $I_{j}$ can be written as

$$
I_{j}:=T_{n-1} T_{n-2} \cdots T_{j}, \quad j=0, \ldots, n-1
$$

and $I_{n}$ is just the $2^{n+1}$ by $2^{n+1}$ identity matrix. Furthermore

$$
\begin{equation*}
\Phi_{j}^{t}=\Phi_{n}^{t} I_{j} \quad j=0, \ldots, n \tag{A.7}
\end{equation*}
$$

For any element $v_{j} \in S_{j}$, we may write it as a linear combination of the basis functions, i.e.

$$
v_{j}=\sum_{k=0}^{2^{j+1}-1} v_{j, k} \varphi_{j, k}
$$

Denote by $\mathbf{v}_{\mathbf{j}}$ the vector of the coefficients $\left\{v_{j, k}\right\}$ associated to $v_{j} \in S_{j}$. Then by Lemma A.1, we have

$$
h_{j} \mathbf{v}_{\mathbf{j}}^{T} \mathbf{v}_{\mathbf{j}} \simeq\|v\|_{L_{2}(0,1)}^{2}
$$

where $\mathbf{v}_{\mathbf{j}_{1}}{ }^{T} \mathbf{w}_{\mathbf{j}_{2}}$ is the usual vector product and $h_{j}$, known as the mesh size of $S_{j}$, is $2^{-j}$.

Let's take a look at $I_{j}$ again. For a given $v_{j} \in S_{j}$, we associate it with its coefficient vector $\mathbf{v}_{\mathbf{j}}$, and $I_{j} v_{j} \in S_{n}$. Here $I_{j}$ is an injection mapping $v_{j} \in S_{j}$ naturally into $S_{n}$. With some ambiguity, we use the same notation $I_{j}$ in (A.7) as a transformation matrix mapping a vector in $\mathbb{R}^{2^{j+1}}$ to $\mathbb{R}^{2^{n+1}}$. More precisely, if $\mathbf{v}_{\mathbf{n}}$ is the vector of coefficients of $v_{j}$ with the basis functions in $S_{n}$ (i.e., $v_{j}=\Phi_{n}^{t} \mathbf{v}_{\mathbf{n}}$ ), then

$$
\mathbf{v}_{\mathbf{n}}=I_{j} \mathbf{v}_{\mathbf{j}}
$$

Thus, vectors $\mathbf{v}_{\mathbf{j}}, \mathbf{v}_{\mathbf{n}}$ can be think of as two representations of the same function in $S$ in terms of bases in $S_{j}, S_{n}$, respectively. $I_{j}$, as a transformation matrix, connects such basis change. It's important in the numerical implementation.

Let the bilinear form $b_{j}(\cdot, \cdot)$ on $S_{j}$ be

$$
b_{j}\left(v_{j}, w_{j}\right)=h_{j}^{-1} \mathbf{v}_{\mathbf{j}}^{T} \mathbf{w}_{\mathbf{j}}, \quad v_{j}, w_{j} \in S_{j}
$$

Then $B_{j}^{-1}$, written in the matrix form, is a $2^{j+1}$ by $2^{j+1}$ identity matrix multiply by $h_{j}^{-1}=2^{j}$.

Now the additive Schwarz preconditioner can be defined in the matrix form by

$$
B=\sum_{j=0}^{n} 2^{-j} I_{j} I_{j}^{T}
$$

Remark 2. The computational work for $B \mathbf{v}_{\mathbf{n}}$ can be calculated as $O\left(2^{n+1}\right)$ if we note that the computational work for $T_{j} \mathbf{v}_{\mathbf{j}}$ is $O\left(2^{j+1}\right)$, and $2^{j+1}$ is the dimension of the space $S_{j}$.

After we introduce the additive Schwarz preconditioner, we may estimate the maximum and the minimum eigenvalues of the product of the matrices $B$ and $A$. By Theorem A.1, we need to estimate the ratio of $a(v, v)$ to $\min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}}{ }^{T} \mathbf{v}_{\mathbf{j}}$.

For Hermite cubic splines, we have the following theorem,

Theorem A.2. For the finite dimensional space $\left\{S_{j}\right\}$ generated by the Hermite cubic splines, we have

$$
\lambda_{\max }:=\max _{0 \neq v \in S} \frac{a(v, v)}{\min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}} \mathbf{v}_{\mathbf{j}}}=O(1)
$$

and

$$
\lambda_{\min }:=\min _{0 \neq v \in S} \frac{a(v, v)}{\min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}} \mathbf{T}_{\mathbf{j}}}=O(1)
$$

Before we prove the theorem, we need three Lemmas. Let $Q_{j}: S_{n} \rightarrow S_{j}, j \geq$ $0\left(Q_{-1}=0\right)$ be the orthogonal projection operator, i.e.,

$$
\left(Q_{j} v, w_{j}\right)=\left(v, w_{j}\right), \quad \forall w_{j} \in S_{j}
$$

where $v$ is an element in $S_{n}$. Then for the sequence of subspaces $\left\{S_{j}\right\}$, we have

Lemma A.2. For $v \in S$, we have

$$
\|v\|_{1}^{2} \simeq \sum_{j=0}^{n} 4^{j}\left\|\left(Q_{j}-Q_{j-1}\right) v\right\|_{L_{2}}^{2}
$$

As a result of norm equivalence, it's well known in the literature. It's proofs may be found in [21, 41].

Lemma A.3. For any $v_{j} \in S_{j}, j=0,1, \ldots, n$, we have

$$
\begin{equation*}
\left\|v_{j}\right\|_{L_{2}}^{2} \simeq \sum_{k=0}^{2^{j+1}} v_{j, k}^{2} h_{j}=2^{-j} \mathbf{v}_{\mathbf{j}}^{t} \mathbf{v}_{\mathbf{j}} \tag{A.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|v_{j}^{\prime}\right\|_{L_{2}}^{2} \leq C \sum_{k=0}^{2^{j+1}}\left(v_{j, k} / h_{j}\right)^{2} h_{j}=C 2^{j} \mathbf{v}_{\mathbf{j}}^{t} \mathbf{v}_{\mathbf{j}} \tag{A.9}
\end{equation*}
$$

Proof. If we note that $\left\|\phi_{j, k}^{e}\right\|_{L_{2}(0,1)} \simeq h_{j}=2^{-j}$, and $\left\|\left(\phi_{j, k}^{e}\right)^{\prime}\right\|_{L_{2}(0,1)} \leq C 2^{j}$, then (A.8, A.9) can be obtained directly by Lemma A.1. (A.8) usually is referred to as the stability of the Hermite cubic splines in the sense of $L_{2}$, and (A.9) is the inverse inequality.

The next Lemma is a Cauchy-Schwarz type inequality,

## Lemma A.4.

$$
\int_{0}^{1} v_{j}^{\prime}(x) w_{k}^{\prime}(x) d x \leq C 2^{-|j-k| / 2}\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}}\right)\left(h_{k}^{-1}\left\|v_{k}\right\|_{L_{2}}\right), \quad \forall v_{j} \in S_{j}, w_{k} \in S_{k}
$$

Proof. For the case $k=j$, by Cauchy-Schwarz inequality

$$
\int_{0}^{1} v_{j}^{\prime}(x) w_{k}^{\prime}(x) d x \leq\left\|v_{j}^{\prime}\right\|_{L_{2}}\left\|w_{k}^{\prime}\right\|_{L_{2}} \leq C\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}}\right)\left(h_{k}^{-1}\left\|v_{k}\right\|_{L_{2}}\right)
$$

where we use the inverse inequality in the last step.
For the case $j<k$, we consider one sub-interval, say $\left(m / 2^{j},(m+1) / 2^{j}\right)$ in the mesh for $S_{j}$. Furthermore, let $\alpha_{1}=v_{j}\left(m / 2^{j}\right), \alpha_{2}=v_{j}\left((m+1) / 2^{j}\right), \beta_{1}=$ $h_{j} v_{j}^{\prime}\left(m / 2^{j}\right), \beta_{2}=h_{j} v_{j}^{\prime}\left((m+1) / 2^{j}\right)$. Then, on the interval $\left(m / 2^{j},(m+1) / 2^{j}\right)$, $v_{j}$ can be written as

$$
v_{j}=\alpha_{1} \phi_{j, m}^{1}+\alpha_{2} \phi_{j, m+1}^{1}+\beta_{1} \phi_{j, m+1}^{2}+\beta_{2} \phi_{j, m+1}^{2}
$$

Let $a=m / 2^{j}, b=(m+1) / 2^{j}$. Then

$$
\int_{a}^{b} v_{j}^{\prime}(x) w_{k}^{\prime}(x) d x=\left.v_{j}^{\prime} w_{k}\right|_{a} ^{b}-\int_{a}^{b} v_{j}^{\prime \prime}(x) w_{k}(x) d x
$$

First, we estimate the term $\left.v_{j}^{\prime} w_{k}\right|_{a} ^{b}$ by

$$
\begin{aligned}
\left.v_{j}^{\prime} w_{k}\right|_{a} ^{b} & \leq \frac{1}{2}\left(\left(v_{j}^{\prime}(a)\right)^{2}+\left(v_{j}^{\prime}(b)\right)^{2}\right)^{1 / 2}\left(\left(w_{k}(a)\right)^{2}+\left(w_{k}(b)\right)^{2}\right)^{1 / 2} \\
& \leq C\left(\left\|v_{j}^{\prime}\right\|_{L_{2}(a, b)} h_{j}^{-1 / 2}\right)\left(\left\|w_{k}\right\|_{L_{2}(a, b)} h_{k}^{-1 / 2}\right)
\end{aligned}
$$

Since by the Lemma A. 1 (after a scale), we have

$$
\left(\left(v_{j}^{\prime}(a)\right)^{2}+\left(v_{j}^{\prime}(b)\right)^{2}\right) h_{j} \leq C\left\|v_{j}^{\prime}\right\|_{L_{2}(a, b)}^{2}
$$

and

$$
\left(\left(w_{k}(a)\right)^{2}+\left(w_{k}(b)\right)^{2}\right) h_{k} \leq C\left\|w_{k}\right\|_{L_{2}(a, b)}^{2} .
$$

Hence, by inverse inequality,

$$
\begin{aligned}
\left(\left\|v_{j}^{\prime}\right\|_{L_{2}(a, b)} h_{j}^{-1 / 2}\right)\left(\left\|w_{k}\right\|_{L_{2}(a, b)} h_{k}^{-1 / 2}\right) & \leq C\left(h_{j}^{-3 / 2}\left\|v_{j}\right\|_{L_{2}(a, b)}\right)\left(h_{k}^{-1 / 2}\left\|w_{k}\right\|_{L_{2}(a, b)}\right) \\
& \leq C\left(h_{k} / h_{j}\right)^{1 / 2}\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}(a, b)}\right)\left(h_{k}^{-1}\left\|v_{k}\right\|_{L_{2}(a, b)}\right)
\end{aligned}
$$

Note that $h_{j}=2^{-j}, h_{k}=2^{-k}$. Then

$$
\left.v_{j}^{\prime} w_{k}\right|_{a} ^{b} \leq C 2^{-|j-k| / 2}\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}}\right)\left(h_{k}^{-1}\left\|v_{k}\right\|_{L_{2}}\right) .
$$

Second, we estimate the term $-\int_{a}^{b} v_{j}^{\prime \prime}(x) w_{k}(x)$.

$$
\begin{aligned}
\left|-\int_{a}^{b} v_{j}^{\prime \prime}(x) w_{k}(x)\right| & \leq C\left\|v_{j}^{\prime \prime}\right\|_{L_{2}(a, b)}\left\|w_{k}\right\|_{L_{2}(a, b)} \\
& \leq C h_{j}^{-2}\left\|v_{j}\right\|_{L_{2}(a, b)}\left\|w_{k}\right\|_{L_{2}(a, b)} \\
& \leq C 2^{-|j-k| / 2}\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}(a, b)}\right)\left(h_{k}^{-1}\left\|v_{k}\right\|_{L_{2}(a, b)}\right) .
\end{aligned}
$$

If we add up the estimates on all intervals and apply the Cauchy-Schwarz inequality, then Lemma A. 4 follows. Thus completes the proof.

Now we give the proof of Theorem

## [theoremA.2]

Proof of Theorem A.2.

Let $v=\sum_{j=0}^{n} I_{j} v_{j}$, then we have

$$
\begin{align*}
a(v, v) & \simeq C\left\|v^{\prime}\right\|_{L_{2}(0,1)}^{2} \\
& =C\left(\sum_{j=0}^{n} v_{j}^{\prime}, \sum_{k=0}^{n} v_{k}^{\prime}\right) \\
& \leq C \sum_{j, k=0}^{n} 2^{-|j-k| / 2}\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}}\right)\left(h_{k}^{-1}\left\|v_{k}\right\|_{L_{2}}\right) \\
& \leq \sum_{j=0}^{n}\left(h_{j}^{-1}\left\|v_{j}\right\|_{L_{2}}\right)^{2} \\
& \simeq \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}} \mathbf{v}_{\mathbf{j}} . \tag{bylemmaA.3}
\end{align*}
$$

Since the splitting of $v$ is arbitrary, this implies that

$$
\begin{equation*}
a(v, v) \leq C \min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}}^{T} \mathbf{v}_{\mathbf{j}} \tag{A.10}
\end{equation*}
$$

Now let $v_{j}=\left(Q_{j}-Q_{j-1}\right) v, j=0, \ldots, n$, then by Lemma A.2, we have

$$
\begin{equation*}
a(v, v) \simeq \sum_{j=0}^{n} 4^{j}\left\|v_{j}\right\|_{L_{2}}^{2} \simeq \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}}{ }^{T} \mathbf{v}_{\mathbf{j}} . \tag{A.11}
\end{equation*}
$$

Combing (A.10) with (A.11) yields that

$$
a(v, v) \simeq \min _{v=\sum_{j=0}^{n} I_{j} v_{j}, v_{j} \in S_{j}} \sum_{j=0}^{n} 2^{j} \mathbf{v}_{\mathbf{j}}^{T} \mathbf{v}_{\mathbf{j}}
$$

Thus completes the proof.

Remark 3. For quasi-uniform mesh, note that $h_{j} \simeq 2^{-j}, j=0, \ldots, n$. Then we can obtain the same result on the additive Schwarz preconditioner for the Hermite cubic splines using the same proof in the note.

Finally, we show the numerical results of the additive Schwarz preconditioner for the model problem

$$
u^{\prime \prime}(x)=f(x) \quad x \in(0,1)
$$

with the boundary conditions

$$
u(0)=u(1)=0
$$

The bilinear form arising from the elliptic problem is $a(v, w)=\int_{0}^{1} v^{\prime}(x) w^{\prime}(x) d x$, and thus $A$ is the stiffness matrix with $\left(k_{1}, k_{2}\right)$ entry $\int_{0}^{1} \varphi_{n, k_{1}}^{\prime}(x) \varphi_{n, k_{2}}^{\prime}(x) d x$. To obtain better results on the condition number, we normalize $\varphi_{j, k}, j=$ $0, \ldots, n, k=0, \ldots, 2^{j+1}-1$, such that

$$
\int_{0}^{1}\left|\varphi_{j, k}^{\prime}(x)\right|^{2} d x=2^{-j}
$$

The additive Schwarz preconditioner is given in (A.1). The condition numbers with respect to different $n$ are listed in the following table. The numerical results confirm the claims in theorem A.2.

| $n$ | 6 | 7 | 8 | 9 | 10 | 11 | 12 |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\kappa$ | 4.62 | 4.71 | 4.78 | 4.83 | 4.86 | 4.89 | 4.89 |

Table A.1: Condition numbers( $\kappa$ ) of $B A$ with different $n$

