# Mean value properties of harmonic functions on Sierpinski gasket type fractals 

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#### Abstract

In this paper, we establish an analogue of the classical mean value property for both the harmonic functions and some general functions in the domain of the Laplacian on the Sierpinski gasket. Furthermore, we extend the result to some other p.c.f. fractals with Dihedral-3 symmetry.


Keywords. Sierpinski gasket, Laplacian, harmonic function, mean value property, analysis on fractals.
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## 1 Introduction

It is well known that harmonic functions i.e. solutions of the Laplace equation $\Delta u=0$, where $\Delta=\sum_{i=1}^{d} \frac{\partial^{2}}{\partial x_{i}^{2}}$, possess the mean value property. Namely, if $u$ is harmonic on a domain $\Omega \subset \mathbf{R}^{d}$, then for every closed ball $B_{r}(x) \subset \Omega$ of a center $x \in \Omega$ and radius $r>0$ the average of $u$ over $B_{r}(x)$ equals to the value of $x$, i.e.,

$$
\frac{1}{\left|B_{r}(x)\right|} \int_{B_{r}(x)} u(y) d y=u(x),
$$

where $\left|B_{r}(x)\right|$ is the volume of the ball $B_{r}(x)$. There is a similar statement for mean values on spheres. More generally, if $u$ is not assumed harmonic but $\Delta u$ is a continuous function, then

$$
\begin{equation*}
\lim _{r \rightarrow 0} \frac{1}{r^{2}}\left(\frac{1}{\left|B_{r}(x)\right|} \int_{B_{r}(x)} u(y) d y-u(x)\right)=c_{n} \Delta u(x) \tag{1.1}
\end{equation*}
$$

for the appropriate dimensional constant $c_{n}$.
What are the fractal analogs of these results? The analytic theory on p.c.f. fractals was developed by Kigami [2, 3, 4] following the work of several probabilists who constructed stochastic processes analogous to Brownian motion, thus obtaining a Laplacian

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indirectly as the generator of the process. See the book of Barlow [1] for an account of this development. Since analysis on fractals has been made possible by the analytic definition of Laplacian, it is natural to explore the properties of these fractal Laplacians that are natural analogs of results that are known for the usual Laplacian. As for the fractal analog of the mean value property, we do not want to specify in advance the nature of the sets on which we do the averaging. So if $K$ is a fractal and $x \in K$, we would like to know that there is a sequence of sets $B_{k}(x)$ containing $x$ with $\bigcap_{k} B_{k}(x)=\{x\}$ such that

$$
\frac{1}{\mu\left(B_{k}(x)\right)} \int_{B_{k}(x)} u(y) d y=u(x)
$$

for every harmonic function $u$. Moreover, for general $u$ not assumed harmonic, is there a formula analogous to (1.1)?

In the present paper, we will mainly deal with the Sierpinski gasket $\mathcal{S G}$. This set is a key example of fractals on which a well established theory of Laplacian exists [1, 2, 3, 4, 15, 6]. Since the mean value property plays a very important role in the usual theory of harmonic functions, it is of independent interest to understand the similar property of harmonic functions on the Sierpinski gasket. We will prove that for each point $x \in \mathcal{S G} \backslash V_{0}$, ( $V_{0}$ is the boundary of $\mathcal{S G}$.) there is a sequence of mean value neighborhoods $B_{k}(x)$ depending only on the location of $x$ in $\mathcal{S G} .\left\{B_{k}(x)\right\}$ forms a system of neighborhoods of the point $x$ satisfying $\bigcap_{k} B_{k}(x)=\{x\}$. On such sequences, we get the fractal analogs of the mean value properties of both the harmonic functions and the general functions which belong to the domain of the fractal Laplacian satisfying some natural continuity assumption. We also investigate the extent to which our method can be applicable to other p.c.f. selfsimilar sets, but it seems that it strongly depends on the symmetric properties of both the geometric structure and the harmonic structure of the fractals.

The paper is organized as follows: In Section 2 we briefly introduce some key notions from analysis on the Sierpinski gasket. In Section 3 and Section 4, we prove the mean value property for harmonic functions and general functions on $\mathcal{S G}$ respectively. Section 5 contains a further extension of the mean value property to p.c.f. self-similar fractals with Dihedral-3 symmetry. An interesting open question is to what extent the results of Section 4 can be extended to this class of fractals.

## 2 Analysis on the Sierpinski gasket

For the convenience of the reader, we collect some key facts from analysis on $\mathcal{S G}$ that we need to state and prove our results. These come from Kigami's theory of analysis
on fractals, and may be found in [2, 3, 4]. An elementary exposition may be found in [5. 6]. Recall that $\mathcal{S G}$ is the attractor of the ifs (iterated function system) in the plane consisting of three homotheties $\left\{F_{0}, F_{1}, F_{2}\right\}$ with contraction ratio $1 / 2$ and fixed points equal to the three vertices $\left\{q_{0}, q_{1}, q_{2}\right\}$ of an equilateral triangle. Then $\mathcal{S G}$ is the unique nonempty compact set satisfying

$$
\begin{equation*}
\mathcal{S G}=\bigcup_{i=0}^{2} F_{i}(\mathcal{S G}) \tag{2.1}
\end{equation*}
$$

We refer to the sets $F_{i}(\mathcal{S G})$ as cells of level one, and by iterating (2.1) we obtain the splitting of $\mathcal{S G}$ into cells of higher level. For a word $w=\left(w_{1}, w_{2}, \cdots, w_{m}\right)$ of length $m$, the set $F_{w}(\mathcal{S G})=F_{w_{1}} \circ F_{w_{2}} \circ \cdots \circ F_{w_{m}}(\mathcal{S G})$ with $w_{i} \in\{0,1,2\}$, is called a $m$-cell. The fractal $\mathcal{S G}$ can be realized as the limit of a sequence of graphs $\Gamma_{0}, \Gamma_{1}, \cdots$ with vertices $V_{0} \subseteq V_{1} \subseteq \cdots$. The initial graph $\Gamma_{0}$ is just the complete graph on $V_{0}=\left\{q_{0}, q_{1}, q_{2}\right\}$, which is considered the boundary of $\mathcal{S G}$. See Fig. 2.1. Note that $\mathcal{S G}$ is connected, but just barely: there is a dense set of points $\mathcal{J}$, called junction points, defined by the condition that $x \in \mathcal{J}$ if and only if $U \backslash\{x\}$ is disconnected for all sufficiently small neighborhoods $U$ of $x$. It is easy to see that $\mathcal{J}$ consists of all images of $\left\{q_{0}, q_{1}, q_{2}\right\}$ under iterates of the ifs. The vertices $\left\{q_{0}, q_{1}, q_{2}\right\}$ are not junction points. All other points in $\mathcal{S G}$ will be called generic points. In the $\mathcal{S G}$ case, $\mathcal{J}=V_{*} \backslash V_{0}$, where $V_{*}=\bigcup_{m} V_{m}$. However, it is not true for general p.c.f. self-similar sets. In all that follows, we assume that $\mathcal{S G}$ is equipped with


Fig. 2.1. The first 3 graphs, $\Gamma_{0}, \Gamma_{1}, \Gamma_{2}$ in the approximation to the Sierpinski gasket.
the probability measure $\mu$ that assigns the measure $3^{-m}$ to each $m$-cell.
We define the unrenormalized energy of a function $u$ on $\Gamma_{m}$ by

$$
E_{m}(u)=\sum_{x \sim_{m} y}(u(x)-u(y))^{2}
$$

The energy renormalization factor is $r=\frac{3}{5}$, so the renormalized graph energy on $\Gamma_{m}$ is

$$
\mathcal{E}_{m}(u)=r^{-m} E_{m}(u),
$$

and we can define the fractal energy $\mathcal{E}(u)=\lim _{m \rightarrow \infty} \mathcal{E}_{m}(u)$. We define dom $\mathcal{E}$ as the space of continuous functions with finite energy. Then $\mathcal{E}$ extends by polarization to a bilinear form $\mathcal{E}(u, v)$ which serves as an inner product in this space.

The standard Laplacian may then be defined using the weak formulation: $u \in \operatorname{dom} \Delta$ with $\Delta u=f$ if $f$ is continuous, $u \in \operatorname{dom} \mathcal{E}$, and

$$
\mathcal{E}(u, v)=-\int f v d \mu
$$

for all $v \in \operatorname{dom}_{0} \mathcal{E}$, where $\operatorname{dom}_{0} \mathcal{E}=\left\{v \in \mathcal{E}:\left.v\right|_{V_{0}}=0\right\}$. There is also a pointwise formula (which is proven to be equivalent in [6]) which, for points in $V_{*} \backslash V_{0}$ computes

$$
\Delta u(x)=\frac{3}{2} \lim _{m \rightarrow \infty} 5^{m} \Delta_{m} u(x)
$$

where $\Delta_{m}$ is a discrete Laplacian associated to the graph $\Gamma_{m}$, defined by

$$
\Delta_{m} u(x)=\sum_{y \sim_{m} x}(u(y)-u(x))
$$

for $x$ not on the boundary.
It is not necessary to invoke the measure to define harmonic functions, although it is true that these are just the solutions of $\Delta h=0$. The more direct definition is that

$$
h(x)=\frac{1}{4} \sum_{y \sim_{m} x} h(y)
$$

for every nonboundary point and every $m$. This can be viewed as a mean value property of $h$ at the junction points. The space of harmonic functions is 3-dimensional and the values at the 3 boundary points may be freely assigned. Moreover, there is a simple efficient algorithm, the " $\frac{1}{5}-\frac{2}{5}$ rule", for computing the values of a harmonic function exactly at all vertex points in terms of the boundary values. The harmonic functions satisfy the maximum principle, i.e., the maximum and minimum are attained on the boundary and only on the boundary if the function is not constant. We call a continuous function $h$ a piecewise harmonic spline of level $m$ if $h \circ F_{w}$ is harmonic for all $|w|=m$.

The Laplacian satisfies the scaling property

$$
\Delta\left(u \circ F_{i}\right)=\frac{1}{5}(\Delta u) \circ F_{i}
$$

and by iteration

$$
\Delta\left(u \circ F_{w}\right)=\frac{1}{5^{m}}(\Delta u) \circ F_{w}
$$

for $F_{w}=F_{w_{1}} \circ F_{w_{2}} \circ \cdots \circ F_{w_{m}}$.
Although there is no satisfactory analogue of gradient, there is normal derivative $\partial_{n} u\left(q_{i}\right)$ defined at boundary points by

$$
\partial_{n} u\left(q_{i}\right)=\lim _{m \rightarrow \infty} \sum_{y \sim \sim_{m} q_{i}} r^{-m}\left(u\left(q_{i}\right)-u(y)\right),
$$

the limit existing for all $u \in \operatorname{dom} \Delta$. The definition may be localized to boundary points of cells. For each point $x \in V_{m} \backslash V_{0}$, there are two cells containing $x$ as a boundary point, hence two normal derivatives at $x$. For $u \in \operatorname{dom} \Delta$, the normal derivatives at $x$ satisfy the matching condition that their sum is zero. The matching conditions allow us to glue together local solutions to $\Delta u=f$.

As is shown in [2, 3], the Dirichlet problem for the Laplacian can be solved by integrating against an explicitly given Green's function. Recall that the Green's function $G(x, y)$ is a uniform limit of $G_{M}(x, y)$ as $M$ goes to the infinity, with $G_{M}$ defined by

$$
G_{M}(x, y)=\sum_{m=0}^{M} \sum_{z, z^{\prime} \in V_{m+1} \backslash V_{m}} g\left(z, z^{\prime}\right) \psi_{z}^{(m+1)}(x) \psi_{z^{\prime}}^{(m+1)}(y)
$$

and
$\left\{\begin{array}{l}g(z, z)=\frac{9}{50} r^{m} \text { for } z \in V_{m+1} \backslash V_{m}, \\ g\left(z, z^{\prime}\right)=\frac{3}{50} r^{m} \text { for } z, z^{\prime} \in V_{m+1} \backslash V_{m} \text { with } z, z^{\prime} \in F_{w}(\mathcal{S G}) \text { for }|w|=m, \text { and } z \neq z^{\prime},\end{array}\right.$ where $\psi_{z}^{m}(x)$ denotes a piecewise harmonic spline of level $m$ satisfying $\psi_{z}^{(m)}(x)=\delta_{z x}$ for $x \in V_{m}$.

## 3 Mean value property of harmonic functions on $\mathcal{S G}$

Lemma 3.1. (a) Let $C$ be any cell with boundary points $p_{0}, p_{1}, p_{2}$, and $h$ any harmonic function. Then

$$
\frac{1}{\mu(C)} \int_{C} h d \mu=\frac{1}{3}\left(h\left(p_{0}\right)+h\left(p_{1}\right)+h\left(p_{2}\right)\right) .
$$

(b) Let $p$ be any junction point, and $C_{1}, C_{2}$ the two m-cells containing $p$. Then

$$
\frac{1}{\mu\left(C_{1} \cup C_{2}\right)} \int_{C_{1} \cup C_{2}} h d \mu=h(p)
$$

Proof. (a) follows by symmetry. (b) follows by combining (a) for $C=C_{1}$ and $C=C_{2}$ with the mean value property of $h$ at $p$.

Note that (b) gives a trivial solution to the problem of finding mean value neighborhoods for junction points.


Fig. 3.1. $C_{w}$ and its three neighboring cells.

Given a point $x$ in $\mathcal{S G} \backslash V_{0}$, consider any cell $F_{w}(\mathcal{S G})=C_{w}$ containing the point $x$, with boundary points $F_{w} q_{i}=p_{i}$. Choose the cell $C_{w}$ small enough, such that it does not intersect $V_{0}$. Then it must have three neighboring cells $C_{0}, C_{1}$ and $C_{2}$ of the same level with $C_{i}$ intersecting $C_{w}$ at $p_{i}$. Denote by $D_{w}$ the union of $C_{w}$ and its three neighbors. See Fig. 3.1. In this section, we will describe a method to find a subset $B$ of $D_{w}$, containing $C_{w}$, such that for any harmonic function $h$, the mean value of $h$ over $B$ is equal to it's value at $x$, i.e., $M_{B}(h)=h(x)$ where $M_{B}(h)$ is defined by

$$
M_{B}(h)=\frac{1}{\mu(B)} \int_{B} h d \mu .
$$

Then we call the set $B$ a $k$ level mean value neighborhood of $x$ associated to $C_{w}$ where $k$ is the length of $w$. Let $k_{0}$ be the smallest value of $k$ such that there exists a $k$ level cell $C_{w}$ containing $x$ but not intersecting $V_{0}$. ( $k_{0}$ depends on the location of $x$ in $\mathcal{S G}$.) Then we can find a sequence of words $w^{(k)}$ of length $k\left(k \geq k_{0}\right)$ and a sequence of mean value neighborhoods $B_{k}(x)$ associated to $C_{w^{(k)}}$. Obviously, $\left\{B_{k}(x)\right\}_{k \geq k_{0}}$ will form a system of neighborhoods of the point $x$ satisfying $\bigcap_{k \geq k_{0}} B_{k}(x)=\{x\}$.

Let $h$ be a harmonic function on $\mathcal{S G}$. The harmonic extension algorithm implies that there exist coefficients $\left\{a_{i}(x)\right\}$ depending only on the relative position of $x$ and $C_{w}$ such that

$$
h(x)=\sum_{i} a_{i}(x) h\left(p_{i}\right) .
$$

Moreover, since constants are harmonic we must have

$$
\sum_{i} a_{i}(x)=1
$$

and by the maximum principle all $a_{i}(x) \geq 0$. Let $W$ denote the triangle in $\mathbf{R}^{3}$ with boundary points $(1,0,0),(0,1,0)$ and $(0,0,1)$ and $\pi_{W}$ the plane containing $W$. See Fig. 3.2. So $\left\{a_{i}(x)\right\} \in W$ for any $x \in C_{w}$. Of course, not every point in $W$ occurs in this way.


Fig. 3.2. The triangle $W$ with 3 boundary points $P_{0}, P_{1}, P_{2}$.

On the other hand, given a set $B$ such that $C_{w} \subset B \subset D_{w}$, by linearity we have

$$
M_{B}(h)=\sum_{i} a_{i} h\left(p_{i}\right)
$$

for some coefficients ( $a_{0}, a_{1}, a_{2}$ ) depending only on the relative geometry of $B$ and $C_{w}$. Again we must have $\sum a_{i}=1$ by considering $h \equiv 1$. So $\left(a_{0}, a_{1}, a_{2}\right) \in \pi_{W}$. (Later we will show that $\left(a_{0}, a_{1}, a_{2}\right)$ does not have to belong to $W$ for some sets $B$.) Thus we have a map, denoted by $\mathcal{T}$ from the collection of $B$ 's to $\pi_{W}$. If we can show that the image of the map $\mathcal{T}$ covers the triangle $W$ for some reasonable class of sets $B$, then we can get a set $B$ over which the mean value property holds for all harmonic functions. Moreover, if we can prove $\mathcal{T}$ is one-to-one, then we get a mean value neighborhood $B$ of $x$ associated to $C_{w}$, that is unique within the collection of sets we are considering.

This is the basic idea of our method. Hence the remaining task in this section is to find a suitable class of sets $B$ on which $\mathcal{T}$ is one-to-one and such that the image of the map $\mathcal{T}$ covers the triangle $W$. Comparing with the usual mean value neighborhoods, (they are just balls in Euclidean case) it is reasonable to require $B$ to be as simple as possible. They should be connected, possess some symmetry properties, depend only on
the relative geometry of $x$ and $C_{w}$, and be independent of the level of $C_{w}$ and the location of $C_{w}$.

Consider a set $B$ containing $C_{w}$ and contained in $D_{w}$. $B$ must be made up of four parts, i.e.,

$$
B=C_{w} \cup E_{0} \cup E_{1} \cup E_{2}
$$

where $E_{i}=B \cap C_{i}$. Suppose the height of $C_{w}$ is $\rho$. Due to the above consideration, we restrict each $E_{i}$ to be a triangle obtained by cutting $C_{i}$ symmetrically with a line at any height $c_{i} \rho$ below the top vertex $p_{i}$ with $0 \leq c_{i} \leq 1$. Write the set $B=B\left(c_{0}, c_{1}, c_{2}\right)$. See Fig. 3.3. for a sketch of $B$. For example, $B(0,0,0)=C_{w}$ and $B(1,1,1)=D_{w}$. Denote by

$$
\mathcal{B}=\left\{B\left(c_{0}, c_{1}, c_{2}\right): 0 \leq c_{i} \leq 1\right\}
$$

the family of all such sets.


Fig. 3.3. The relative geometry of $B$ and $C_{w}$.

Then on $\mathcal{B}$ the map $\mathcal{T}$ described above can be viewed as a nonlinear vector valued function from $[0,1]^{3}$ to $\pi_{W}$. For simplicity, we may write $\mathcal{T}\left(c_{0}, c_{1}, c_{2}\right)=\left(a_{0}, a_{1}, a_{2}\right)$ for each set $B\left(c_{0}, c_{1}, c_{2}\right)$. The following lemma shows that the value $\mathcal{T}\left(c_{0}, c_{1}, c_{2}\right)$ is independent of the particular choice of $C_{w}$, which benefits from the symmetric properties of the set $B\left(c_{0}, c_{1}, c_{2}\right)$.

Lemma 3.2. $\mathcal{T}\left(c_{0}, c_{1}, c_{2}\right)$ is independent of the particular choice of $C_{w}$.
Proof. Let $h$ be a harmonic function. First we consider the integral $\int_{E_{i}} h d \mu$. Denote by $\left\{s_{i}, t_{i}, p_{i}\right\}$ the boundary points of $C_{i}$. By linearity, $\frac{1}{\mu\left(C_{w}\right)} \int_{E_{i}} h d \mu$ can be expressed as a non-negative linear combination of $\left\{h\left(s_{i}\right), h\left(t_{i}\right), h\left(p_{i}\right)\right\}$, which by symmetry must have the form

$$
\begin{equation*}
\int_{E_{i}} h d \mu=\left(m_{i} h\left(p_{i}\right)+n_{i}\left(h\left(s_{i}\right)+h\left(t_{i}\right)\right)\right) \mu\left(C_{w}\right) \tag{3.1}
\end{equation*}
$$

for some appropriate non-negative coefficients $m_{i}, n_{i}$. Notice that in (3.1), the coefficients $m_{i}, n_{i}$ are independent of the location of $C_{i}$ in $\mathcal{S G}$. Actually, they only depend on the relative position of $E_{i}$ in $C_{i}$, i.e., $m_{i}, n_{i}$ depend only on $c_{i}$. Using the mean value property at $p_{i}$, namely

$$
4 h\left(p_{i}\right)=h\left(p_{i-1}\right)+h\left(p_{i+1}\right)+h\left(s_{i}\right)+h\left(t_{i}\right),
$$

we obtain

$$
\begin{aligned}
\int_{E_{i}} h d \mu & =\left(m_{i} h\left(p_{i}\right)+n_{i}\left(4 h\left(p_{i}\right)-h\left(p_{i-1}\right)-h\left(p_{i+1}\right)\right)\right) \mu\left(C_{w}\right) \\
& =\left(\left(4 n_{i}+m_{i}\right) h\left(p_{i}\right)-n_{i}\left(h\left(p_{i-1}\right)+h\left(p_{i+1}\right)\right)\right) \mu\left(C_{w}\right)
\end{aligned}
$$

Notice that the ratio of $\mu\left(E_{i}\right)$ to $\mu\left(C_{i}\right)$ also depends only on $c_{i}$. Combined with Lemma 3.1(a), we see that $\left(a_{0}, a_{1}, a_{2}\right)=\mathcal{T}\left(c_{0}, c_{1}, c_{2}\right)$ is independent of the particular choice of $C_{w}$, depending only on $\left(c_{0}, c_{1}, c_{2}\right)$.

We will show the image of the map $\mathcal{T}$ covers the triangle $W$. More precisely, $\mathcal{T}\left(c_{0}, c_{1}, c_{2}\right)$ will fill out a set $\widetilde{W}$ which is a bit larger than $W$. Denote by $P_{0}=(1,0,0), P_{1}=(0,1,0)$ and $P_{2}=(0,0,1)$ the three boundary points of the triangle $W$ in $\mathbf{R}^{3}$ and by $O$ the center point of $W$.


Fig. 3.4. a $1 / 6$ region of $\widetilde{W}$ surrounded by $\overline{O Q_{0}}, \overline{O P_{2}}$ and $\widehat{P_{2} Q_{0}}$.

Lemma 3.3. $\mathcal{T}(0,0,1)=P_{2}$ and $\mathcal{T}(0,1,1)=Q_{0}$ where $Q_{0}=\left\{-\frac{1}{9}, \frac{5}{9}, \frac{5}{9}\right\}$ is a point in $\pi_{W}$ located outside of $W$.

Proof. This follows by a direct computation. We omit it.
Lemma 3.4. $\mathcal{T}(\{(0, c, 1): 0 \leq c \leq 1\})$ is a continuous curve lying outside of $W$, joining $P_{2}$ and $Q_{0}$. (See Fig. 3.4.)

Proof. From Lemma 3.3, by varying c continuously between 0 and 1 we trace a continuous curve $\widehat{P_{2} Q_{0}}$ joining $P_{2}$ and $Q_{0}$. So we only need to prove the curve $\widehat{P_{2} Q_{0}}$ lies
outside of $W$. To prove this, we consider the set $B=B(0, c, 1)$ for $0 \leq c \leq 1$. In this case

$$
B=C_{w} \cup E_{1} \cup C_{2}
$$

Given a harmonic function $h$, by the proof of Lemma 3.2, we have

$$
\int_{E_{1}} h d \mu=\left(\left(4 n_{1}+m_{1}\right) h\left(p_{1}\right)-n_{1}\left(h\left(p_{0}\right)+h\left(p_{2}\right)\right)\right) \mu\left(C_{w}\right),
$$

for some appropriate non-negative coefficients $m_{1}, n_{1}$ depending only on $c$.
On the other hand, we have

$$
\int_{C_{w} \cup C_{2}} h d \mu=2 h\left(p_{2}\right) \mu\left(C_{w}\right)
$$

by Lemma 3.1(b).
Hence

$$
\begin{aligned}
\int_{B} h d \mu & =\int_{E_{1}} h d \mu+\int_{C_{w} \cup C_{2}} h d \mu \\
& =\left(-n_{1} h\left(p_{0}\right)+\left(4 n_{1}+m_{1}\right) h\left(p_{1}\right)+\left(2-n_{1}\right) h\left(p_{2}\right)\right) \mu\left(C_{w}\right) .
\end{aligned}
$$

The coefficient of $h\left(p_{0}\right)$ is always less than 0 . Moreover, it equals to 0 if and only if $E_{1}=\emptyset(\mathrm{c}=0)$. Hence $\mathcal{T}(0, c, 1)$ will always lie on the outside of the triangle $W$ as $c$ varies between 0 and 1 .

Now we come to the main result of this section.
Theorem 3.1. The map $\mathcal{T}$ from $\mathcal{B}$ to $\pi_{W}$ fills out a region $\widetilde{W}$ which contains the triangle $W$.

Proof. We only need to prove the map $\mathcal{T}$ from $\mathcal{B}$ to $\pi_{W}$ fills out a $1 / 6$ region surrounded by the line segments $\overline{O Q_{0}}, \overline{O P_{2}}$ and the curve $\widehat{P_{2} Q_{0}}$ as shown in Fig. 3.4. Then we will get the desired result by exploiting the symmetry.

Consider a subfamily $\mathcal{B}_{1}=\{B(0,0, c): 0 \leq c \leq 1\}$ of $\mathcal{B}$. If we restrict the map $\mathcal{T}$ to $\mathcal{B}_{1}$, by varying $c$ continuously between 0 and 1 we trace a curve (it is a line segment, which follows from the symmetry of $E_{2}$ ) in $W$ joining the center $O$ and the vertex point $P_{2}$.

Consider another subfamily $\mathcal{B}_{2}=\{B(0, c, c): 0 \leq c \leq 1\}$ of $\mathcal{B}$. If we restrict the map $\mathcal{T}$ to $\mathcal{B}_{2}$, by varying $c$ continuously between 0 and 1 we trace a curve (it is also a line segment, which follows from the symmetric effect of $E_{1}$ and $E_{2}$ ) in $W$ joining the center $O$ and the point $Q_{0}$ across the boundary line $\overline{P_{1} P_{2}}$ with $Q_{0}$ located outside of $W$, where $Q_{0}$ is the point defined in Lemma 3.3.

Fix a number $0 \leq y \leq 1$. Consider a subfamily $\mathcal{C}_{y}=\{B(0, c, y): 0 \leq c \leq y\}$ of $\mathcal{B}$. If we restrict the map $\mathcal{T}$ to $\mathcal{C}_{y}$, by varying $c$ continuously between 0 and $y$ we trace a curve $\Gamma_{y}$ joining the two points $\mathcal{T}(0,0, y)$ and $\mathcal{T}(0, y, y)$. The first endpoint $\mathcal{T}(0,0, y)$ lies on the line segment $\overline{O P_{2}}$ and the second endpoint $\mathcal{T}(0, y, y)$ lies on the line segment $\overline{O Q_{0}}$. (See Fig. 3.4. for $\Gamma_{y}$.) When $y=0$, the curve $\Gamma_{0}$ draws back to the single center point $O$. When $y=1$, by Lemma 3.4, the curve $\Gamma_{1}$ is a continuous curve located outside of the triangle $W$. Moreover, $P_{2}$ is the only common points of $\Gamma_{1}$ and $W$. Hence if we vary $y$ continuously between 0 and 1 , we can fill out the $1 / 6$ region surrounded by the line segments $\overline{O Q_{0}}, \overline{O P_{2}}$ and the curve $\widehat{P_{2} Q_{0}}$.

Remark. In the proof of the above theorem, we actually only consider those sets $B$ in $\mathcal{B}$ which are contained in the union of $C_{w}$ and subsets of only two neighbors. See Fig. 3.5. Of course, the map $\mathcal{T}$ restricted to this subfamily is one-to-one, which can be easily seen from the proof. Hence instead of $\mathcal{B}$, the map $\mathcal{T}$ is one-to-one from $\mathcal{B}^{*}$ onto $\widetilde{W}$, where

$$
\mathcal{B}^{*}=\left\{B\left(0, c_{1}, c_{2}\right): 0 \leq c_{i} \leq 1\right\} \cup\left\{B\left(c_{0}, 0, c_{2}\right): 0 \leq c_{i} \leq 1\right\} \cup\left\{B\left(c_{0}, c_{1}, 0\right): 0 \leq c_{i} \leq 1\right\}
$$



Fig. 3.5. The 3 shapes of $B \in \mathcal{B}^{*}$ associated to $C_{w}$ shown in Fig. 3.1.

Based on the discussion in the beginning of this section, we then have
Theorem 3.2. For each point $x \in \mathcal{S G} \backslash V_{0}$, there exists a system of mean value neighborhoods $B_{k}(x)$ with $\bigcap_{k} B_{k}(x)=\{x\}$.

## 4 Mean value property of general functions on $\mathcal{S G}$

In this section, we extend the mean value property to more general functions on $\mathcal{S G}$. Given a point $x$ in $\mathcal{S G} \backslash V_{0}$ and a cell $C_{w}$ containing $x$, for each mean value neighborhood
$B$ of $x$ associated to $C_{w}$, we assign a constant $c_{B}$ to $B$. We want

$$
M_{B}(u)-u(x) \approx c_{B} \Delta u(x)
$$

for $u$ in $\operatorname{dom} \Delta$. More precisely, let $\left\{B_{k}(x)\right\}_{k \geq k_{0}}$ be the system of mean value neighborhoods of the point $x$; we want

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \frac{1}{c_{B_{k}(x)}}\left(M_{B_{k}(x)}-u(x)\right)=\Delta u(x) \tag{4.1}
\end{equation*}
$$

for appropriate functions in the domain of $\Delta$, which is the desired fractal analog of (1.1).
For this purpose, let $v$ be a function on $\mathcal{S G}$ satisfying $\Delta v \equiv 1$. For each point $x$ in $\mathcal{S G} \backslash V_{0}$, and each mean value neighborhood $B$ of $x$, define $c_{B}$ by

$$
c_{B}=M_{B}(v)-v(x) .
$$

Note that the result is independent of which $v$, because any two such functions differ by a harmonic function and the equality $M_{B}(h)-h(x)=0$ always holds for any harmonic function $h$. So we can choose

$$
v(x)=-\int G(x, y) d \mu(y)
$$

which vanishes on the boundary of $\mathcal{S G}$. Here $G$ is the Green's function.
We will prove that $c_{B}$ depends only on the relative geometry of $B$ and $C_{w}$ and the size of $C_{w}$, not on the location of $x$ or $C_{w}$ in $\mathcal{S G}$. More precisely, we will prove:

Theorem 4.1. Let $x \in \mathcal{S G} \backslash V_{0}$ and $B$ be a $k$ level mean value neighborhood of $x$. Then

$$
c_{0} \frac{1}{5^{k}} \leq c_{B} \leq c_{1} \frac{1}{5^{k}}
$$

for some constant $c_{0}, c_{1}$ which are independent of $x$.
On the other hand, given a point $x$ and $C_{w}$ a $k$ level neighborhood of $x$, for any $u \in \operatorname{dom} \Delta^{2}$, we write

$$
u=h^{(k)}+(\Delta u(x)) v+R^{(k)}
$$

on $C_{w}$, where $h^{(k)}$ is a harmonic function defined by

$$
h^{(k)}+\left.(\Delta u(x)) v\right|_{\partial C_{w}}=\left.u\right|_{\partial C_{w}} .
$$

It is not hard to prove the following estimate:
Lemma 4.1. The remainder $R^{(k)}$ satisfies

$$
R^{(k)}=O\left(\left(\frac{3}{5} \cdot \frac{1}{5}\right)^{k}\right)
$$

on $C_{w}$ (hence also on $B_{k}(x)$ ).
This looks like a Taylor expansion remainder estimate of $u$ at $x$. See more details on this topic in 7 .

Proof of Lemma 4.1.
It is easy to check that $\Delta_{y} R^{(k)}(y)=\Delta_{y} u(y)-\Delta_{y} u(x)$ and $R^{(k)}(y)$ vanishes on the boundary of $C_{w}$. Hence $R^{(k)}$ is given by the integral of $\Delta_{y} u(y)-\Delta_{y} u(x)$ on $C_{w}$ against a scaled Green's function. Since the scaling factor is $\left(\frac{1}{5}\right)^{k}$, and

$$
\left|\Delta_{y} u(y)-\Delta_{y} u(x)\right| \leq c(3 / 5)^{k}
$$

( $\Delta u$ satisfies a Holder condition with $\gamma=\frac{3}{5}$ ), we get the desired result.
A more general version of Lemma 4.1 is the following:
Lemma 4.1.' Let $u \in$ dom $\Delta$ with $g=\Delta u$ satisfying the following Holder condition

$$
|g(y)-g(x)| \leq c \gamma^{k}, \quad(0<\gamma<1)
$$

for all $y \in C_{w}$. Then the remainder satisfies

$$
R^{(k)}=O\left(\left(\frac{\gamma}{5}\right)^{k}\right)
$$

on $C_{w}$ (hence also on $B_{k}(x)$ ).
Using the Taylor expansion of $u$ at $x$ and Theorem 4.1, we have

$$
\begin{aligned}
& \frac{1}{c_{B_{k}(x)}}\left(M_{B_{k}(x)}(u)-u(x)\right)-\Delta u(x) \\
= & \frac{1}{c_{B_{k}(x)}}\left(M_{B_{k}(x)}\left(R^{(k)}\right)-R^{(k)}(x)\right) \\
= & \frac{1}{c_{B_{k}(x)}} O\left(\left(\frac{\gamma}{5}\right)^{k}\right)=O\left(\gamma^{k}\right) .
\end{aligned}
$$

Hence we get (4.1) by letting $k$ go to infinity.
Thus Theorem 4.1 implies the following result:
Theorem 4.2. Let $u \in \operatorname{dom} \Delta$ with $g=\Delta u$ satisfying the Holder condition $\mid g(y)-$ $g(x) \mid \leq c \gamma^{k}$ for some $\gamma$ with $0<\gamma<1$, for all $x, y$ belonging to the same $k$ level cell. Then

$$
\lim _{k \rightarrow \infty} \frac{1}{c_{B_{k}(x)}}\left(M_{B_{k}(x)}(u)-u(x)\right)=\Delta u(x)
$$

To prove Theorem 4.1, we need the explicit expression for the function $v$. Recall from Section 2 that $v(x)$ is the uniform limit of $v_{M}(x)$ for

$$
v_{M}(x)=-\int G_{M}(x, y) d \mu(y)
$$

Interchanging the integral and summation,

$$
v_{M}(x)=-\sum_{m=0}^{M} \sum_{z, z^{\prime} \in V_{m+1} \backslash V_{m}} g\left(z, z^{\prime}\right) \int \psi_{z^{\prime}}^{(m+1)}(y) d \mu(y) \psi_{z}^{(m+1)}(x)
$$

Since an easy computation shows that for each $z \in V_{m+1} \backslash V_{m}, \int \psi_{z}^{(m+1)}(y) d \mu(y)=\frac{2}{3^{m+2}}$, we find

$$
v_{M}(x)=-\frac{2}{9} \sum_{m=0}^{M} \frac{1}{3^{m}} \sum_{z, z^{\prime} \in V_{m+1} \backslash V_{m}} g\left(z, z^{\prime}\right) \psi_{z}^{(m+1)}(x)
$$

Substituting the exact value of $g\left(z, z^{\prime}\right)$ into it, we get

$$
\begin{aligned}
v_{M}(x) & =-\frac{2}{9} \sum_{m=0}^{M} \frac{1}{3^{m}} \sum_{z \in V_{m+1} \backslash V_{m}}\left(\frac{9}{50} r^{m}+2 \frac{3}{50} r^{m}\right) \psi_{z}^{(m+1)}(x) \\
& =-\frac{1}{15} \sum_{m=0}^{M} \frac{1}{5^{m}} \phi_{m}(x)
\end{aligned}
$$

for

$$
\phi_{m}(x)=\sum_{z \in V_{m+1} \backslash V_{m}} \psi_{z}^{(m+1)}(x)
$$

Thus

$$
v(x)=-\frac{1}{15} \sum_{m=0}^{\infty} \frac{1}{5^{m}} \phi_{m}(x)
$$

Remark. The function $v$ is invariant under Dihedral-3 symmetry.
This is a direct corollary of the fact that each $\phi_{m}(x)$ is invariant under $D_{3}$ symmetry.
Due to the above remark, we may assume that $D_{w}$ associated to $C_{w}$ has a fixed shape as shown in Fig. 3.1 without loss of generality. We now show that $c_{B}$ depends only on the relative geometry of $B$ and $C_{w}$ and the size of $C_{w}$, not on the location of $x$ or $C_{w}$ in $\mathcal{S G}$.

Lemma 4.2. Let $x, x^{\prime}$ be two distinct points in $\mathcal{S G} \backslash V_{0}$. Let $C_{w}$ and $C_{w^{\prime}}$ be two $k$ and $k^{\prime}$ level neighboring cells of $x$ and $x^{\prime}$ respectively. Denote by $B$ and $B^{\prime}$ two mean value neighborhoods of $x$ and $x^{\prime}$ respectively. If $B$ and $B^{\prime}$ have the same shapes (the same relative locations associated to $C_{w}$ and $C_{w^{\prime}}$ respectively), then

$$
c_{B}=5^{k^{\prime}-k} c_{B^{\prime}} .
$$

In particular, if $B$ and $B^{\prime}$ have the same levels and same shapes, then $c_{B}=c_{B^{\prime}}$.

Proof. $D_{w}$ can be decomposed into a union of a $k$ level cell $D_{w}^{(1)}$ and a $(k-1)$ level cell $D_{w}^{(2)}$ as shown in Fig. 3.1. Denote by $q$ the junction point connecting $D_{w}^{(1)}$ and $D_{w}^{(2)}$. Similarly, $D_{w^{\prime}}$ can also be written as a union of a $k^{\prime}$ cell $D_{w^{\prime}}^{(1)}$ and a $\left(k^{\prime}-1\right)$ cell $D_{w^{\prime}}^{(2)}$ with a junction point $q^{\prime}$ connecting them.

Let $\tau$ be the linear function mapping $D_{w}$ onto $D_{w^{\prime}}$. Suppose $D_{w}^{(1)}=F_{\alpha}(\mathcal{S G})$ and $D_{w}^{(2)}=F_{\beta}(\mathcal{S G})$ where $\alpha$ and $\beta$ are the corresponding words of $D_{w}^{(1)}$ and $D_{w}^{(2)}$ respectively. Similarly, denote by $\alpha^{\prime}$ and $\beta^{\prime}$ the corresponding words of $D_{w^{\prime}}^{(1)}$ and $D_{w^{\prime}}^{(2)}$. Hence we can write $\tau$ as $\tau(z)=F_{\alpha^{\prime}} \circ F_{\alpha}^{-1}(z)$ if $z \in D_{w}^{(1)}$, and $\tau(z)=F_{\beta^{\prime}} \circ F_{\beta}^{-1}(z)$ if $z \in D_{w}^{(2)}$. In particular, $\tau(q)=q^{\prime}$ and $\tau(x)=x^{\prime}$.

Consider the function $\left(v \circ F_{\alpha}-5^{k^{\prime}-k} v \circ F_{\alpha^{\prime}}\right)$ defined on $\mathcal{S G}$. Noting that $|\alpha|=k$ and $\left|\alpha^{\prime}\right|=k^{\prime}$, we have

$$
\Delta\left(v \circ F_{\alpha}-5^{k^{\prime}-k} v \circ F_{\alpha^{\prime}}\right)=r^{|\alpha|} \mu_{\alpha} \Delta v \circ F_{\alpha}-5^{k^{\prime}-k} r^{\left|\alpha^{\prime}\right|} \mu_{\alpha^{\prime}} \Delta v \circ F_{\alpha^{\prime}}=0,
$$

which shows that the difference between $v \circ F_{\alpha}$ and $5^{k^{\prime}-k} v \circ F_{\alpha^{\prime}}$ is a harmonic function. Hence the difference between $v$ and $5^{k^{\prime}-k} v \circ \tau$ on $D_{w}^{(1)}$ is harmonic. A similar discussion will show that the difference between $v$ and $5^{k^{\prime}-k} v \circ \tau$ on $D_{w}^{(2)}$ is also harmonic. Since the matching condition on normal derivatives of $\left(v-5^{k^{\prime}-k} v \circ \tau\right)$ at $q$ holds obviously, we have proved that $\Delta\left(v-5^{k^{\prime}-k} v \circ \tau\right)=0$ on $D_{w}$, i.e., the function $\left(v-5^{k^{\prime}-k} v \circ \tau\right)$ is harmonic on $D_{w}$.

By the definition $c_{B}=M_{B}(v)-v(x)$ and $c_{B^{\prime}}=M_{B^{\prime}}(v)-v\left(x^{\prime}\right)$. Notice that for the second equality, by changing variables we can write $c_{B^{\prime}}=M_{B}(v \circ \tau)-v \circ \tau(x)$. Hence

$$
c_{B}-5^{k^{\prime}-k} c_{B^{\prime}}=M_{B}\left(v-5^{k^{\prime}-k} v \circ \tau\right)-\left(v-5^{k^{\prime}-k} v \circ \tau\right)(x)=0,
$$

since $\left(v-5^{k^{\prime}-k} v \circ \tau\right)$ is a harmonic function on $D_{w}$.
Proof of Theorem 4.1.

## Estimate of $c_{B}$ from above.

From Lemma 4.2, since $c_{B}$ depends only on the relative geometry of $B$ and $C_{w}$ and the size of $C_{w}$, but not on the location of $C_{w}$, we may assume that $D_{w}$ is contained in a $(k-2)$ level cell $C$ in $\mathcal{S G}$ without loss of generality.

By the definition of $c_{B}$, we may write

$$
c_{B}=M_{B}(v)-v(x)=\lim _{M \rightarrow \infty}\left(\frac{1}{\mu(B)} \int_{B} v_{M} d \mu-v_{M}(x)\right) .
$$

Substituting the exact formula of $v_{M}$ into it, we get

$$
c_{B}=-\frac{1}{15} \sum_{m=0}^{\infty} \frac{1}{5^{m}}\left(M_{B}\left(\phi_{m}\right)-\phi_{m}(x)\right),
$$

for

$$
\phi_{m}=\sum_{z \in V_{m+1} \backslash V_{m}} \psi_{z}^{(m+1)}
$$

Notice that each $\phi_{m}$ is a piecewise harmonic spline of level $m+1$. So when $m+1 \leq k-2$, $\phi_{m}$ is harmonic in the cell $C$, which yields that $M_{B}\left(\phi_{m}\right)-\phi_{m}(x)=0$. So the first $k-2$ terms in the infinite series of $v$ will contribute 0 to $c_{B}$. Hence

$$
c_{B}=-\frac{1}{15} \sum_{m=k-2}^{\infty} \frac{1}{5^{m}}\left(M_{B}\left(\phi_{m}\right)-\phi_{m}(x)\right) .
$$

It is easy to see that this implies

$$
\left|c_{B}\right| \leq \frac{1}{15} \sum_{m=k-2}^{\infty} \frac{1}{5^{m}} \frac{1}{\mu(B)} \int_{B}\left|\phi_{m}(y)-\phi_{m}(x)\right| d \mu(y)
$$

Then by the maximum principle, we finally get

$$
\left|c_{B}\right| \leq \frac{1}{15} \sum_{m=k-2}^{\infty} \frac{1}{5^{m}}=\frac{25}{12} \cdot \frac{1}{5^{k}}
$$

## Estimate of $c_{B}$ from below.



Fig. 4.1. a $1 / 3$ region of $C_{w}$.

Without loss of generality, we assume that $x$ is located in the $1 / 3$ region of $C_{w}$ as shown in Fig. 4.1, i.e., $x$ is contained in the triangle $T_{p_{1}, p_{2}, o}$, where $o$ is the geometric center of $C_{w}$. Then by the proof of Theorem 3.1, B is a subset of the union of $C_{w}$ and two of its neighbors $C_{1}$ and $C_{2}$. Hence we can write $B=C_{w} \cup E_{1} \cup E_{2}$, where $E_{i}=B \cap C_{i}$.

Claim 1. Let $\widetilde{B}=F_{0}(\mathcal{S G}) \cup \widetilde{E}_{1} \cup \widetilde{E}_{2}$, where $\widetilde{E}_{i}$ is a triangle obtained by cutting $F_{i}(\mathcal{S G})$ symmetrically with a line below the top vertex $F_{i} q_{0}$.(see Fig. 4.2.) If $\widetilde{B}$ and $B$ have the same shapes, then

$$
c_{B}=5^{1-k} c_{\widetilde{B}}
$$



Fig. 4.2. a sketch of $\widetilde{B}$.

This is a direct corollary of Lemma 4.2.
We only need to prove that $c_{\widetilde{B}}$ for $\widetilde{B}$ defined in Claim 1 has a positive lower bound. For simplicity of notation, in all that follows, we write $B$ instead of $\widetilde{B}$. In other words, we only need to consider $B$ whose associate cell $C_{w}$ is $F_{0}(\mathcal{S G})$. In this setting, $p_{i}=F_{0} q_{i}$, $C_{1}=F_{1}(\mathcal{S G})$ and $C_{2}=F_{2}(\mathcal{S G})$.

We write $v=-\frac{1}{15} \widetilde{v}$ where $\widetilde{v}$ is the non-negative function defined by

$$
\widetilde{v}(y)=\sum_{m=0}^{\infty} \frac{1}{5^{m}} \phi_{m}(y) .
$$

We have the following two claims on $\widetilde{v}$.
Claim 2. For each $x$ contained in the triangle $T_{p_{1}, p_{2}, o}, \widetilde{v}(x) \geq \frac{24}{25}$.
Proof. Denote by

$$
\widetilde{v}_{M}(y)=\sum_{m=0}^{M} \frac{1}{5^{m}} \phi_{m}(y)
$$

the partial sum of the first $M+1$ terms of $\widetilde{v}$. A direct computation shows that

$$
\widetilde{v}\left(F_{u} q_{0}\right)=\widetilde{v}_{2}\left(F_{u} q_{0}\right)=\frac{24}{25},
$$

and

$$
\widetilde{v}\left(F_{u} q_{1}\right)=\widetilde{v}_{2}\left(F_{u} q_{1}\right)=\widetilde{v}\left(F_{u} q_{2}\right)=\widetilde{v}_{2}\left(F_{u} q_{2}\right)=1,
$$

for $u=(0,1,1),(0,1,2),(0,2,1)$ and $(0,2,2)$. (In fact, $\widetilde{v}$ takes the constant value 1 along the line segment joining $p_{1}$ and $p_{2}$.) Notice that for each point $x$ in the triangle $T_{p_{1}, p_{2}, o}$, $x$ is contained in one of the four 3 level cells $F_{011}(\mathcal{S G}), F_{012}(\mathcal{S G}), F_{021}(\mathcal{S G})$ and $F_{022}(\mathcal{S G})$. Since $\widetilde{v}_{2}$ is harmonic in each such cell, by using the maximal principle, we get that

$$
\widetilde{v}_{2}(x) \geq \frac{24}{25}
$$

Hence $\widetilde{v}(x) \geq \frac{24}{25}$ since each term in the infinite series of $\widetilde{v}$ is non-negative.
Claim 3. $M_{B}(\widetilde{v}) \leq \frac{17}{18}$.
Proof. First of all, we prove that

$$
\int_{F_{0}(\mathcal{S G})} \widetilde{v}(y) d \mu(y)=\frac{5}{18} .
$$

We need to compute $\int_{F_{0}(\mathcal{S G})} \phi_{m}(y) d \mu(y)$ for each non-negative integer $m$. When $m=0$,

$$
\int_{F_{0}(\mathcal{S G})} \phi_{m}(y) d \mu(y)=\int_{F_{0}(\mathcal{S G})} \phi_{0}(y) d \mu(y)=2 \cdot \frac{1}{3} \cdot \frac{1}{3}=\frac{2}{9} .
$$

When $m \geq 1$, also an easy computation shows that

$$
\int_{F_{0}(\mathcal{S G})} \phi_{m}(y) d \mu(y)=3^{m} \cdot \frac{2}{3} \cdot \frac{1}{3^{m+1}}=\frac{2}{9} .
$$

Hence

$$
\int_{F_{0}(\mathcal{S G})} \widetilde{v}(y) d \mu(y)=\frac{2}{9} \sum_{m=0}^{\infty} \frac{1}{5^{m}}=\frac{5}{18} .
$$

By our assumption, the mean value neighborhood $B$ can be written as

$$
B=F_{0}(\mathcal{S G}) \cup E_{1} \cup E_{2},
$$

where $E_{i}=B \cap C_{i}$. Hence we have

$$
\begin{aligned}
M_{B}(\widetilde{v}) & =\frac{1}{\mu\left(F_{0}(\mathcal{S G})\right)+\mu\left(E_{1}\right)+\mu\left(E_{2}\right)}\left(\int_{F_{0}(\mathcal{S G})} \widetilde{v}(y) d \mu(y)+\int_{E_{1}} \widetilde{v}(y) d \mu(y)+\int_{E_{2}} \widetilde{v}(y) d \mu(y)\right) \\
& \leq \frac{1}{\mu\left(F_{0}(\mathcal{S G})\right)+\mu\left(E_{1}\right)+\mu\left(E_{2}\right)}\left(\int_{F_{0}(\mathcal{S G})} \widetilde{v}(y) d \mu(y)+\int_{E_{1}} 1 \cdot d \mu(y)+\int_{E_{2}} 1 \cdot d \mu(y)\right) \\
& =\frac{5 / 18+\mu\left(E_{1}\right)+\mu\left(E_{2}\right)}{1 / 3+\mu\left(E_{1}\right)+\mu\left(E_{2}\right)} .
\end{aligned}
$$

Since $0 \leq \mu\left(E_{1}\right)+\mu\left(E_{2}\right) \leq \frac{2}{3}$, an easy calculation shows that

$$
\frac{5 / 18+\mu\left(E_{1}\right)+\mu\left(E_{2}\right)}{1 / 3+\mu\left(E_{1}\right)+\mu\left(E_{2}\right)} \leq \frac{5 / 18+2 / 3}{1 / 3+2 / 3}=\frac{17}{18}
$$

Hence we always have

$$
M_{B}(\widetilde{v}) \leq \frac{17}{18}
$$

Now we turn to estimate $c_{B}$. Obviously,

$$
c_{B}=M_{B}(v)-v(x)=-\frac{1}{15}\left(M_{B}(\widetilde{v})-\widetilde{v}(x)\right)
$$

By Claim 2 and Claim 3, we notice that $M_{B}(\widetilde{v})-\widetilde{v}(x) \leq \frac{17}{18}-\frac{24}{25}=-\frac{7}{450}$. Hence

$$
c_{B} \geq \frac{1}{15} \cdot \frac{7}{450}>0
$$

## 5 p.c.f. fractals with Dihedral-3 symmetry

The results for $\mathcal{S G}$ should extend to other p.c.f. fractals which possess symmetric properties of both the geometric structure and the harmonic structure. We assume that a regular harmonic structure is given on a p.c.f. self-similar fractal $K$. The reader is referred to [3, 6] for exact definitions and any unexplained notations. We assume now that $\sharp V_{0}=3$ and all structures possess full $D_{3}$ symmetry. This means there exists a group $\mathcal{G}$ of homeomorphisms of $K$ isomorphic to $D_{3}$ that acts as permutations on $V_{0}$, and $\mathcal{G}$ preserves the self-similar and harmonic structures and the self-similar measure. We must have $r_{0}=r_{1}=r_{2}$ and $\mu_{0}=\mu_{1}=\mu_{2}$, but in general it is not necessary that all $r$ 's and all $\mu$ 's be the same.

We begin by assuming that the fractal $K$ is the invariant set of a finite iterated function system of contractive similarities. We denote these maps $\left\{F_{i}\right\}_{i=1, \cdots, N}$ with $N \geq 3$. We denote $V_{0}=\left\{q_{0}, q_{1}, q_{2}\right\}$ the set of boundary points.

Examples. (i) The Sierpinski gasket $\mathcal{S G}$. In this case all $r_{i}=3 / 5$ and all $\mu_{i}=1 / 3$.
(ii) The hexagasket, or fractal Star of David, can be generated by 6 maps with simultaneously rotate and contract by a factor of $1 / 3$ in the plane. Thus $V_{0}$ consists of 3 points of an equilateral triangle, and $V_{1}$ consists of the vertices of the Star of David, as shown in Fig. 5.1. Although the same geometric fractal can be constructed by using contractions which do not rotate, this gives rise to a different self-similar structure (in particular with $\sharp V_{0}=6$ ). Our choice of self-similar structure destroys the $D_{6}$ symmetry of the geometric fractal, but it has the advantage of easier computation. In this case, all $r_{i}=3 / 7$ and all $\mu_{i}=1 / 6$. Note that in this example there exist points in $V_{1}$ that are not junction points.
(iii) The level 3 Sierpinski gasket $\mathcal{S G}_{3}$, obtained by taking 6 contractions of ratio $1 / 3$ as shown in Fig. 5.2. Here we have all $r_{i}=7 / 15$ and $\mu_{i}=1 / 6$. Note that all seven

$\Gamma_{0}$

$\Gamma_{1}$

Fig. 5.1. The first 2 graphs, $\Gamma_{0}, \Gamma_{1}$ in the approximation to the hexagasket.
vertices in $V_{1} \backslash V_{0}$ are junction points, but the one in the middle intersects three 1-cells. In a similar manner we could define $\mathcal{S G}_{n}$ for any value of $n \geq 2$.


Fig. 5.2. The graph of the $V_{1}$ vertices of the level 3 Sierpinski gasket.

We prove that there are results analogous to Theorem 3.1, which yield the existence of mean value neighborhoods associated to $K$.

Given a point $x$ in $K \backslash V_{0}$, consider any cell $F_{w} K=C_{w}$ with boundary points $p_{0}, p_{1}, p_{2}$ containing the point $x$. Without losing of generality, we may require that the cell $C_{w}$ does not intersect $V_{0}$. For each $i$, denote by $C_{i, 1}, \cdots, C_{i, l_{i}}$ the neighboring cells of $C_{w}$ of the same size, intersecting $C_{w}$ at $p_{i}$, where $l_{i}$ is the number of such cells. It is possible that $l_{i}=0$ for some $i$ since $p_{i}$ may be a non-junction point. If this is true, the matching condition says that the normal derivative of any harmonic function $h$ must be zero at this point, which yields that the value of $h$ at this point is the mean value of the values of $h$ at the other two boundary points of $C_{w}$. In other words, the restriction of all global
harmonic functions in $C_{w}$ is two dimensional. Denote by $D_{w}$ the union of $C_{w}$ and all its neighboring cells, i.e.,

$$
D_{w}=C_{w} \cup \bigcup_{i, j} C_{i, j}
$$

Two cells $C_{w}$ and $C_{w^{\prime}}$ are said to have the same neighborhood type if they have the same relative geometry with respect to $D_{w}$ and $D_{w^{\prime}}$ respectively. It is obvious that there only exist finitely many distinct types. For example, for $\mathcal{S G}$, all cells have exactly only one neighborhood type. For $\mathcal{S G}_{3}$, the number of the finite types is 3 . For $\mathcal{S G}_{n}(n \geq 4)$, the number of the finite types becomes 4. For the hexagasket gasket, the number of the finite types is 2 .

Let $h$ be a harmonic function on $K$. Given a set $B$ containing $C_{w}$, define

$$
M_{B}(h)=\frac{1}{\mu(B)} \int_{B} h d \mu
$$

the mean value of $h$ over $B$. We are interested in an identity

$$
\begin{equation*}
M_{B}(h)=\sum_{i} a_{i} h\left(p_{i}\right) \tag{5.1}
\end{equation*}
$$

for some coefficients $\left(a_{0}, a_{1}, a_{2}\right)$ satisfying $\sum a_{i}=1$. Notice that this is true for $\mathcal{S G}$. In that setting, a harmonic function is uniquely determined by its values on the boundary of any given cell $C_{w}$ because the harmonic extension matrix associated with $C_{w}$ is invertible. However, in general case, the harmonic extension matrices may not be invertible. So we can not prove (5.1) for every set $B$ simply by linearity. However, it will suffice to show that the equality (5.1) holds for certain specified sets $B$.

Consider a set $B$ which is a subset of $D_{w}$, containing $C_{w}$. $B$ must be made up of four parts, i.e.,

$$
B=C_{w} \cup E_{0} \cup E_{1} \cup E_{2}
$$

where $E_{i}=B \cap C_{i}$ with $C_{i}=\cup_{j=1}^{l_{i}} C_{i, j}$. It is possible that $C_{i}$ may be empty since $p_{i}$ may be a nonjunction point. We can also subdivide each $E_{i}$ into $l_{i}$ small pieces, i.e., $E_{i}=\cup_{j} E_{i, j}$ for $E_{i, j}=E_{i} \cap C_{i, j}$. For each $i$, we require that $E_{i, 1}, \cdots, E_{i, l_{i}}$ be of the same size and shape. Moreover, in analogy with the $\mathcal{S G}$ case, we require that each $E_{i, j}$ to be a symmetric (under the reflection symmetry that fixes $p_{i}$ ) cutoff sub-triangle of $C_{i, j}$, containing $p_{i}$ as one of its vertex points. This means that there is a straight line $L_{i, j}$, symmetric under the reflection symmetry fixing $p_{i}$, cutting $C_{i, j}$ into two parts, and $E_{i, j}$ is the one containing $p_{i}$. For each $E_{i, j}$, define the distance between $p_{i}$ and the line $L_{i, j}$ the height of $E_{i, j}$. Of course, for each fixed $i, E_{i, 1}, \cdots, E_{i, l_{i}}$ have the same heights. We
call the common value the height of $E_{i}$. Suppose the height of every $C_{i, j}$ is $\rho$. (Of course, they all equal.) Then for each $i$, the height of $E_{i}$ is $c_{i} \rho$ where the coefficient $0 \leq c_{i} \leq 1$. Hence we can write the set $B=B\left(c_{0}, c_{1}, c_{2}\right)$. (If $p_{i}$ is a nonjunction point, then $c_{i}$ should always be 0.) For example, suppose that the boundary points of $C_{w}$ consist of junction points, then $B(0,0,0)=C_{w}$ and $B(1,1,1)=D_{w}$. Denote by

$$
\mathcal{B}=\left\{B\left(c_{0}, c_{1}, c_{2}\right): 0 \leq c_{i} \leq 1\right\}
$$

the family of all such sets. Then we can show that the formula (5.1) holds for each $B \in \mathcal{B}$.
Proposition 5.1. Let $B \in \mathcal{B}$, then for any harmonic function $h$, we have (5.1) for some coefficients $\left(a_{0}, a_{1}, a_{2}\right)$ independent of $h$. Moreover, $\sum_{i} a_{i}=1$.

Proof. Each $B \in \mathcal{B}$ can be written as $B=C_{w} \cup E_{0} \cup E_{1} \cup E_{2}$. Given a harmonic function $h$ on $K$, for fixed $i$, we first consider the integral $\int_{E_{i}} h d \mu$. Obviously,

$$
\int_{E_{i}} h d \mu=\sum_{j} \int_{E_{i, j}} h d \mu .
$$

For each $1 \leq j \leq l_{i}$, denote by $\left\{z_{i, j}, w_{i, j}, p_{i}\right\}$ the boundary points of $C_{i, j}$. Since each $E_{i, j}$ is contained in $C_{i, j}, \frac{1}{\mu\left(C_{w}\right)} \int_{E_{i, j}} h d \mu$ can be expressed as a linear combination of $h\left(p_{i}\right), h\left(z_{i, j}\right)$ and $h\left(w_{i, j}\right)$ with non-negative coefficients independent of the harmonic function $h$. Since the set $E_{i, j}$ is symmetric under the reflection symmetry fixing $p_{i}$, the two coefficients with respect to $h\left(z_{i, j}\right)$ and $h\left(w_{i, j}\right)$ must be equal. In other words, we can write

$$
\int_{E_{i, j}} h d \mu=\left(m_{i, j} h\left(p_{i}\right)+n_{i, j} h\left(z_{i, j}\right)+n_{i, j} h\left(w_{i, j}\right)\right) \mu\left(C_{w}\right)
$$

for $m_{i, j}, n_{i, j} \geq 0$. Moreover, since for each fixed $i, E_{i, j}$ are in the same relative position associated to $C_{i, j}$ for different $j$ 's, $\int_{E_{i, j}} h d \mu$ can be expressed as a linear combination of $h\left(p_{i}\right), h\left(z_{i, j}\right), h\left(w_{i, j}\right)$ with the same coefficients for different $j$ 's. Hence we can write

$$
\int_{E_{i}} h d \mu=\left(m_{i} h\left(p_{i}\right)+n_{i} \sum_{j}\left(h\left(z_{i, j}\right)+h\left(w_{i, j}\right)\right)\right) \mu\left(C_{w}\right),
$$

for suitable coefficients $m_{i}, n_{i} \geq 0$. The mean value property at the point $p_{i}$ says that

$$
\sum_{j}\left(h\left(z_{i, j}\right)+h\left(w_{i, j}\right)\right)=\left(2 l_{i}+2\right) h\left(p_{i}\right)-\left(h\left(p_{i-1}\right)+h\left(p_{i+1}\right)\right) .
$$

Combining the above two equalities, we get

$$
\int_{E_{i}} h d \mu=\left(\left(m_{i}+2 l_{i} n_{i}+2 n_{i}\right) h\left(p_{i}\right)-n_{i} h\left(p_{i-1}\right)-n_{i} h\left(p_{i+1}\right)\right) \mu\left(C_{w}\right) .
$$

On the other hand, by the linearities and symmetries of both the harmonic structure and the self-similar measure,

$$
\int_{C_{w}} h d \mu=\frac{\mu\left(C_{w}\right)}{3}\left(h\left(p_{0}\right)+h\left(p_{1}\right)+h\left(p_{2}\right)\right) .
$$

Since the ratio of $\mu\left(E_{i, j}\right)$ to $\mu\left(C_{w}\right)$ depends only on $c_{i}$, we have proved that $M_{B}(h)$ can be viewed as a linear combination of the values of $h$ on the boundary points of $C_{w}$, i.e.,

$$
M_{B}(h)=\sum_{i} a_{i} h\left(p_{i}\right),
$$

where the combination coefficients are independent of $h$. Moreover, we must have $\sum a_{i}=1$ by considering $h \equiv 1$.

Remark 1. This means that $M_{B}(h)$ is a weighted average of the values $h\left(p_{0}\right), h\left(p_{1}\right)$ and $h\left(p_{2}\right)$. Moreover, if one of the boundary points, for example $p_{2}$, is a nonjunction point, then by the fact that $h\left(p_{2}\right)=\frac{1}{2}\left(h\left(p_{0}\right)+h\left(p_{1}\right)\right)$, we have

$$
M_{B}(h)=a_{0} h\left(p_{0}\right)+a_{1} h\left(p_{1}\right)+\frac{1}{2} a_{2}\left(h\left(p_{0}\right)+h\left(p_{1}\right)\right)=\widetilde{a}_{0} h\left(p_{0}\right)+\widetilde{a}_{1} h\left(p_{1}\right)
$$

for $\widetilde{a}_{0}=a_{0}+\frac{1}{2} a_{2}$ and $\widetilde{a}_{1}=a_{1}+\frac{1}{2} a_{2}$. We also have $\widetilde{a}_{0}+\widetilde{a}_{1}=1$. Hence in this case, we can also view $M_{B}(h)$ as a weighted average of the values of $h\left(p_{0}\right)$ and $h\left(p_{1}\right)$.

Remark 2. The proof of Proposition 5.1 shows that $\left(a_{0}, a_{1}, a_{2}\right)$ depends only on the neighborhood type of $C_{w}$ and the relative position of $B$ associated to $C_{w}$, and does not depend on the particular choice of $C_{w}$. In other words, if we consider a cell $C_{w}$ with a given neighborhood type, then for each set $B \in \mathcal{B}$ with the expression $B=B\left(c_{0}, c_{1}, c_{2}\right)$, the coefficients $\left(a_{0}, a_{1}, a_{2}\right)$ of $B$ depend only on $\left(c_{0}, c_{1}, c_{2}\right)$.

To prove the desired result of this section, we classify the distinct neighborhood types into three cases according to the number of nonjunction points in the set of boundary points of $C_{w}$.

## Case 1. All boundary points of $C_{w}$ are junction points.

This case is similar to what we have described in the $\mathcal{S G}$ setting. Let $W$ denote the triangle in $\mathbf{R}^{3}$ with boundary points $P_{0}=(1,0,0), P_{1}=(0,1,0)$ and $P_{2}=(0,0,1)$ and $\pi_{W}$ the plane containing $W$. Notice that from Proposition 5.1, $\left(a_{0}, a_{1}, a_{2}\right) \in \pi_{W}$ for each $B$. We use $\mathcal{T}$ to denote the map from $\mathcal{B}$ to $\pi_{W}$. From Remark 2 of Proposition 5.1, the map $\mathcal{T}$ is unique determined by the neighborhood type of $C_{w}$. Let $\mathcal{B}^{*}$ be a subfamily contained in $\mathcal{B}$ defined by

$$
\mathcal{B}^{*}=\left\{B\left(0, c_{1}, c_{2}\right): 0 \leq c_{i} \leq 1\right\} \cup\left\{B\left(c_{0}, 0, c_{2}\right): 0 \leq c_{i} \leq 1\right\} \cup\left\{B\left(c_{0}, c_{1}, 0\right): 0 \leq c_{i} \leq 1\right\}
$$

i.e., those elements $B$ in $\mathcal{B}$ which have the decomposition form $B=C_{w} \cup E_{1} \cup E_{2}$ or $B=C_{w} \cup E_{0} \cup E_{2}$, or $B=C_{w} \cup E_{0} \cup E_{1}$. Then we have

Claim 1. The map $\mathcal{T}$ from $\mathcal{B}$ to $\pi_{W}$ fills out a region $\widetilde{W}$ which contains the triangle $W$. Moreover, $\mathcal{T}$ is one-to-one from $\mathcal{B}^{*}$ onto $\widetilde{W}$.

Proof. The proof is similar to the $\mathcal{S G}$ case. The only difference is the line segments $\overline{O Q_{0}}$ and $\overline{O P_{2}}$ described in the proof of Theorem 3.1 may become continuous curves $\widehat{O Q_{0}}$ and $\overline{O P_{2}}$ in the general setting.

Case 2. There is one nonjunction point (for example, $p_{2}$ ) among the boundary points of $C_{w}$.

In this case, there is no neighboring cell intersecting $C_{w}$ at the point $p_{2}$. Hence $E_{2}$ will always be empty. So $\mathcal{B}=\left\{B\left(c_{0}, c_{1}, 0\right): 0 \leq c_{i} \leq 1\right\}$ for this case.

As shown in Remark 1 of Proposition 5.1, for any harmonic function $h$ on $K, B \in \mathcal{B}$, $M_{B}(h)$ is a weighted average of $h\left(p_{0}\right)$ and $h\left(p_{1}\right)$, i.e.,

$$
M_{B}(h)=a_{0} h\left(p_{0}\right)+a_{1} h\left(p_{1}\right)
$$

with $a_{0}, a_{1}$ independent of $h$, satisfying $a_{0}+a_{1}=1$. Let $I$ denote the line segment in $\mathbf{R}^{2}$ with endpoints $P_{0}=(1,0), P_{1}=(0,1)$ and $\rho_{I}$ the line containing $I$. Notice that from Remark 1 of Proposition 5.1, $\left(a_{0}, a_{1}\right) \in \rho_{I}$ for each $B$. We still use $\mathcal{T}$ to denote the map from $\mathcal{B}$ to $\rho_{I}$. From Remark 2 of Proposition 5.1, the map $\mathcal{T}$ is unique determined by the neighborhood type of $C_{w}$. For simplicity, we may write $\mathcal{T}\left(c_{0}, c_{1}\right)=\left(a_{0}, a_{1}\right)$ for each set $B\left(c_{0}, c_{1}, 0\right)$. We will show the image of the map $\mathcal{T}$ covers the line segment $I$. Similar to Case 1 , let $\mathcal{B}^{*}$ be a subfamily contained in $\mathcal{B}$ defined by

$$
\mathcal{B}^{*}=\left\{B\left(c_{0}, 0,0\right): 0 \leq c_{0} \leq 1\right\} \cup\left\{B\left(0, c_{1}, 0\right): 0 \leq c_{1} \leq 1\right\},
$$

i.e., those elements $B$ in $\mathcal{B}$ which have the decomposition form $B=C_{w} \cup E_{0}$ or $B=$ $C_{w} \cup E_{1}$. Then we have

Claim 2. The map $\mathcal{T}$ from $\mathcal{B}$ to $\rho_{I}$ fills out the line segment $I$. Moreover, $\mathcal{T}$ is a one-to-one map on $\mathcal{B}^{*}$.

Proof. The proof is similar to Case 1. Denote by $O=\left(\frac{1}{2}, \frac{1}{2}\right)$ the midpoint of $I$. We only prove the $\operatorname{map} \mathcal{T}$ from $\mathcal{B}$ to $\rho_{I}$ fills out half of the line segment $I$. Then we will get the desired result by symmetry.

Let $h$ be a harmonic function on $K$. We consider $\mathcal{T}(\{(c, 0): 0 \leq c \leq 1\})$. When $c=0$, $B(0,0,0)=C_{w}$ and $M_{C_{w}}(h)=\frac{1}{3}\left(h\left(p_{0}\right)+h\left(p_{1}\right)+h\left(p_{2}\right)\right)$. Combining this with the fact that

$$
h\left(p_{2}\right)=\frac{1}{2}\left(h\left(p_{0}\right)+h\left(p_{1}\right)\right),
$$

we get

$$
M_{C_{w}}(h)=\frac{1}{2}\left(h\left(p_{0}\right)+h\left(p_{1}\right)\right) .
$$

Hence $\mathcal{T}(0,0)$ is the midpoint $O$ of $I$. When $c=1, B(1,0,0)=C_{w} \cup C_{0}$, and an easy calculation gives that $M_{C_{w} \cup C_{0}}=h\left(p_{0}\right)$. Hence $\mathcal{T}(1,0)$ is the endpoint $P_{0}$. So if we vary $c$ continuously between 0 and 1 , we can fill out the line segment joining $O$ and $P_{0}$, which is half of $I$.

Case 3. There are two nonjunction points (for example, $p_{1}$ and $p_{2}$ ) among the boundary points of $C_{w}$.

In this case, let $h$ be any harmonic function on $K$. By the matching condition on both points $p_{1}$ and $p_{2}, h$ must be constant on the whole cell $C_{w}$. Hence for every point $x \in C_{w}$, we could view $C_{w}$ itself as the mean value neighborhood of $x$.

We now summarize what we have accomplished:
Theorem 5.1. Given a point $x \in K \backslash V_{0}$, let $C_{w}$ be a cell containing $x$, not intersecting $V_{0}$, and let $D_{w}$ be the union of $C_{w}$ and its neighboring cells of the same size. Then there exists a mean value neighborhood $B$ of $x$ satisfying $C_{w} \subset B \subset D_{w}$. (Note that $B$ is unique within the collection of sets described above.) Moreover, for each point $x \in K \backslash V_{0}$, there exists a system of mean value neighborhoods $B_{k}(x)$ with $\bigcap_{k} B_{k}(x)=\{x\}$.

We should mention here that the result can also be extended to some other p.c.f. fractals including the 3-dimensional Sierpinski gasket. However, it seems that some strong symmetric conditions of both the geometric and the harmonic structures should be required.

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