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# Masayoshi Hata <br> On continuous functions with no unilateral derivatives 

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# ON CONTINUOUS FUNCTIONS WITH NO UNILATERAL DERIVATIVES 

by Masayoshi HATA

## 1. Introduction.

It is known that A. S. Besicovitch in 1925 gave the first example of a continuous function $B(x)$ which has nowhere a unilateral derivative finite or infinite by geometrical process. E. D. Pepper [9] has examined this same function $B(x)$, giving a different exposition. The graph of his function is illustrated in Figure 1. Later, A. N. Singh [12, 13] gave the arithmetical definition of $B(x)$ and constructed an infinite class of such non-differentiable functions. On the other hand, A. P. Morse [8] gave an example of a continuous function $f(x)$ satisfying

$$
\liminf _{s \rightarrow x \pm}\left|\frac{f(s)-f(x)}{s-x}\right|<\limsup _{s \rightarrow x \pm}\left|\frac{f(s)-f(x)}{s-x}\right|=\infty
$$

respectively, for every $x \in(0,1)$, by arithmetical process.
It seems, however, that their methods are somewhat complicated and inappropriate to the study concerning further properties of such functions. In the present paper we shall develop a simple but powerful method to construct and analyze such singular functions by using certain one-dimensional dynamical systems.

The difficulties of finding such functions may be explained by the fact that the set of functions which have nowhere a unilateral derivative finite or infinite is of only the first category in the space of continuous functions (S. Saks [11]), while the set of functions which have nowhere a finite unilateral derivative is of the second category (S. Banach [1], S. Mazurkiewicz [7] and V. Jarnik [5]).

Key-words : Non-differentiable functions - Knot points - Functional equations.


Fig. 1.

## 2. Main Result.

To state our main theorem, we need some definitions and notations. We denote, as usual, the upper and lower derivatives at $x$ of a realvalued function $f(x)$ on the right by $D^{+} f(x), D_{+} f(x)$ respectively. Similarly the upper and lower derivatives, on the left, are denoted by $D^{-} f(x), D_{-} f(x)$ respectively. A point $x$ is said to be a knot point of $f(x)$ provided that

$$
D^{+} f(x)=D^{-} f(x)=\infty \quad \text { and } \quad D_{+} f(x)=D_{-} f(x)=-\infty
$$

The set of knot points of $f(x)$ is denoted by $\operatorname{Knot}(f)$. For a measurable
set $E$, we denote by $|E|$ the Lebesgue measure of $E$. Our theorem can now be stated as follows :

Theorem 2.1. - For any $\alpha \in[0,1)$ and $\varepsilon \in(0,1)$, there exists a continuous function $\psi_{\alpha, \varepsilon}(x)$ defined on the unit interval I sayisfying the following properties:
(1) $\psi_{\alpha, \varepsilon}(x)$ has nowhere a unilateral derivative finite or infinite;
(2) $\left|\operatorname{Knot}\left(\psi_{\alpha, \varepsilon}\right)\right|=\alpha$;
(3) $\psi_{\alpha, \varepsilon}(x)$ satisfies Hölder's condition of order $1-\varepsilon$.

Remark. - K. M. Garg [3] has shown that the set of knot points of Besicovitch's function is of measure zero. He also showed that, for every continuous function defined on $I$ which has nowhere a unilateral derivative finite or infinite, the set of points at which the upper derivative on one side is $+\infty$, the lower derivative on the other side is $-\infty$, and the other two derivatives are finite and equal has a positive measure in every subinterval of $I$; therefore the constant $\alpha$ in our theorem can not be taken to be 1 . Note that the set $\operatorname{Knot}(f)$ is of the second category if $f(x)$ is a continuous function which has nowhere a finite or infinite derivative (W. H. Young [14]).

As a corollary, we have immediately
Corollary 2.2. - For any $\alpha \in[0,2 \pi)$ and $\varepsilon \in(0,1)$, there exists an absolutely convergent cosine Fourier series

$$
\Psi_{\alpha, \varepsilon}(x)=\sum_{n=0}^{\infty} a_{\alpha, \varepsilon, n} \cos n x
$$

satisfying the following properties:
(1) $\Psi_{\alpha, \varepsilon}(x)$ has nowhere a unilateral derivative finite or infinite;
(2) $\left|\operatorname{Knot}\left(\left.\Psi_{\alpha, \varepsilon}\right|_{[0,2 \pi]}\right)\right|=\alpha$;
(3) $\sum_{n=1}^{\infty}\left|a_{\alpha, \varepsilon, n}\right|^{2} n^{2-\varepsilon}<\infty$.

For the proof of Theorem 2.1, we shall introduce a symbol space in section 3 and certain functional equations in section 4. The fundamental properties of the solution are investigated in sections 5 and 6 . We then prove Theorem 2.1 in section 7 using Cantor sets of positive measure.

## 3. Preliminaries.

We first divide the unit interval $I$ into $m$ subintervals

$$
I_{1}=\left[c_{0}, c_{1}\right], I_{2}=\left[c_{1}, c_{2}\right], \ldots, I_{m}=\left[c_{m-1}, c_{m}\right]
$$

where $0=c_{0}<c_{1}<c_{2}<\cdots<c_{m}=1, m \geqslant 2$ and define the address $A(x)$ of a point $x \in I$ by setting $A(x)=j$ for $c_{j-1} \leqslant x<c_{j}, 1 \leqslant j \leqslant m$ and $A\left(c_{m}\right)=m$. Let $g_{j}(x)$ be a strictly monotone, either increasing or decreasing, continuous function defined on the subinterval $I_{j}$ such that $g_{j}\left(I_{j}\right)=I$ for $1 \leqslant j \leqslant m$. Define the sign $\varepsilon_{j}$ to be either +1 or -1 according as $g_{j}$ is monotone increasing or monotone decreasing on $I_{j}$. We assume, in addition, that $g_{1}(x)$ and $g_{m}(x)$ are monotone increasing; so $\varepsilon_{1}=\varepsilon_{m}=+1$.

Let $\Sigma=\{1,2, \ldots, m\}^{N}$ be the one-sided symbol space endowed with the metric

$$
d(w, z)=\sum_{n=1}^{\infty} 2^{-n}\left|w_{n}-z_{n}\right| \quad \text { for } \quad w=\left(w_{n}\right), \quad z=\left(z_{n}\right) \in \Sigma
$$

It is known that $\Sigma$ is a totally disconnected compact metric space. Let $G(x)=g_{A(x)}(x)$ for brevity. Note that the function $G: I \rightarrow I$ is not necessarily continuous. We then define the itinerary $v(x)$ of a point $x \in I$ by setting

$$
v(x)=\left(A_{0}(x), A_{1}(x), \ldots, A_{n}(x), \ldots\right)
$$

where $A_{n}(x)=A\left(G^{n}(x)\right)$ for $n \geqslant 0$. Put $e_{0}=\{0,1\}$ and define the set $e_{n+1}$ inductively by setting $e_{n+1}=\left\{0<x<1 ; G(x) \in e_{n}\right\}$ for $n \geqslant 0$. Obviously \# $e_{n}=m^{n-1}(m-1)$ for $n \geqslant 1$. Let $e=\bigcup_{n \geqslant 0} e_{n}$. Then it is easily verified that the set of discontinuity points of $v$ is precisely equal to the set $e-e_{0}$.

Put $\Lambda_{0}=\left\{v(x) ; x \in e_{0}\right\}$. For $N \geqslant 1$, let $\Lambda_{N}$ be the set of words $w=\left(w_{n}\right) \in \Sigma$ such that either $w_{n}=1$ for $n>N, w_{N} \neq 1$ or $w_{n}=m$ for $n>N, w_{N} \neq m$. Let $\Lambda=\bigcup_{n \geqslant 0} \Lambda_{n}$. Then it is easily seen that for $x \in e-e_{0}$ there exist the limits

$$
\lim _{\varepsilon \rightarrow 0 \pm} v(x+\varepsilon)=\left(A_{0}(x \pm), A_{1}(x \pm), \ldots\right)
$$

in $\Lambda-\Lambda_{0}$ respectively. Note that $v(x)$ is equal to either $v(x+)$ or $v(x-)$. Thus the set $\Lambda_{n}$ consists of the following $2 m^{n-1}(m-1)$ distinct words :

$$
\left\{v(x+) ; x \in e_{n}\right\}+\left\{v(x-) ; x \in e_{n}\right\}
$$

for $n \geqslant 1$. Therefore we have $\Lambda=\Lambda_{0}+\Sigma_{+}+\Sigma_{-}$, where $\Sigma_{+}=\left\{v(x+) ; x \in e-e_{0}\right\}$ and $\Sigma_{-}=\left\{v(x-) ; x \in e-e_{0}\right\}$.

We assume further that each function $h_{j}=g_{j}^{-1}: I \rightarrow I_{j}$ is a contraction; namely the Lipschitz constant

$$
\operatorname{Lip}\left(h_{j}\right)=\sup _{x \neq y \in I}\left|\frac{h_{j}(x)-h_{j}(y)}{x-y}\right|
$$

satisfies $\operatorname{Lip}\left(h_{j}\right)<1$. Let $\gamma=\max _{1 \leqslant j \leqslant m} \operatorname{Lip}\left(h_{j}\right) \in[1 / m, 1)$. We then define the mapping $\mu: \Sigma \rightarrow I$ by setting

$$
\mu(w)=\lim _{n \rightarrow \infty} h_{w_{1}} \circ h_{w_{2}} \circ \cdots \circ h_{w_{n}}(I) \quad \text { for } \quad w=\left(w_{n}\right) \in \Sigma
$$

Clearly $\mu$ is continuous. Then it follows that $X=\mu(\Sigma)$ is a compact subset of $I$ and satisfies the following equality :

$$
X=h_{1}(X) \cup h_{2}(X) \cup \cdots \cup h_{m}(X)
$$

It is known that the above equation possesses a unique non-empty compact solution [4, p. 384]; thus we have $\mu(\Sigma)=X=I$, since $h_{j}(I)=I_{j}$ for $1 \leqslant j \leqslant m$. It also follows that the set $e$ is a dense subset of $I$; therefore the mapping $v$ is one to one.

Let $S_{n}=\bigcup_{0 \leqslant j \leqslant n} e_{j}$ for $n \geqslant 1$ and let

$$
H_{n, x}(y)=h_{A_{0}(x)} \circ h_{A_{1}(x)} \circ \cdots \circ h_{A_{n-1}(x)}(y)
$$

for $n \geqslant 1$ and $x, y \in I$. Obviously $H_{n, x}$ is a contraction satisfying $\operatorname{Lip}\left(H_{n, x}\right) \leqslant \gamma^{n}$. We first consider an arbitrary point $x \in I-e$. Put $K_{n, x}=H_{n, x}(I)$ for $n \geqslant 1$. Since $K_{n, x}$ is the connected component of $I-S_{n}$ containing $x$ and $\left|K_{n, x}\right| \leqslant \gamma^{n}$, we have

$$
\lim _{n \rightarrow \infty} \bar{K}_{n, x}=x
$$

that is, $\mu \circ v(x)=x$. Thus $v$ maps $I-e$ homeomorphically onto
$v(I-e)$. We next consider an arbitrary point $x \in e_{N}, N \geqslant 1$. Put $K_{n, x}^{ \pm}=H_{n, x \pm}(I)$ for $n \geqslant N$, respectively. Since $K_{n, x}^{ \pm}$are the two consecutive connected components of $I-S_{n}$ such that the left end point of $K_{n, x}^{+}$is $x$ and the right end point of $K_{n, x}^{-}$is also $x$, we have

$$
\lim _{n \rightarrow \infty} \bar{K}_{n, x}^{+}=\lim _{n \rightarrow \infty} \bar{K}_{n, x}^{-}=x
$$

so $\mu \circ v(x)=\mu \circ v(x \pm)=x$. Similarly we can define $K_{n, 0}^{+}$and $K_{n, 1}^{-}$ for $n \geqslant 1$; thus $\mu \circ v(0)=0$ and $\mu \circ v(1)=1$. Then we have

Lemma 3.1. $-v(I-e)=\Sigma-\Lambda$; namely, $w=\left(w_{n}\right) \in v(I-e)$ if and only if

$$
\#\left\{n \geqslant 1 ; w_{n} \neq 1\right\}=\infty=\#\left\{n \geqslant 1 ; w_{n} \neq m\right\}
$$

Proof. - Suppose that $w=v(x) \in \Lambda$ for some $x \in I-e$. Since $v$ is one to one, we have $v(I-e) \cap v(e)=\phi$; thus $w \in \Sigma_{+}+\Sigma_{-}$. Hence there exists $y \in e-e_{0}$ such that either $w=v(y+)$ or $w=v(y-)$. Therefore $x=\mu \circ v(x)=\mu(w)=\mu \circ v(y \pm)=y$. This contradiction implies that $\Lambda \cap v(I-e)=\phi$; that is, $v(I-e) \subset \Sigma-\Lambda$. Thus it suffices to show that $\Sigma-\Lambda \subset v(I-e)$.

Suppose now that there exists a word $w=\left(w_{n}\right) \in \Sigma-\Lambda$ such that $w \notin v(I-e)$. Put $z=\left(z_{n}\right) \equiv v \circ \mu(w)$. Then it follows that $w \neq z$. For otherwise, we have $\mu(w) \in e$; thus, $w \in v(e) \subset \Lambda$, contrary to $w \in \Sigma-\Lambda$. Let $N \geqslant 1$ be the smallest integer such that $w_{N} \neq z_{N}$. Since $\mu(w)=\mu \circ v \circ \mu(w)=\mu(z)$, it follows that

$$
h_{w_{N}} \circ h_{w_{N+1}} \circ \cdots=h_{z_{N}} \circ h_{z_{N+1}} \circ \cdots, \text { say } p
$$

Then we have $p \in e_{1}$ and $w, z \in \Lambda_{N}$, contrary to $w \in \Sigma-\Lambda$. This completes the proof.

## 4. Functional Equations.

Let $f_{j}: I \rightarrow I$ be a contraction for $1 \leqslant j \leqslant m$. We assume that $c_{0}=0$ and $c_{m}=1$ are unique fixed points of $f_{1}(x)$ and $f_{m}(x)$ respectively. The following lemma is a special case of the general theorem obtained by the author [4, p. 397], but we include the proof for completeness.

Lemma 4.1. - The functional equations

$$
\begin{equation*}
\psi(x)=f_{j}\left(\psi\left(g_{j}(x)\right)\right) \quad \text { for } \quad x \in I_{j}, 1 \leqslant j \leqslant m \tag{4.1}
\end{equation*}
$$

possess a unique continuous solution $\psi(x)$ if and only if

$$
\begin{equation*}
f_{j}\left(\frac{1+\varepsilon_{j}}{2}\right)=f_{j+1}\left(\frac{1-\varepsilon_{j+1}}{2}\right) \quad \text { for } \quad 1 \leqslant j \leqslant m-1 \tag{4.2}
\end{equation*}
$$

Remark. - This is a generalization of the theorem obtained by G. de Rham [10]; indeed he has shown that the equations

$$
M\left(\frac{x}{2}\right)=F_{0}(M(x)), M\left(\frac{1+x}{2}\right)=F_{1}(M(x)) \quad \text { for } \quad x \in I
$$

possess a unique continuous solution $M(x)$ if and only if $F_{1}\left(p_{0}\right)=F_{0}\left(p_{1}\right)$ where $p_{0}, p_{1}$ are unique fixed points of the contractions $F_{0}, F_{1}$ respectively. Lebesgue's singular functions and Pólya's space-filling curves satisfy the above equations for certain affine contractions $F_{0}$ and $F_{1}$.

Proof. - The conditions (4.2) are obviously necessary; thus it suffices to show the sufficiency. Let $\mathscr{F}$ be the set of continuous functions $u(x)$ defined on $I$ satisfying $u(0)=0$ and $u(1)=1$; obviously $\mathscr{F}$ is a closed subset of the Banach space $C([0,1])$ with the usual uniform norm. We now consider the following operator :

$$
T u(x)=f_{A(x)}(u(G(x)))
$$

Then it is easily seen that the conditions (4.2) imply that $T(\mathscr{F}) \subset \mathscr{F}$; moreover $T$ is a contraction, since

$$
\|T u-T v\| \leqslant \lambda \max _{x \in I}|u(G(x))-v(G(x))| \leqslant \lambda\|u-v\|
$$

where $\lambda=\max _{1 \leqslant j \leqslant m} \operatorname{Lip}\left(f_{j}\right) \in[1 / m, 1)$, for any $u, v \in \mathscr{F}$. Hence $T$ has a unique fixed point $\psi$ in $\mathscr{F}$; namely

$$
\psi(x)=f_{j}\left(\psi\left(g_{j}(x)\right)\right) \quad \text { for } \quad c_{j-1} \leqslant x<c_{j}, \quad 1 \leqslant j \leqslant m
$$

Obviously this equality holds also true for $x=c_{j}$. This completes the proof.

For $n \geqslant 1$ and $x, y \in I$, we define

$$
F_{n, x}(y)=f_{A_{0}(x)} \circ f_{A_{1}(x)} \circ \cdots \circ f_{A_{n-1}(x)}(y)
$$

The function $F_{n, x}$ is a contraction satisfying $\operatorname{Lip}\left(F_{n, x}\right) \leqslant \lambda^{n}$. Put $\beta=\max _{1 \leqslant j \leqslant m} \operatorname{Lip}\left(g_{j}\right) \in[m, \infty]$. Then we have

Lemma 4.2. - Suppose that $\left\{f_{j}\right\}$ satisfy the conditions (4.2). If $\beta<\infty$, then the continuous solution $\psi(x)$ satisfies Hölder's condition of order $\log (1 / \lambda) / \log \beta$.

Proof. - Consider arbitrary two points $x<y$ in $I$. Let $N \geqslant 0$ be the smallest integer satisfying $\#\left\{S_{N+1} \cap(x, y)\right\} \geqslant 2$. We now distinguish two cases: (a) $S_{N} \cap(x, y)=\phi$; (b) $S_{N} \cap(x, y)$ consists of a single point, say $p$. In case ( $a$ ), it follows that

$$
\begin{aligned}
|\psi(x)-\psi(y)| & =\lim _{\varepsilon \rightarrow 0^{+}}|\psi(x+\varepsilon)-\psi(y-\varepsilon)| \\
& =\lim _{\varepsilon \rightarrow 0^{+}}\left|F_{N, x+\varepsilon}\left(\psi\left(G^{N}(x+\varepsilon)\right)\right)-F_{N, x+\varepsilon}\left(\psi\left(G^{N}(y-\varepsilon)\right)\right)\right| \leqslant \lambda^{N}
\end{aligned}
$$

Similarly we have $|\psi(x)-\psi(y)| \leqslant 2 \lambda^{N}$ in case (b), since $(x, p) \cap S_{N}=(p, y) \cap S_{N}=\phi$. Now let $s<t$ be any two consecutive points of $e_{N+1}$ contained in $(x, y)$. Then it follows that $|x-y|>|s-t| \geqslant \beta^{-N-1} ;$ thus

$$
|\psi(x)-\psi(y)| \leqslant 2 \lambda^{N}=\frac{2}{\lambda} \beta^{-\xi(N+1)} \leqslant \frac{2}{\lambda}|x-y|^{\xi}
$$

where $\xi=\log (1 / \lambda) / \log \beta$, which obviously completes the proof.

## 5. Some Properties.

The continuous solution $\psi(x)$ of the equations (4.1) is not necessarily singular in general; for example, if we take

$$
g_{j}(x)=m x-j+1 \quad \text { and } \quad f_{j}(x)=\frac{x}{m}+\frac{j-1}{m}
$$

for $1 \leqslant j \leqslant m$, then obviously $\psi(x) \equiv x$ is a smooth solution of (4.1). In this paper, to discuss the singularities of $\psi(x)$, we shall restrict ourselves to the following case:

$$
\begin{equation*}
\varepsilon_{j}=1+2\left[\frac{j}{4}\right]-2\left[\frac{j+1}{4}\right] \tag{5.1}
\end{equation*}
$$

and

$$
f_{j}(x)=\frac{1}{2 k}\left\{(-1)^{[j / 2]} x+\left[\frac{j}{2}\right]-\left[\frac{j}{4}\right]+\left[\frac{j-1}{4}\right]\right\}
$$

for $1 \leqslant j \leqslant m=4 k$, where $k$ is a positive integer ; so $\lambda=1 / 2 k$. Then it is easily seen that the functions $\left\{f_{j}\right\}$ satisfy the conditions (4.2); therefore the equations (4.1) possess a unique continuous solution $\psi(x)$, which depends only on the functions $\left\{g_{j}\right\}$ satisfying the conditions (5.1). Let $\eta_{j}$ be the sign of the function $f_{j}$; namely $\eta_{j}=(-1)^{[j / 2]}$, for $1 \leqslant j \leqslant 4 k$. For brevity, put

$$
\varepsilon_{n, x}=\prod_{j=0}^{n-1} \varepsilon_{A_{j}(x)} \quad \text { and } \quad \eta_{n, x}=\prod_{j=0}^{n-1} \eta_{A_{j}(x)}
$$

for $n \geqslant 1, x \in I$.
Consider now an arbitrary point $x \in I-e$. We define

$$
p_{j, n, x}=H_{n, x}\left(c_{j}\right) \quad \text { for } \quad n \geqslant 1, \quad 0 \leqslant j \leqslant 4 k
$$

Obviously $p_{j, n, x} \neq x$. Since $p_{j, n, x} \in G^{-n}\left(c_{j}\right) \subset e_{n+1}$ for $1 \leqslant j \leqslant 4 k-1$, we have

$$
G^{n}\left(p_{j, n, x}\right)=c_{j} \quad \text { for } \quad 1 \leqslant j \leqslant 4 k-1
$$

The points $p_{0, n, x}$ and $p_{4 k, n, x}$ are two end points of $K_{n, x}$ and do not satisfy the above equality in general ; however,

$$
\lim _{\substack{y \rightarrow p_{j, n, x} \\ y \in \mathbb{K}_{n, x}}} G^{n}(y)=c_{j} \quad \text { for } \quad j=0,4 k
$$

Note that $0<\left|x-p_{j, n, x}\right|<\gamma^{n}$ for any $n \geqslant 1$. Then we have

Lemma 5.1. - Suppose that $x \in I-e$. Then the points $\left\{p_{j, n, x}\right\}$ satisfy the following properties :
(1) $\operatorname{sign}\left(x-p_{j, n, x}\right)=\varepsilon_{n, x} \operatorname{sign}\left\{A_{n}(x)-j-\frac{1}{2}\right\}$,
(2) $\psi(x)-\psi\left(p_{j, n, x}\right)=\frac{\eta_{n, x}}{(2 k)^{n}}\left\{\psi\left(G^{n}(x)\right)-\frac{1-(-1)^{j}}{4 k}-\frac{1}{k}\left[\frac{j}{4}\right]\right\}$
for $n \geqslant 1$ and $0 \leqslant j \leqslant 4 k$.
Proof. - Since $p_{j, n, x}=H_{n, x}\left(c_{j}\right)$, we have
$\operatorname{sign}\left(x-p_{j, n, x}\right)=\operatorname{sign}\left\{H_{n, x}\left(G^{n}(x)\right)-H_{n, x}\left(c_{j}\right)\right\}=\varepsilon_{n, x} \operatorname{sign}\left\{G^{n}(x)-c_{j}\right\} ;$
thus the property (1) follows immediately. Since $K_{n, x} \cap S_{n}=\phi$,

$$
\psi\left(p_{j, n, x}\right)=\lim _{\substack{y \rightarrow p_{j, n, x} \\ y \in K_{n, x}}} \psi(y)=\lim _{\substack{y \rightarrow p_{j, n, x} \\ y \in K_{n, x}}} F_{n, x}\left(\psi\left(G^{n}(y)\right)\right)=F_{n, x}\left(\psi\left(c_{j}\right)\right)
$$

for $0 \leqslant j \leqslant 4 k$; hence

$$
\psi(x)-\psi\left(p_{j, n, x}\right)=\mathrm{F}_{n, x}\left(\psi\left(\mathrm{G}^{n}(x)\right)\right)-\mathrm{F}_{n, x}\left(\psi\left(c_{j}\right)\right) \frac{\eta_{n, x}}{(2 k)^{n}}\left\{\psi\left(G^{n}(x)\right)-\psi\left(c_{j}\right)\right\}
$$

which obviously completes the proof.
We now consider an arbitrary point $x \in e_{N}, N \geqslant 1$. Then it is easily seen that, for $1 \leqslant j \leqslant 4 k-1$, each of the sets $K_{n, x}^{ \pm}$contains exactly one point of $G^{-n}\left(c_{j}\right) \subset e_{n+1}$, say $q_{j, n, x}^{ \pm}$respectively. Obviously $q_{j, n, x}^{ \pm} \neq x$. Similarly we can define $\left\{q_{j, n, 0}^{+}\right\}$and $\left\{q_{j, n, 1}^{-}\right\}$for $n \geqslant 0$, $1 \leqslant j \leqslant 4 k-1$. Note that $0<\left|x-q_{j, n, x}^{ \pm}\right|<\gamma^{n}$ for any $n \geqslant N$. It also follows that

$$
\lim _{\varepsilon \rightarrow 0 \pm} G^{n}(x+\varepsilon)=\frac{1}{2}\left(1 \mp \varepsilon_{N, x \pm}\right)
$$

for every $n \geqslant N$, respectively. We, of course, adopt the rule: $\varepsilon_{0,0^{+}}=\varepsilon_{0,1_{-}}=\eta_{0,0^{+}}=\eta_{0,1_{-}}=1$. Then we have

Lemma 5.2. - Suppose that $x \in e_{N}, N \geqslant 0$. Then the points $\left\{q_{j, n, x}^{ \pm}\right\}$satisfy the following :

$$
\psi(x)-\psi\left(q_{j, n, x}^{ \pm}\right)=\frac{\eta_{N, x \pm}}{(2 k)^{n}}\left\{\frac{1}{2}\left(1 \mp \varepsilon_{N, x \pm}\right)-\frac{1-(-1)^{j}}{4 k}-\frac{1}{k}\left[\frac{j}{4}\right]\right\}
$$

for $n \geqslant N$ and $1 \leqslant j \leqslant 4 k-1$, respectively.
Proof. - Since $K_{n, x}^{ \pm} \cap S_{n}=\phi$, we have

$$
\begin{aligned}
& \psi(x)-\psi\left(q_{j, n, x}^{ \pm}\right)=\lim _{\varepsilon \rightarrow 0 \pm}\left\{\psi(x+\varepsilon)-\psi\left(q_{j, n, x}^{ \pm}\right)\right\}= \\
& \lim _{\varepsilon \rightarrow 0 \pm}\left\{F_{n, x+\varepsilon}\left(\psi\left(G^{n}(x+\varepsilon)\right)\right)-F_{n, x+\varepsilon}\left(\psi\left(c_{j}\right)\right)\right\}=\frac{\eta_{N, x \pm}}{(2 k)^{n}}\left\{\frac{1}{2}\left(1 \mp \varepsilon_{N, x \pm}\right)-\psi\left(c_{j}\right)\right\}
\end{aligned}
$$

for every $n \geqslant N$, respectively. This completes the proof.

## 6. Singularities.

For any $x \neq y \in I$, we define $\Delta \psi(x, y)=(\psi(x)-\psi(y)) /(x-y)$. Let $W$ be the set of points $x \in I$ at which $A_{n}(x) \equiv 2$ or $3(\bmod 4)$ for infinitely many $n$ 's. Obviously $W \subset I-e$. First of all, we have

Theorem 6.1. - Suppose that $\gamma \leqslant 1 / 2 k$. Then we have

$$
D^{ \pm} \psi(x) \geqslant 0 \geqslant D_{ \pm} \psi(x) \quad \text { and } \quad D^{ \pm} \psi(x)-D_{ \pm} \psi(x) \geqslant 1 / 4 k
$$

respectively, for every $x \in W$.
Proof. - We distinguish two cases (not exclusive) as follows :
Case $A . A_{n}(x) \equiv 3(\bmod 4)$ for infinitely many $n ' s$.
Let $0<n_{1}<n_{2}<\cdots$ be the subsequence of integers such that $A_{n_{i}}(x)=4 N_{i}+3$, where $0 \leqslant N_{i}<k$. From the functional equations (4.1), we have

$$
\frac{N_{i}}{k} \leqslant \psi\left(G^{n_{i}}(x)\right) \leqslant \frac{2 N_{i}+1}{k}
$$

therefore $\left\{\psi(x)-\psi\left(P_{i, 1}\right)\right\}\left\{\psi(x)-\psi\left(P_{i, 2}\right)\right\} \leqslant 0 \quad$ by (2) of Lemma 5.1, where $p_{i, j}=p_{4 N_{i}+j, n_{i}, x}$ for $0 \leqslant j \leqslant 4$. On the order hand, we have $\operatorname{sign}\left(x-P_{i, 1}\right)=\operatorname{sign}\left(x-P_{i, 2}\right)=\varepsilon_{n_{i}, x}$ by (1) of Lemma 5.1. Since $\varepsilon_{n_{i}, x}$ changes the sign infinitely many times as $i$ increases, it follows that $D^{ \pm} \psi(x) \geqslant 0 \geqslant D_{ \pm} \psi(x)$. It also follows that

$$
\left|\Delta \psi\left(x, P_{i, 1}\right)\right|+\left|\Delta \psi\left(x, P_{i, 2}\right)\right| \geqslant \frac{(2 k)^{-n_{i}-1}}{\left|x-P_{i, 1}\right|}>\frac{1}{2 k}(2 k \gamma)^{-n_{i}} \geqslant \frac{1}{2 k}
$$

therefore $D^{ \pm} \psi(x)-D_{ \pm} \psi(x) \geqslant 1 / 4 k$ respectively, as required.
Case $B . A_{n}(x) \equiv 2(\bmod 4)$ for infinitely many $n$ 's.
Let $0<n_{1}<n_{2}<\cdots$ be the subsequence of integers such that $A_{n_{i}}(x)=4 N_{i}+2$, where $0 \leqslant N_{i}<k$. Since

$$
\frac{N_{i}}{k} \leqslant \psi\left(G^{n_{i}}(x)\right) \leqslant \frac{2 N_{i}+1}{k}
$$

it is easily seen that $\left\{\psi(x)-\psi\left(P_{i, 0}\right)\right\}\left\{\psi(x)-\psi\left(P_{i, 1}\right)\right\} \leqslant 0$ and $\left\{\psi(x)-\psi\left(P_{i, 2}\right)\right\}\left\{\psi(x)-\psi\left(P_{i, 3}\right)\right\} \leqslant 0$. On the other hand, we have
$\operatorname{sign}\left(x-P_{i, 0}\right)=\operatorname{sign}\left(x-P_{i, 1}\right)=\operatorname{sign}\left(P_{i, 2}-x\right)=\operatorname{sign}\left(P_{i, 3}-x\right) ;$ therefore $D^{ \pm} \psi(x) \geqslant 0 \geqslant D_{ \pm} \psi(x)$. Moreover,

$$
\left|\Delta \psi\left(x, P_{i, 0}\right)\right|+\left|\Delta \psi\left(x, P_{i, 1}\right)\right| \geqslant \frac{(2 k)^{-n_{i}-1}}{\left|x-P_{i, 0}\right|}>\frac{1}{2 k}(2 k \gamma)^{-n_{i}} \geqslant \frac{1}{2 k}
$$

The same estimate holds true if we replace $P_{i, 0}, P_{i, 1}$ by $P_{i, 2}, P_{i, 3}$, respectively; thus $D^{ \pm} \psi(x)-D_{ \pm} \psi(x) \geqslant 1 / 4 k$ respectively. This completes the proof.

Let $W_{0} \subset W$ be the set of points $x \in I$ at which $A_{n}(x) \equiv 2$ or 3 $(\bmod 4)$ and $A_{n+1}(x) \equiv 2$ or $3(\bmod 4)$ for infinitely many $n$ 's. Then we have

Theorem 6.2. - Suppose that $\gamma \leqslant 1 / 2 k$. Then $W_{0}$ is contained in the set $\operatorname{Knot}(\psi)$ except for a set of measure zero.

Proof. - We consider an arbitrary point $x$ of $W_{0}$. Let $0 \leqslant n_{1}<n_{2}<\cdots$ be the subsequence of integers such that $A_{n_{i}}(x)=4 N_{i}+\delta_{i}$ and $A_{n_{i}+1}(x)=4 L_{i}+\omega_{i}$, where $0 \leqslant N_{i}, L_{i}<k$ and $2 \leqslant \delta_{i}, \omega_{i} \leqslant 3$. Then it is easily seen that

$$
\frac{2 N_{i}+1}{2 k}-\frac{2 L_{i}+1}{(2 k)^{2}} \leqslant \psi\left(G^{n_{i}}(x)\right) \leqslant \frac{2 N_{i}+1}{k}-\frac{L_{i}}{2 k^{2}} ;
$$

therefore by (2) of Lemma 5.1,

$$
\eta_{n_{i}, x}(2 k)^{n_{i}}\left\{\psi(x)-\psi\left(P_{i, 0}\right)\right\}=\psi\left(G^{n_{i}}(x)\right)-\frac{N_{i}}{k} \geqslant \frac{1}{2 k}-\frac{2 L_{i}+1}{(2 k)^{2}} \geqslant(2 k)^{-2} .
$$

Similarly we have

$$
\eta_{n_{i}, x}(2 k)^{n_{i}}\left\{\psi\left(P_{i, 4}\right)-\psi(x)\right\}=\frac{N_{i}+1}{k}-\psi\left(G^{n_{i}}(x)\right) \geqslant \frac{1}{2 k}+\frac{L_{i}}{2 k^{2}} \geqslant \frac{1}{2 k}
$$

Therefore, since $\operatorname{sign}\left(x-P_{i, 0}\right)=\operatorname{sign}\left(P_{i, 4}-x\right)$, it follows that

$$
\operatorname{sign}\left(\Delta \psi\left(x, P_{i, 0}\right)\right)=\operatorname{sign}\left(\Delta \psi\left(x, P_{i, 4}\right)\right)
$$

and

$$
\left|\Delta \psi\left(x, P_{i, 0}\right)\right|>(2 k)^{-2}, \quad\left|\Delta \psi\left(x, P_{i, 4}\right)\right|>\frac{1}{2 k}
$$

Hence the set $\left[D_{+} \psi(x), D^{+} \psi(x)\right] \cap\left[D_{-} \psi(x), D^{-} \psi(x)\right]$ contains an interval of length $(2 k)^{-2}$ by Theorem 6.1. Thus it follows from Denjoy's theorem
[2, p. 105] that except for a set of measure zero, every point of $W_{0}$ is a knot point of $\psi(x)$. This completes the proof.

For $N \geqslant 0$, let $Y_{N}$ be the set of points $x \in I$ at which $A_{n}(x) \equiv 0$ or $1(\bmod 4)$ for all $n \geqslant N$ and $A_{N-1}(x) \equiv 2$ or $3(\bmod 4)$. Obviously $I-W=\bigcup_{n \geqslant 0} Y_{n}$. For brevity, put $Y_{n}^{*}=Y_{n} \cap(I-e)$ for $n \geqslant 0$. Then the unit interval $I$ is decomposed as follows:

$$
I=W+e+\bigcup_{n \geqslant 0} Y_{n}^{*}
$$

For $n \geqslant 1$, let $\Xi_{n}$ be the set of finite words ( $w_{1}, \ldots w_{n}$ ) of length $n$ such that $1 \leqslant w_{j} \leqslant 4 k$ and $w_{j} \equiv 0$ or $1(\bmod 4)$ for $1 \leqslant j \leqslant n$. Then we have

Theorem 6.3. - Suppose that there exists a positive constant $C_{0}$, independent of $n$, satisfying

$$
\min _{\left(w_{1} \ldots w_{n}\right) \in \Xi_{n}}\left|h_{w_{1}} \circ \ldots \circ h_{w_{n}}(I)\right| \geqslant C_{0}(2 k)^{-n}
$$

for all $n \geqslant 1$. Suppose further that $\beta<\infty$. Then we have

$$
D^{ \pm} \psi(x)-D_{ \pm} \psi(x) \geqslant \frac{1}{2 k}
$$

respectively, for every $x \in I-W$.
Proof: - We distinguish two cases as follows :
Case A. $x \in Y_{N}^{*}$ for some $N \geqslant 0$.
By Lemma 3.1, we have $A_{n}(x) \neq 1$ for infinitely many $n$ 's. Let $N \leqslant n_{1}<n_{2}<\cdots$ be the subsequence of integers such that $A_{n_{i}}(x) \geqslant 4$. Put $Q_{i, j}=p_{j, n_{i}, x}$ for $0 \leqslant j \leqslant 2$. Since

$$
\psi\left(G^{n_{i}}(x)\right) \geqslant \frac{1}{2 k}
$$

and $\operatorname{sign}\left(x-Q_{i, 1}\right)=\operatorname{sign}\left(x-Q_{i, 2}\right)=\operatorname{sign}\left(Q_{i, 2}-Q_{i, 1}\right)=\varepsilon_{N, x}$, we have $\left|\Delta \psi\left(x, Q_{i, 1}\right)-\Delta \psi\left(x, Q_{i, 2}\right)\right|=$

$$
(2 k)^{-n_{i}}\left|\psi\left(G^{n_{i}}(x)\right)\left\{\frac{1}{x-Q_{i, 2}}-\frac{1}{x-Q_{i, 1}}\right\}+\frac{1 \cdots}{2 k\left(x-Q_{i, 1}\right)}\right| \geqslant \frac{(2 k)^{-n_{i}-1}}{\left|x-Q_{i, 1}\right|}>\frac{1}{2 k}
$$

On the other hand, it follows that
$\left|x-Q_{i, 0}\right|>\left|Q_{i, 1}-Q_{i, 0}\right| \geqslant \beta^{-N}\left|h_{A_{N^{(x)}}} \circ \ldots \circ h_{A_{n_{i}-1}(x)} \circ h_{1}(I)\right| \geqslant$

$$
C_{0} \beta^{-N}(2 k)^{-n_{i}+N-1}
$$

therefore

$$
\left|\Delta \psi\left(x, Q_{i, 0}\right)\right|=(2 k)^{-n_{i}}\left|\frac{\psi\left(G^{n_{i}}(x)\right)}{x-Q_{i, 0}}\right| \leqslant \frac{2 k}{C_{0}}\left(\frac{\beta}{2 k}\right)^{N}
$$

Since $\operatorname{sign}\left(x-Q_{i, 0}\right)=\varepsilon_{N, x}$, we conclude that either $\left[D_{+} \psi(x), D^{+} \psi(x)\right]$ or $\left[D_{-} \psi(x), D^{-} \psi(x)\right]$ contains an interval of length $1 / 2 k$ according as $\varepsilon_{N, x}=-1$ or +1 .

It also follows from Lemma 3.1 that $A_{n}(x) \neq 4 k$ for infinitely many $n$ 's. Let $N \leqslant n_{1}<n_{2}<\ldots$ be the subsequence of integers such that $A_{n_{i}}(x) \leqslant 4 k-3$. Put $R_{i, j}=p_{4 k-j, n_{i}, x}$ for $0 \leqslant j \leqslant 3$. Since

$$
\psi\left(G^{n_{i}}(x)\right) \leqslant \frac{2 k-1}{2 k}
$$

and $\operatorname{sign}\left(x-R_{i, 2}\right)=\operatorname{sign}\left(x-R_{i, 3}\right)=\operatorname{sign}\left(R_{i, 3}-R_{i, 2}\right)=-\varepsilon_{N, x}$, we have $\left|\Delta \psi\left(x, R_{i, 2}\right)-\Delta \psi\left(x, R_{i, 3}\right)\right|=$

$$
\begin{array}{r}
(2 k)^{-n_{i}}\left|\left\{\frac{2 k-1}{2 k}-\psi\left(G^{n_{i}}(x)\right)\right\}\left\{\frac{1}{x-R_{i, 3}}-\frac{1}{x-R_{i, 2}}\right\}+\frac{1}{2 k\left(x-R_{i, 2}\right)}\right| \geqslant \\
\frac{(2 k)^{-n_{i}-1}}{\left|x-R_{i, 2}\right|}>\frac{1}{2 k}
\end{array}
$$

On the other hand, $\left|x-R_{i, 0}\right|>\left|R_{i, 1}-R_{i, 0}\right| \geqslant C_{0} \beta^{-N}(2 k)^{-n_{i}+N-1}$; thus

$$
\left|\Delta \psi\left(x, R_{i, 0}\right)\right|=(2 k)^{-n_{i}}\left|\frac{\psi\left(G^{n_{i}}(x)\right)-1}{x-R_{i, 0}}\right| \leqslant \frac{2 k}{C_{0}}\left(\frac{\beta}{2 k}\right)^{N}
$$

Since $\operatorname{sign}\left(x-R_{i, 0}\right)=-\varepsilon_{N, x}$, it follows that either $\left[D_{+} \psi(x), D^{+} \psi(x)\right]$ or $\left[D_{-} \psi(x), D^{-} \psi(x)\right]$ contains an interval of length $1 / 2 k$ according as $\varepsilon_{N, x}=+1$ or -1 . Hence $D^{ \pm} \psi(x)-D_{ \pm} \psi(x) \geqslant 1 / 2 k$ respectively.

Case B. $x \in e_{N}$ for some $N \geqslant 0$.
For $n \geqslant N$, let $Q_{n}^{+}=\max \left\{q_{1, n, x}^{+}, q_{3, n, x}^{+}\right\}, Q_{n}^{-}=\min \left\{q_{1, n, x}^{-}, q_{3, n, x}^{-}\right\}$ and let $R_{n}^{+}=q_{2, n, x}^{+}$respectively. Then $Q_{n}^{-}<R_{n}^{-}<x<R_{n}^{+}<Q_{n}^{+}$.

Since $\operatorname{sign}\left(x-Q_{n}^{ \pm}\right)=\operatorname{sign}\left(Q_{n}^{ \pm}-R_{n}^{ \pm}\right)= \pm 1$ respectively, it follows from Lemma 5.2 that

$$
\begin{aligned}
& \left|\Delta \psi\left(x, R_{n}^{ \pm}\right)-\Delta \psi\left(x, Q_{n}^{ \pm}\right)\right|= \\
& \quad(2 k)^{-n}\left|\frac{1}{2}\left(1 \mp \varepsilon_{N, x \pm}\right)\left\{\frac{1}{x-R_{n}^{ \pm}}-\frac{1}{x-Q_{n}^{ \pm}}\right\}+\frac{1}{2 k\left(x-Q_{n}^{ \pm}\right)}\right| \geqslant \frac{(2 k)^{-n-1}}{\left|x-Q_{n}^{ \pm}\right|}>\frac{1}{2 k},
\end{aligned}
$$

respectively. On the other hand, we have

$$
\left|x-R_{n}^{ \pm}\right|>\left|K_{n+1, x}^{ \pm}\right| \geqslant \beta^{-N}\left|h_{A_{N}(x \pm)} \circ \cdots \circ h_{A_{n}(x \pm)}(I)\right| \geqslant C_{0} \beta^{-N}(2 k)^{-n+N-1}
$$

therefore

$$
\left|\Delta \psi\left(x, R_{n}^{ \pm}\right)\right| \leqslant \frac{(2 k)^{-n}}{\left|x-R_{n}^{ \pm}\right|}<\frac{2 k}{C_{0}}\left(\frac{\beta}{2 k}\right)^{N}
$$

Hence $D^{ \pm} \psi(x)-D_{ \pm} \psi(x) \geqslant 1 / 2 k$ respectively. This completes the proof.

Let $Y^{*}=\bigcup_{n \geqslant 0} Y_{n}^{*}$ for brevity. Then we have
Theorem 6.4. $-\operatorname{Knot}(\psi) \cap Y^{*}=\phi$.
Proof. - We consider an arbitrary point $x$ of $Y_{N}^{*}$ for some $N \geqslant 0$. Let $s_{n}=p_{0, n, x}$ for $n \geqslant N$. Since $\operatorname{sign}\left(x-s_{n}\right)=\varepsilon_{N, x}$ is independent of $n \geqslant N$, the sequence $\left\{s_{n}\right\}$ is monotone, either increasing or decreasing, and converges to $x$. Note that $s_{n}=s_{n+1}$ if and only if $A_{n}(x)=1$. Put $J_{n}=\left[s_{n}, s_{n+1}\right] \subset \bar{K}_{n, x}$ for $n \geqslant N$. Then it is easily seen that

$$
\left(x, s_{N}\right]=\bigcup_{n \geqslant N} J_{n} .
$$

Since the function $G^{n}(x)$ maps $K_{n, x}$ homeomorphically onto $(0,1)$, we have $A_{n}(x)>A_{n}(y)$ for all $y \in \dot{J}_{n}$. Therefore

$$
\psi\left(G^{n}(x)\right) \geqslant f_{A_{n}(x)}(0) \geqslant \max _{j<A_{n}(x)}\left\|f_{j}\right\| \geqslant \psi\left(G^{n}(y)\right) ;
$$

thus

$$
\begin{array}{r}
\eta_{N, x} \operatorname{sign}\{\psi(x)-\psi(y)\}=\eta_{N, x} \operatorname{sign}\left\{F_{n, x}\left(\psi\left(G^{n}(x)\right)\right)-F_{n, x}\left(\psi\left(G^{n}(y)\right)\right)\right\}= \\
\operatorname{sign}\left\{\psi\left(G^{n}(x)\right)-\psi\left(G^{n}(y)\right)\right\} \geqslant 0 .
\end{array}
$$

By the continuity of $\psi$, we conclude that

$$
\eta_{N, x} \operatorname{sign}\{\psi(x)-\psi(y)\} \geqslant 0 \quad \text { for every } \quad y \in\left[x, s_{N}\right]
$$

This means that $x$ is not a knot point of $\psi(x)$.

## 7. Proof of Theorem 2.1.

First of all, for any integer $k \geqslant 1$ and positive numbers $\sigma, \tau, \rho$ satisfying

$$
\begin{equation*}
2 k(\sigma+\tau)<1 \quad \text { and } \quad \sigma \geqslant \rho \tag{7.1}
\end{equation*}
$$

we shall construct two Cantor sets $E_{0} \equiv E_{0}(k, \sigma, \tau)$ and $E_{1} \equiv E_{1}(k, \sigma, \rho)$. The set $E_{0}(k, \sigma, \tau)$ is obtained from the unit interval $I$ by a sequence of deletions of open intervals known as middle thirds, as follows : First divide $I$ into $k$ equal parts, say

$$
I_{1,1}=\left[0, \frac{1}{k}\right], \quad I_{1,2}=\left[\frac{1}{k}, \frac{2}{k}\right], \quad \ldots, \quad I_{1, k}=\left[\frac{k-1}{k}, 1\right],
$$

and remove from each closed interval $I_{1, j}$ the open interval $U_{1, j}$ centered at $(2 j-1) / 2 k$ and of length $2 \sigma$. We subdivide each of the $2 k$ remaining closed intervals into $k$ equal parts, say $I_{2, j}, 1 \leqslant j \leqslant 2 k^{2}$, ordered from left to right, each of length $(1-2 k \sigma) /\left(2 k^{2}\right)$. Then remove from each closed interval $I_{2, j}$ the middle open interval $U_{2, j}$ of length $2 \sigma \tau$, leaving the $4 k^{2}$ closed intervals, each of length $\left(1-2 k \sigma-4 k^{2} \sigma \tau\right) /\left(4 k^{2}\right)$. This process is permitted to continue indefinitely. At the $n$th stage of deletion, each length of the $2^{n-1} k^{n}$ open intervals $U_{n, j}$ is $2 \sigma \tau^{n-1}$, and therefore the measure of the union of the open intervals removed in the entire sequence of removal operations is $2 k \sigma /(1-2 k \tau)$. The set $E_{0}$ is defined to be the closed set remaining; thus

$$
\left|E_{0}\right|=\frac{1-2 k(\sigma+\tau)}{1-2 k \tau}
$$

We next define the set $E_{1}(k, \sigma, \rho)$, which is slightly different from $E_{0}$ defined above, as follows : First divide the unit interval $I$ into $k$ equal parts, say

$$
J_{1,1}=\left[0, \frac{1}{k}\right], \quad J_{1,2}=\left[\frac{1}{k}, \frac{2}{k}\right], \quad \ldots, \quad J_{1, k}=\left[\frac{k-1}{k}, 1\right] .
$$

Then remove from each closed interval $J_{1, j}$ the two intervals

$$
V_{1, j}^{-}=\left[\frac{j-1}{k}, \frac{2 j-1-2 k \sigma}{2 k}\right), V_{1, j}^{+}=\left(\frac{2 j-1+2 k \sigma}{2 k}, \frac{j}{k}\right]
$$

each of length $(1-2 k \sigma) / 2 k$. We subdivide each of the $k$ remaining closed intervals into $2 k$ equal parts, say $J_{2, j}, 1 \leqslant j \leqslant 2 k^{2}$, ordered
from left to right, each of length $\sigma / k$. Then delete from each closed interval $J_{2, j}$ the two intervals $V_{2, j}^{ \pm}$of length $\rho(1-2 k \sigma) / 2 k$, leaving the $2 k^{2}$ middle closed intervals, each of length $(\sigma-\rho+2 k \sigma \rho) / k$. At the $n$th stage of deletion, we have $\left|V_{n, j}^{ \pm}\right|=\rho^{n-1}(1-2 k \sigma) / 2 k$; therefore the measure of the union of the removed intervals in the entire sequence of removal operations is $(1-2 k \sigma) /(1-2 k \rho)$. The set $E_{1}$ is defined to be the closed set remaining ; thus

$$
\left|E_{1}\right|=\frac{2 k(\sigma-\rho)}{1-2 k \rho}
$$

Note that the set $E_{1}$ is contained in $\left[\frac{1-2 k \sigma}{2 k(1-\rho)}, 1-\frac{1-2 k \sigma}{2 k(1-\rho)}\right]$.
We now define the continuous function $\zeta_{0}(x) \equiv \zeta_{0}(k, \sigma, \tau ; x)$ by setting

$$
\zeta_{0}(x)=\int_{0}^{x} d_{0}(s) d s \quad \text { for } \quad 0 \leqslant x \leqslant 1
$$

where $d_{0}(s)=1 / 2 k$ if $s \in E_{0}(k, \sigma, \tau)$ and $d_{0}(s)=\tau$ otherwise. We also define the continuous function $\zeta_{1}(x) \equiv \zeta_{1}(k, \sigma, \rho ; x)$ by setting

$$
\zeta_{1}(x)=\frac{1}{2 k}-\sigma+\int_{0}^{x} d_{1}(s) d s \quad \text { for } \quad 0 \leqslant x \leqslant 1
$$

where $d_{1}(s)=1 / 2 k$ if $s \in E_{1}(k, \sigma, \rho)$ and $d_{1}(s)=\rho$ otherwise. Then it is easily seen that $\zeta_{0}(I)=[0,(1-2 k \sigma) / 2 k], \zeta_{1}(I)=[(1-2 k \sigma) / 2 k, 1 / 2 k]$ and $\zeta_{i}\left(E_{i}\right)=E_{i} \cap \zeta_{i}(I)$ for $i=0,1$.

We next define, for $0 \leqslant i<k$,

$$
\begin{aligned}
& g_{4 i+1}(x)=\zeta_{0}^{-1}\left(x-\frac{i}{k}\right) \quad \text { for } \quad x \in I_{4 i+1}=\left[\frac{i}{k}, \frac{2 i+1}{2 k}-\sigma\right] \text {, } \\
& g_{4 i+2}(x)=\zeta_{1}^{-1}\left(x-\frac{i}{k}\right) \quad \text { for } \quad x \in I_{4 i+2}=\left[\frac{2 i+1}{2 k}-\sigma, \frac{2 i+1}{2 k}\right] \text {, } \\
& g_{4 i+3}(x)=\zeta_{1}^{-1}\left(\frac{i+1}{k}-x\right) \quad \text { for } \quad x \in I_{4 i+3}=\left[\frac{2 i+1}{2 k}, \frac{2 i+1}{2 k}+\sigma\right] \text {, } \\
& g_{4 i+4}(x)=\zeta_{0}^{-1}\left(x-\frac{2 i+1}{2 k}-\sigma\right) \text { for } x \in I_{4 i+4}=\left[\frac{2 i+1}{2 k}+\sigma, \frac{i+1}{k}\right] ;
\end{aligned}
$$

thus the unit interval $I$ is divided into $m=4 k$ subintervals $I_{j}=\left[c_{j-1}, c_{j}\right]$. We have $\left|I_{4 i+1}\right|=\left|I_{4 i+4}\right|=(1-2 k \sigma) / 2 k$ and $\left|I_{4 i+2}\right|=\left|I_{4 i+3}\right|=\sigma$. Obviously the functions $g_{j}(x)$ satisfy the conditions (5.1) and we denote
by $\psi(k, \sigma, \tau, \rho ; x)$ the corresponding continuous solution of the equations (4.1).

It follows from Theorems 6.1 and 6.3 that $\psi(k, \sigma, \tau, \rho ; x)$ has nowhere a unilateral derivative finite or infinite for any integer $k$ and positive numbers $\sigma, \tau, \rho$ satisfying (7.1), since we have

$$
\gamma=\frac{1}{2 k}, \quad \beta=\max \left\{\frac{1}{\rho}, \frac{1}{\tau}\right\}
$$

and

$$
\left|h_{w_{1}} \circ \cdots \circ h_{w_{n}}(I)\right|=\frac{1}{(2 k)^{n}}-\frac{\sigma}{(2 k)^{n-1}}-\frac{\sigma \tau}{(2 k)^{n-2}}-\cdots-\sigma \tau^{n-1}>\frac{\left|E_{0}\right|}{(2 k)^{n}},
$$

for every finite word $\left(w_{1} \ldots w_{n}\right) \in \Xi_{n}$.
Since the Cantor set $E_{0}$ is a unique compact subset of $I$ satisfying

$$
E_{0}=h_{1}\left(E_{0}\right) \cup h_{4}\left(E_{0}\right) \cup h_{5}\left(E_{0}\right) \cup \cdots \cup h_{4 k}\left(E_{0}\right)
$$

and since the mapping $v$ maps $Y_{0}^{*}$ homeomorphically onto $v\left(Y_{0}^{*}\right)$, it follows that $\bar{Y}_{0}^{*}=E_{0}$. On the other hand, for every $x \in W+\bigcup_{n \geqslant 1} Y_{n}^{*}$, there exist $n=n(x)$ and $j=j(x)$ such that $x \in U_{n, j}$; thus $E_{0} \subset Y_{0}^{*}+e$. Therefore $\left|Y_{0}^{*}\right|=\left|E_{0}\right|$, since $e$ is countable. Let $\Omega_{n}$ be the set of finite words $\left(w_{1} \ldots w_{n}\right)$ of length $n$ such that $1 \leqslant w_{j} \leqslant 4 k$ for $1 \leqslant j \leqslant n$. Then for any $n \geqslant 0$, the set $Y_{n+1}^{*}$ is decomposed as follows:

$$
Y_{n+1}^{*}=\bigcup_{\substack{\left(w_{1} \ldots w_{n}\right) \in \Omega_{n} \\ j \in \Omega_{1}-\Xi_{1}}} h_{w_{1}} \circ \cdots \circ h_{w_{n}} \circ h_{j}\left(Y_{0}^{*}\right)
$$

On each interval $V_{1, j}^{ \pm}$, for any $\left(w_{1} \ldots w_{n}\right) \in \Omega_{n}$ and $j \in \Omega_{1}-\Xi_{1}$, the function $h_{w_{1}} \circ \cdots \circ h_{w_{n}} \circ h_{j}(x)$ is a linear contraction; more precisely we have

$$
\left|\frac{d}{d x}\left(h_{w_{1}} \circ \cdots \circ h_{w_{n}} \circ h_{j}\right)(x)\right|=\rho^{n+1-r(w)} \tau^{r(w)} \quad \text { for } \quad x \in \stackrel{\circ}{V}_{1, j}^{ \pm}
$$

where $r(w) \equiv r\left(w_{1}, \ldots, w_{n}\right)=\frac{1}{2} \sum_{j=1}^{n}\left(1+\eta_{w_{j}}\right)$. Since $Y_{0}^{*} \cap U_{1, j}=\phi$ for all $j$, we have

$$
\left|Y_{n+1}^{*}\right|=2 k\left|Y_{0}^{*}\right| \sum_{\left(w_{1} \ldots w_{n}\right) \in \Omega_{n}} \rho^{n+1-r(w)} \tau^{r(w)}=2 k \rho\left|E_{0}\right|(2 k(\rho+\tau))^{n}
$$

Therefore it follows that

$$
\left|Y^{*}\right|=\sum_{n=0}^{\infty}\left|Y_{n}^{*}\right|=\left|E_{0}\right|+2 k \rho\left|E_{0}\right| \sum_{n=0}^{\infty}(2 k(\rho+\tau))^{n}=\frac{1-2 k(\sigma+\tau)}{1-2 k(\rho+\tau)}
$$

For $N \geqslant 0$, let $Z_{N}$ be the set of points $x \in I$ at which $A_{n}(x) \equiv 2$ or $3(\bmod 4)$ for all $n \geqslant N$ and $A_{N-1}(x) \equiv 0$ or $1(\bmod 4)$. Put $Z=\bigcup_{n \geqslant 0} Z_{n}$. Obviously $Z \subset W_{0} \subset I-e$. Then it is easily seen that the set $Z_{0}$ is a compact subset of $I$ satisfying

$$
Z_{0}=h_{2}\left(Z_{0}\right) \cup h_{3}\left(Z_{0}\right) \cup h_{6}\left(Z_{0}\right) \cup \cdots \cup h_{4 k-1}\left(Z_{0}\right) ;
$$

therefore $Z_{0}=E_{1}$. For any $n \geqslant 0$, the set $Z_{n+1}$ is decomposed as follows :

$$
Z_{n+1}=\bigcup_{\substack{\left(w_{1} \ldots w_{n}\right) \in \Omega_{n} \\ j \in \Xi_{1}}} h_{w_{1}} \circ \cdots \circ h_{w_{n}} \circ h_{j}\left(Z_{0}\right)
$$

On each open interval $U_{1, j}$, for any $\left(w_{1} \ldots w_{n}\right) \in \Omega_{n}$ and $j \in \Xi_{1}$, the function $h_{w_{1}} \circ \cdots \circ h_{w_{n}} \circ h_{j}(x)$ is a linear contraction such that

$$
\left|\frac{d}{d x}\left(h_{w_{1}} \circ \cdots \circ h_{w_{n}} \circ h_{j}\right)(x)\right|=\rho^{n-r(w)} \tau^{1+r(w)} \quad \text { for } \quad x \in U_{1, j} .
$$

Since $Z_{0} \cap V_{1, j}^{ \pm}=\phi$ for all $j$, we have

$$
\left|Z_{n+1}\right|=2 k\left|Z_{0}\right| \sum_{\left(w_{1} \ldots w_{n}\right) \in \Omega_{n}} \rho^{n-r(w)} \tau^{1+r(w)}=2 k \tau\left|E_{1}\right|(2 k(\rho+\tau))^{n}
$$

therefore

$$
\begin{aligned}
& |Z|=\sum_{n=0}^{\infty}\left|Z_{n}\right|=\left|E_{1}\right|+2 k \tau\left|E_{1}\right| \sum_{n=0}^{\infty}(2 k(\rho+\tau))^{n}= \\
& \frac{2 k(\sigma-\rho)}{1-2 k(\rho+\tau)}=1-\left|Y^{*}\right| .
\end{aligned}
$$

Then it follows from Theorems 6.2 and 6.4 that

$$
|Z| \leqslant\left|W_{0}\right| \leqslant|\operatorname{Knot}(\psi)| \leqslant 1-\left|Y^{*}\right|=|Z| ;
$$

hence we obtain

$$
|\operatorname{Knot}(\psi)|=\frac{2 k(\sigma-\rho)}{1-2 k(\rho+\tau)}
$$

Thus if we take, for a fixed number $\alpha \in[0,1)$,

$$
\sigma_{0}=\frac{1+\alpha}{8 k}, \quad \tau_{0}=\frac{1}{4 k} \quad \text { and } \quad \rho_{0}=\frac{1}{8 k}
$$

then the function $\psi_{0}(x) \equiv \psi\left(k, \sigma_{0}, \tau_{0}, \rho_{0} ; x\right)$ satisfies $\left|\operatorname{Knot}\left(\psi_{0}\right)\right|=\alpha$ and Hölder's condition of order $\log (2 k) / \log (8 k)$ by Lemma 4.2, which obviously converges to 1 as $k$ tends to infinity. This completes the proof of Theorem 2.1.

Remark. - Besicovitch's function $B(x)$ illustrated in Figure 1 is precisely equal to the function $\psi(1,1 / 8,1 / 4,1 / 8 ; x)$; thus $B(x)$ satisfies Hölder's condition of order $1 / 3$.

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