# A BEURLING-TYPE THEOREM 

BY

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## § 1. Introduction and statement of the main result

In this paper we shall be concerned primarily with the linear topological space $A^{-\infty}$ whose elements are holomorphic functions

$$
f(z)=\sum_{0}^{\infty} a_{\nu} z^{\nu}
$$

in the unit disk $U=\{z:|z|<1\}$ satisfying
or equivalently,

$$
\begin{equation*}
|f(z)| \leqslant C_{f}(1-|z|)^{-n_{f}} \quad(z \in U) \tag{1.1}
\end{equation*}
$$

$$
\log ^{+}\left|a_{\nu}\right|=O(\log \nu) \quad(\nu \rightarrow \infty)
$$

$A^{-\infty}$ can be thought of as the union of Banach spaces $A^{-n}(n>0)$, the norm in each $A^{-n}$ being defined as follows:

$$
\begin{equation*}
\|f\|_{-n}=\sup _{z \in U}\left\{|f(z)|(1-|z|)^{n}\right\}<\infty \tag{1.2}
\end{equation*}
$$

The topology in $A^{-\infty}$ is introduced in a standard way [6]. Clearly, $A^{-\infty}$ is a topological algebra under pointwise multiplication. It is the smallest algebra containing the disk algebra $A\left({ }^{1}\right)$ and closed under differentiation.

The dual of $A^{-\infty}$ is the topological algebra $A^{\infty}$ whose elements are functions $F(z)$ holomorphic in $U$ and infinitely differentiable in $\bar{U}$ :

$$
\begin{equation*}
F(z)=\sum_{0}^{\infty} b_{\nu} z^{\nu} \quad\left(b=O\left(\nu^{-k}\right) \quad \forall k>0\right) \tag{1.3}
\end{equation*}
$$

The linear functionals in $A^{-\infty}$ are given by the formula

$$
\begin{equation*}
F(f)=\frac{1}{2 \pi i} \lim _{r \rightarrow 1-} \int_{\partial U} \bar{F}(\zeta) f(r \zeta) \frac{d \zeta}{\zeta}=\sum_{0}^{\infty} b_{\nu} a_{\nu} \tag{1.4}
\end{equation*}
$$

[^0]Let $T$ denote the (continuous) linear operator in $A^{-\infty}$ of multiplication by the argument:

$$
\begin{equation*}
(T f)(z)=z f(z) \quad\left(f \in A^{-\infty}\right) \tag{1.5}
\end{equation*}
$$

Since the set of all polynomials is dense in $A^{-\infty}$ it is readily seen that every invariant subspace for the operator $T$ is a closed ideal in the algebra $A^{-\infty}$, and vice versa.

For every element $0 \neq f \in A^{-\infty}$ let $Z_{f}$ denote the sequence $\left\{\alpha_{\nu}\right\}$ of its zeros, $0 \leqslant\left|\alpha_{1}\right| \leqslant\left|\alpha_{2}\right| \leqslant \ldots<1$, each zero repeated according to its multiplicity. $Z_{f}$ will be called the zero set of $f$. For every closed ideal $0 \neq I \subset A^{-\infty}$ the zero set $Z_{I}$ is defined as the sequence of common zeros for all elements $f \in I$, each zero repeated according to its minimal multiplicity. For a complete description of $A^{-\infty}$-zero sets see [6] where a certain condition ( $T$ ) was established and proved to be necessary and sufficient for a sequence $\alpha=\left\{\alpha_{\nu}\right\} \subset U$ to be an $A^{-\infty}$-zero set. This condition ( $T$ ) implies in particular that ever subset of an $A^{-\infty}$-zero set is an $A^{-\infty}$-zero set itself. Therefore sets $Z_{I}$ are in fact not different from $Z_{j}$ : for every ideal $0 \neq I \subset A^{-\infty}$ there is an element $f \in A^{-\infty}$ such that $Z_{f}=Z_{I}$.

Our main result (Theorem 1.1) concerns the description of closed ideals $0 \neq I \subset A^{-\infty}$. It states that every such ideal is uniquely determined by its zero set $Z_{I}$ and by its so-called $x$-singular measure $\sigma_{I}$. Now, the notion of a $x$-singular measure $\sigma_{f}$ associated with functions $f(z)$ of the class $A^{-\infty}$ (and, for that matter, with those of the larger class $n=A^{-\infty} / A^{-\infty}$ ) was introduced in [6], but the definition adopted there depended heavily on a number of other concepts, in particular on that of a premeasure of bounded $\kappa$-variation. There is, however, an alternative definition for $\sigma_{f}\left(f \in A^{-\infty}\right)$ which is (at least formally) quite independent of the theory expounded in [6]. We shall use that definition to state our main result (Theorem 1.1) but we do not see how to prove it without making extensive use of the results from [6].

In this section we confine ourselves only to those preliminary notions and porpositions which are indispensible for introducing the concept of a $x$-singular measure and for formulating Theorem 1.l.

Definition 1.1. A subset $F$ of the circumference $\partial U$ is called a Beurling-Carleson (B.-C.) set if
(i) $F$ is closed;
(ii) $F$ is of Lebesgue measure zero, $|F|=0$;
(iii) $\sum_{v}\left|I_{\nu}\right| \log \frac{1}{\left|I_{\nu}\right|}<\infty$,
where $\left\{I_{\nu}\right\}$ are the complementary aros of $F$ (i.e. the components of $\partial U \backslash F$ ) and $\left|I_{\nu}\right|$ is the length of $I_{\nu}$.

It is well known $[2 ; 3]$ that $\mathrm{B} .-\mathrm{C}$. sets coincide with null sets for the classes $A^{n}=$ $\left\{f: f^{(n)} \in A\right\} \quad(n=1,2, \ldots)$. Moreover $[7 ; 8]$, if $F$ is a $B$. -C. set, then an outer function $\Phi(z) \in A^{\infty}$ exists such that $F=\left\{\zeta \in \partial U: \Phi^{(n)}(\zeta)=0 \forall n \geqslant 0\right\}$.

The set of all B.-C. sets will be denoted $\mathcal{F}$, and the set of all Borel sets $B$ such that $\bar{B} \in \mathcal{F}$ will be denoted $\mathcal{B}$.

Definition 1.2. A function $\sigma: \mathcal{B} \rightarrow \mathbf{R}$ is called a $x$-singular measure ( $x-$ s.m.) if
(i) $\sigma$ is a finite Borel measure on every B.-C. set $F \subset \partial U$;
(ii) there is a constant $C>0$ such that

$$
\begin{equation*}
|\sigma(F)| \leqslant C \sum_{v}\left|I_{v}\right|\left(\log \frac{2 \pi}{\left|I_{v}\right|}+1\right) \quad(\forall F \in \mathcal{F}), \tag{1.6}
\end{equation*}
$$

where $\left\{I_{\nu}\right\}$ are the conplementary arcs of $F$.
It is clear that a $x$-singular measure $\sigma$ is completely determined by the values $\sigma(F)(F \in \mathcal{F})$; in other words, a function $\sigma: \mathfrak{Z} \rightarrow \mathbf{R}$ possesses (if at all) only one extension to a $x$-s.m.

The total variation $|\sigma|$ of a $x$-s.m. $\sigma$ satisfying (1.6) is a non-negative $x$-s.m. with the constant not exceeding 2 C .

Notations like $\max \left\{\sigma_{1}, \sigma_{2}\right\}, \min \left\{\sigma_{1}, \sigma_{2}\right\}$, l.u.b. $\left\{\sigma_{\nu}\right\}, \sigma_{1} \geqslant \sigma_{2}$ have their familiar meaning accepted in the measure theory.

Proposition 1.1. Let $0 \neq f \in A^{-\infty}, F \in \mathcal{F}$. Let further $\Phi \in A^{\infty}$ be an outer function such that $F=\left\{\zeta \in \partial U: \Phi^{(n)}(\zeta)=0 \forall n \geqslant 0\right\}$ and $\mu$ be a non-negative Borel measure on $F$. Consider

$$
\begin{equation*}
f_{F, \mu}(z)=f(z) \Phi(z) \exp \left\{\int_{F} \zeta+z-z(|d \zeta|)\right\} \quad(z \in U) \tag{1.7}
\end{equation*}
$$

and define

$$
\begin{equation*}
m_{F, f}=\left\{\mu: f_{F, \mu} \in A^{-\infty}\right\} . \tag{1.8}
\end{equation*}
$$

Then
(i) $\prod_{F, f}$ does not depend on $\Phi$, i.e. for any given $F, \mu$ all functions (1.7) (with different $\Phi$ ) either belong to $A^{-\infty}$ or are outside $A^{-\infty}$;
(ii) $M_{F, f}$ has a maximal element $\mu_{0}$, so that

$$
m_{F, f}=\left\{\mu: 0 \leqslant \mu \leqslant \mu_{0}\right\} ;
$$

(iii) there is a constant $C$ such that

$$
\begin{equation*}
\mu_{0}(F) \leqslant C \sum_{\nu}\left|I_{\nu}\right|\left(\log \frac{2 \pi}{\left|I_{\nu}\right|}+1\right) \quad(\forall F \in \mathcal{F}) . \tag{1.9}
\end{equation*}
$$

This proposition will be proved in section 4 in the course of proving Theorem l.l.

Definition 1.3. With every element $0 \neq f \in A^{-\infty}$ a non-positive $x$-singular measure $\sigma_{f}$ will be associated defined as follows:

$$
\begin{equation*}
\sigma_{f}(F)=-\mu_{0}(F)=-\max _{\mu \in m_{F, f}} \mu(F) \quad(\forall F \in \mathcal{F}) . \tag{1.10}
\end{equation*}
$$

For $f=0$ we set formally $\sigma_{0}(F)=-\infty(\forall F \in \mathcal{F})$.
For every closed ideal $0 \neq I \in A^{-\infty}$ we define

$$
\begin{equation*}
\sigma_{I}=1 . \underset{f \in I}{ } \text { u.b. } \sigma_{f} \tag{1.11}
\end{equation*}
$$

In section 4 when proving Theorem 1.1 it will be shown that Definition 1.3 is equivalent to another definition of $\sigma_{f}$ as the $\varkappa$-singular part of a premeasure [6].

We are now in a position to formulate our main result.
Theorem 1.1. Let $I \neq\{0\}$ be a closed ideal in $A^{-\infty}$; let $Z_{I}$ and $\sigma_{I}$ be respectively its zero set and its $x$-singular measure. Then

$$
\begin{equation*}
I=\left\{f \in A^{-\infty}: Z, \supseteq Z_{1}, \sigma_{f} \leqslant \sigma_{I}\right\} . \tag{1.12}
\end{equation*}
$$

Conversely, let $\alpha=\left\{\alpha_{\nu}\right\}$ be an arbitrary $A^{-\infty}$-zero set and $\sigma_{0}$ be an arbitrary non-positive $x$-singular measure. Then

$$
\begin{equation*}
I\left(\alpha ; \sigma_{0}\right)=\left\{f \in A^{-\infty}: Z_{f} \supseteq \alpha, \sigma_{f} \leqslant \sigma_{0}\right\} \tag{1.13}
\end{equation*}
$$

is a non-trivial closed ideal in $A^{-\infty}$.
Corollary 1.1.1. The necassary and sufficient condition for an element $f \in A^{-\infty}$ to be cyclic ${ }^{(1)}$ is $Z_{f}=\varnothing, \sigma_{f}=0$.

Corallary 1.1.2. Every closed ideal in $A^{-\infty}$ is principal, i.e. generated by a single element.

Corallary 1.1.3. The only "maximal" ideals in $A^{-\infty}$ are those of the form $I_{z_{0}}=$ $\left\{f \in A^{-\infty}: f\left(z_{0}\right)=0\right\}\left(z_{0} \in U\right)$. A closed ideal I such that $Z_{I}=\varnothing, \sigma_{I} \neq 0$ is not contained in any maximal ideal.

In the succeeding pages we shall first (in section 2) carry out a more thorough study of $x$-singular measures and their relationship to premeasures of bounded $x$-variation [6]. In particular, we shall establish the following facts:
(a) Every $x$-singular measure is concentrated on a $x \boldsymbol{F}_{\sigma}$-set, i.e. on the union of a countable set of B.-C. sets.
(b) For every non-positive $x$-s.m. $\sigma$ there is an element $f \in A^{-\infty}$ such that $\sigma=\sigma_{f}$.
${ }^{(1)}$ i.e. for $/ A^{-\infty}$ to be dense in $A^{-\infty}$.

In section 3 we shall prove, using purely real-variable argument, a crucial approximation theorem for premeasures of bounded $\varkappa$-variation. Essentially, this theorem shows that in regard to some general measure-theoretical properties premeasure with a vanishing $x$-singular part comport themselves in some ways like absolutely continuous measures in the classical theory.

Finally, in section 4 we shall prove Theorem 1.1 using the above-mentioned approximation theorem, some standard functional-analytic argument involving the dual space $A^{\infty}$ and the notion of annihilator, and results from [6] concerning holomorphic and meromorphic functions of the class $n=A^{-\infty} / A^{-\infty}$ and their generalized Nevanlinna factorization. Incidentally we shall prove Proposition 1.1 and equivalence of the two definitions for $\sigma_{f}$.

We shall adhere throughout to the following notation: The letter $C$ will be used to denote various positive constants which may differ from one formula to the next. The complement of a set $S \subseteq \partial U$ will be denoted $S^{c}=\partial U \backslash S .|S|$ is always used to designate Lebesgue measure of a set $S \subseteq \partial U$.

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## § 2. Classes of harmonic functions and premeasures.

For the reader's convenience we shall give here (in a slightly modified form) some definitions and results from [6] which will be used later. These results center round the representation of harmonic functions by means of generalized Poisson integrals involving so-called premeasures. Once such a representation is established, the problem arises to describe the class of harmonic functions under consideration in terms of premeasures. For the class $\mathcal{H}$ (see below) of harmonic functions a downright isomorphism exists between $\mathcal{H}$ and the corresponding space of premeasures. It is clear that such a close relationship should make it possible to treat many problems concerning harmonic (and analytic) functions by purely real-variable means.

Note that, in the light of some recent results of W. K. Hayman and the present author [5], it is highly probable that a similar relationship exists for much wider classes of harmonic functions than $\mathcal{H}$.

Definition 2.1. A real-valued harmonic function $u(z)(z \in U), u(0)=0$, is called $x$-bounded above (or just $x$-bounded) if

$$
\begin{equation*}
-\infty<u(z) \leqslant C \log \frac{1}{1-|z|} \quad(z \in U) . \tag{2.1}
\end{equation*}
$$

The least constant $C$ in (2.1) will be called the upper $x$-bound of $u$ and will be denoted $\|u\|^{*}$.

Clearly $\|u\|^{*} \geqslant 0$, and $\|u\|^{*}=0$ implies $u(z) \equiv 0$. The class of all $x$-bounded harmonic functions will be denoted $\boldsymbol{\mathcal { H }}^{+}$.

Definition 2.2. $\boldsymbol{H}=\boldsymbol{H}^{+}-\boldsymbol{H}^{+}$, i.e. every $u(z) \in \mathcal{H}$ possesses a representation

$$
\begin{equation*}
u(z)=u_{1}(z)-u_{2}(z) \quad\left(u_{1}, u_{2} \in \mathcal{H}^{+}\right) \tag{2.2}
\end{equation*}
$$

Proposition 2.1. $\mathcal{H}$ becomes a Banach space if it is provided with the norm

$$
\begin{equation*}
\|u\|=\min \left(\left\|u_{1}\right\|^{*}+\left\|u_{2}\right\|^{*}\right) \tag{2.3}
\end{equation*}
$$

where minimum is taken over all the representations (2.2).
The proof is immediate, by the use of simple compactness theorems for harmonic functions. There is at least one minimal representation such that $\|u\|=\left\|u_{1}\right\|^{*}+\left\|u_{2}\right\|^{*}$.

Next we turn to the notion of a premeasure.
Definition 2.3. Let $\mathcal{K}$ be the set of all open, closed and halfclosed arcs of the circumference $\partial U$, including all one-point sets, $\partial U$ and $\varnothing$. A function $u: \mathcal{K} \rightarrow \mathbf{R}$ is called a premeasure if
(i) $\mu\left(I_{1} \cup I_{2}\right)=\mu\left(I_{1}\right)+\mu\left(I_{2}\right)$ for all $I_{1}, I_{2} \in \mathcal{K}$ such that

$$
I_{1} \cup I_{2} \in \mathcal{K}, I_{1} \cap I_{2}=\varnothing ;
$$

(ii) $\mu(\partial U)=0$;
(iii) $\lim _{\nu=\infty} \mu\left(I_{\nu}\right)=0$ whenever $I_{\nu} \in \mathcal{K}, I_{1} \supset I_{2} \supset \ldots, \bigcap_{\nu} I_{\nu}=\varnothing$.

Clearly, every premeasure is immediately extended by finite additivity to the class of sets

$$
S=\bigcup_{v=1}^{n} I_{v} \quad\left(I_{v} \in \mathcal{K}\right) .
$$

With every premeasure $\mu$ a function $\hat{\mu}(\theta)=\mu\left(I_{\theta}\right)(0<\theta \leqslant 2 \pi)$ will be associated, where $I_{\theta}=\left\{e^{i t}: 0 \leqslant t<\theta\right.$. Thus a $1-1$ correspondence is established between the set of premeasures and the set of all real functions $\hat{\mu}(\theta)(0<\theta \leqslant 2 \pi)$ satisfying the following conditions:
(a) $\hat{\mu}(\theta-)(0<\theta \leqslant 2 \pi)$ and $\hat{\mu}(\theta+)(0 \leqslant \theta<2 \pi)$ exist;
(b) $\hat{\mu}(\theta-)=\hat{\mu}(\theta)(0<\theta \leqslant 2 \pi)$;
(c) $\hat{\mu}(2 \pi)=0$.

Clearly, $\hat{\mu}(\theta)$ has at most a countable set of points of discontinuity, all of them jumps.

In what follows we shall adhere to the following notation: the distance between two points $\zeta_{1}, \zeta_{2} \in \partial U$ is

$$
d\left(\zeta_{1}, \zeta_{2}\right)=\frac{1}{\pi} \min \left\{\arg \frac{\zeta_{2}}{\zeta_{1}}, \arg \frac{\zeta_{1}}{\zeta_{2}}\right\} \quad(0 \leqslant \arg \zeta<2 \pi \quad \forall \zeta \in \partial U),
$$

so that $0 \leqslant d\left(\zeta_{1}, \zeta_{2}\right) \leqslant 1\left(\forall \zeta_{1}, \zeta_{2} \in \partial U\right)$; the distance of a point $\zeta \in \partial U$ from a set $F \subset \partial U$ is

$$
d(\zeta, F)=\inf _{\zeta^{\prime} \in F} d\left(\zeta, \zeta^{\prime}\right) .
$$

The $\delta$-neighbourhood of a set $S \subset \partial U$ is $S^{\delta}=\{\zeta \in \partial U: d(\zeta, S)<\delta\}$.
Definition 2.4. We shall assign to every open set $G \subset \partial U$ the quantity (which may be $+\infty$ )

$$
\begin{equation*}
x(G)=\frac{1}{2 \pi} \int_{G}\left|\log d\left(\zeta, G^{c}\right)\right| \cdot|d \zeta| \tag{2.4}
\end{equation*}
$$

Further, we define $\chi(\varnothing)=\chi(\partial U)=0$.
A straightforward computation hows that

$$
\begin{equation*}
x(G)=\sum_{v} \frac{\left|I_{v}\right|}{2 \pi}\left(\log \frac{2 \pi}{\left|I_{v}\right|}+1\right) \tag{2.5}
\end{equation*}
$$

$\left\{I_{\nu}\right\}$ being the set of components of $G \cdot\left({ }^{1}\right)$
Definition 2.5. The entropy $\hat{\chi}_{G}(F)$ of a closed set $F$ with respect to an open set $G \supset F$ is defined as

$$
\begin{equation*}
\hat{x}_{G}(F)=\frac{1}{2 \pi} \int_{G}\left|\log d\left(\zeta, F \cup G^{c}\right)\right| \cdot|d \zeta|-x(G)=\frac{1}{2 \pi} \int_{G}\left|\log \frac{d\left(\zeta, F \cup G^{c}\right)}{d\left(\zeta, \overline{G^{c}}\right)}\right| \cdot|d \zeta| . \tag{2.6}
\end{equation*}
$$

The entropy with respect to $\partial U$ will be called simply entropy and will be denoted $\hat{x}(F)$ :

$$
\hat{x}(F)=\frac{1}{2 \pi} \int_{\partial U}|\log d(\zeta, F)| \cdot|d \zeta| .
$$

We have $\hat{\varkappa}_{G}(\varnothing)=0$. If $|F|=0$, then $\hat{\varkappa}_{G}(F)=x(G \backslash F)$. According to Definition 1.1 the B.-C. sets are exactly those sets $F$ with $\hat{x}(F)<\infty$. From (2.6) follows easily that if $F_{1}, F_{2} \subset G$, then

$$
\begin{equation*}
\hat{x}_{G}\left(F_{1} \cup F_{2}\right) \leqslant \hat{x}_{G}\left(F_{1}\right)+\hat{x}_{G}\left(F_{2}\right) . \tag{2.7}
\end{equation*}
$$

${ }^{(1)}$ Sometimes we shall use notation (2.5) also for sets $G$, not necessarily open, consisting of a finite number of components $I_{\nu} \in \mathcal{K}$; we set

$$
\frac{|I|}{2 \pi}\left(\log \frac{2 \pi}{|I|}+1\right)=0 \text { if }|I|=0
$$

Definition 2.6. A premeasure $\mu$ (and the associated function $\hat{\mu}(\theta)$ ) is said to be $x$-boundeed above (or simply $x$-bounded) if for all open arcs $I \subset \partial U$

$$
\begin{equation*}
\mu(I) \leqslant C \varkappa(I)=\frac{C|I|}{2 \pi}\left(\log \frac{2 \pi}{|I|}+1\right) . \tag{2.8}
\end{equation*}
$$

The least constant $C$ in (2.8) will be called the upper $x$-bound of the premeasure $\mu$ and will be denoted $\|\mu\|^{+}$. The set of all $x$-bounded premeasures will be denoted $x B^{+}$.

Clearly $\|\mu\|^{+} \geqslant 0$, and $\|\mu\|^{+}=0$ if and only if $\mu=0$.
Definition 2.7. A premeasure $\mu$ (and the associated function $\hat{\mu}(\theta)$ ) is said to be of bounded $x$-variation if for every finite set $\left\{I_{\nu}\right\}$ of non-overlapping open arcs such that $U_{\nu} I_{\nu}=\partial U$

$$
\begin{equation*}
\sum_{\nu}\left|\mu\left(I_{\nu}\right)\right| \leqslant C \sum_{\nu} x\left(I_{\nu}\right)=C \sum_{\nu} \frac{\left|I_{\nu}\right|}{2 \pi}\left(\log \frac{2 \pi}{\left|I_{\nu}\right|}+1\right) \tag{2.9}
\end{equation*}
$$

The minimal constant $C$ in (2.9) will be called the $x$-variation of $\mu$ and will be denoted $\chi \operatorname{Var} \mu$. The set of all premeasures of bounded $x$-variation will be denoted $\chi \mathrm{V}$.

Proposition 2.2 [6]. Every $x$-bounded premeasure is of bounded $x$-variation and

$$
\begin{equation*}
\varkappa \operatorname{Var} \mu \leqslant 2\|\mu\|^{+} . \tag{2.10}
\end{equation*}
$$

Proposition 2.3. [6]. Every $\mu \in x V$ is the difference of two $x$-bounded premeasures $\mu=\mu_{1}-\mu_{2}$ with

$$
\begin{equation*}
\|\mu,\|^{+} \leqslant a \cdot \varkappa \operatorname{Var} \mu \quad(j=1,2) \tag{2.11}
\end{equation*}
$$

where $a$ is an absolute constant.

Proposition 2.4. $x V$ becomes a Banach space if provided with the norm

$$
\begin{equation*}
\|\mu\|=x \operatorname{Var} \mu \tag{2.12}
\end{equation*}
$$

The proof is immediate.
Next comes a theorem which, though not stated explicitely in [6], follows directly from the results of that paper.

Theorem 2.1. There exists a linear operator $u=\bar{D} \mu$ (the generalized Poisson operator) which maps $\varkappa V$ onto $\mathcal{H}$ :

$$
\begin{equation*}
u(z)=(D \mu)(z)=\int_{\partial U} P(\zeta, z) \mu(|d \zeta|) \quad(z \in U) \tag{2.13}
\end{equation*}
$$

where $P(\zeta, z)=\operatorname{Re}(\zeta+z) /(\zeta-z)(\zeta \in \partial U, z \in U)$ is the Poisson kernel and the integral is understood either as a Riemann-Stieltjes integral

$$
u(z)=\int_{0}^{2 \pi} P\left(e^{i \theta}, z\right) d \hat{\mu}(\theta)
$$

or as a Riemann integral

$$
\begin{equation*}
u(z)=-\int_{0}^{2 \pi} \hat{\mu}(\theta)\left[\frac{d}{d \theta} P\left(e^{i \theta}, z\right)\right] d \theta . \tag{2.13"}
\end{equation*}
$$

The inverse operator $\mu=\mathcal{D}^{-1} u$ is given by

$$
\begin{equation*}
\frac{1}{2}[\mu(I)+\mu(I)]=\lim _{r \rightarrow 1-} \frac{1}{2 \pi} \int_{I} u(r \zeta)|d \zeta| \tag{2.14}
\end{equation*}
$$

where $I \subset \partial U$ is an arbitrary open arc.
Corollary 2.1.1. There are two positive constants $\lambda_{1}$ and $\lambda_{2}$ such that

$$
\begin{gathered}
\lambda_{1} \cdot \varkappa \operatorname{Var} \mu \leqslant\|D \mu\|_{\mathcal{H}} \leqslant \lambda_{2} \cdot \varkappa \operatorname{Var} \mu \quad(\forall \mu \in \varkappa V), \\
\lambda_{1} \cdot\|\mu\|^{+} \leqslant\|D \mu\|^{*} \leqslant \lambda_{2}\|\mu\|^{+} \quad\left(\forall \mu \in k B^{+}\right) .
\end{gathered}
$$

Remark. The existence of the limit in (2.14) for $u \in \mathcal{H}^{+}$is the crucial point in the proof of Theorem 2.1. Recently W. K. Hayman and the present author proved [5] that the limit in (2.14) exists for every harmonic function $u(z)(z \in U), u(0)=0$, such that

$$
\int_{0}^{1} \sqrt{\frac{k(r)}{1-r}} d r<\infty
$$

where

$$
k(r)=\max _{|z|-r} u(z)
$$

Proposition 2.5 [6]. Let $\mu$ be a premeasure of bounded $x$-variation. Define for every B.-C. set $F$

$$
\begin{equation*}
\sigma(F)=-\sum_{V} \mu\left(I_{\nu}\right) \tag{2.15}
\end{equation*}
$$

$I_{\nu}$ being the complementary arcs of $F\left({ }^{(1)}\right.$ Then $\sigma$ possesses a unique extension to a $x$-singular measure. Moreover,

$$
\begin{equation*}
|\sigma(F)| \leqslant \chi \operatorname{Var} \mu \cdot \hat{x}(F) \quad(\forall F \in \mathcal{F}) \tag{2.16}
\end{equation*}
$$

If $\mu \in \varkappa B^{+}$, then $\sigma \leqslant 0$.
Definition 2.8. $\sigma$ will be called the $x$-singular part of the premeasure $\mu$.
We prove now a somewhat different form of (2.15) which will be needed later.
Proposition 2.6. Let $\mu$ be a premeasure of bounded $\varkappa$-variation and let $\sigma$ be its $x$-singular part. Then for every $F \in \mathcal{F}$

$$
\begin{equation*}
\sigma(F)=\lim _{\delta \rightarrow 0} \mu\left(F^{\delta}\right)=\lim _{\delta \rightarrow 0} \mu\left(\bar{F}^{\delta}\right) . \tag{2.17}
\end{equation*}
$$

${ }^{(1)}$ The series (2.15) is absolutely convergent (cf. [6]).

Proof. It is enough to prove the former equality. For every complementary interval $I_{\nu}$ of $F$ set

$$
I_{\nu \delta}=\left\{\zeta \in \partial U: d\left(\zeta, I_{v}^{c}\right) \geqslant \delta\right\}=\left(\left(I_{v}^{c}\right)^{\delta}\right)^{c} .
$$

Clearly,

$$
\mu\left(F^{\delta}\right)=-\sum_{\left|I_{\nu}\right| \geqslant 2 \delta} \mu\left(I_{\nu \delta}\right) .
$$

This together with (2.15) yields

$$
\mu\left(F^{\delta}\right)-\sigma(F)=\sum_{\left|I_{\nu}\right|<2 \delta} \mu\left(I_{\nu}\right)+\sum_{\left|I_{\nu}\right| \geqslant 2 \delta} \mu\left(I_{\nu} \backslash I_{\nu \delta}\right) .
$$

The first sum tends to zero as $\delta \rightarrow 0$; what remains to be proved is that

$$
\lim _{\delta \rightarrow 0} \sum_{I_{\nu} \geqslant 2 \delta} \mu\left(I_{\nu} \backslash I_{\nu \delta}\right)=0
$$

If we assume the contrary, then there is sequence $\delta_{n} \downarrow 0$ such that
(i) $\left|\mu\left(G_{n}\right)\right| \geqslant \varepsilon>0 \quad(n=1,2, \ldots)$,
where

$$
G_{n}=\bigcup_{2 \delta_{n}>|n v| \geqslant 2 \delta_{n+1}}\left(I_{v} \backslash I_{v \delta_{n+1}}\right) ;
$$

(ii) $x\left(G_{n}\right) \leqslant 2^{-n} \quad(n=1,2, \ldots)$.

Consider now

$$
G^{(N)}=\bigcup_{n=1}^{N} G_{n}
$$

$G^{\gamma(N)}$ is composed of a finite number of disjoint open arcs, say, $A_{\nu}^{(N)}$, and it is easily seen that for the complementary arcs $B_{v}^{(N)}$

$$
\sum_{\nu} \chi\left(B_{v}^{(N)}\right) \leqslant \hat{x}(F)<\infty .
$$

Therefore

$$
\sum_{\nu} x\left(A_{\nu}^{(N)}\right)+\sum_{v} \chi\left(B_{v}^{(N)}\right) \leqslant \hat{\chi}(F)+1 .
$$

On the other hand

$$
\sum_{v}\left|\mu\left(A_{\nu}^{(N)}\right)\right|+\sum_{v}\left|\mu\left(B_{v}^{(N)}\right)\right| \geqslant \sum_{n=1}^{N}\left|\mu\left(G_{n}\right)\right| \geqslant N \varepsilon \rightarrow \infty \quad(N \rightarrow \infty) .
$$

This clearly contradicts our assumption that $x \operatorname{Var} \mu<\infty$. Thus Proposition 2.6 is proved.
Our next task is to prove that every $x$-s.m. (cf. Definition 1.2) is concentrated on a $\chi F_{\sigma}$-set, i.e. on union of a countable set of B.-C. sets.

Theorem 2.2. Let $\sigma$ be a $x$-s.m. Then there is a sequence $\left\{F_{v}\right\}_{1}^{\infty}$ of B.-C. sets, $F_{1} \subseteq F_{2} \subseteq \ldots$, such that for every $F \in \mathcal{F}$

$$
\begin{equation*}
\sigma(F)=\lim _{\nu \rightarrow \infty} \sigma\left(F \cap F_{\nu}\right), \quad|\sigma|(F)=\lim _{v \rightarrow \infty}|\sigma|\left(F \cap F_{\nu}\right) . \tag{2.18}
\end{equation*}
$$

Proof. It is enough to prove the latter equality. Since $|\sigma|$ is a $x$-s.m., it satisfies (1.6) or, equivalently,

$$
\begin{equation*}
|\sigma|(F) \leqslant C \hat{\varkappa}(F) \quad(\forall F \in \mathcal{F}) . \tag{2.19}
\end{equation*}
$$

We shall prove the theorem by organizing a transfinite "process of exhaustion" and by showing that this process stops after a countable number of steps. With this goal in view we shall introduce certain parameters associated with a $x$-s.m.

Let $G \subset \partial U$ be an open set such that $\bar{G} \backslash G \in \mathcal{F}$, or equivalently $|\bar{G} \backslash G|=0$, $x(G)<\infty, x(\partial U \backslash \bar{G})<\infty$. Define

$$
\begin{equation*}
m(\sigma ; G)=\sup _{G \supset F \xi \mathcal{Z}}\left\{|\sigma|(F)-C \hat{\varkappa}_{G}(F)\right\}, \tag{2.20}
\end{equation*}
$$

where $C$ is the constant in (2.19). In view of (2.19) we have

$$
|\sigma|(F)+|\sigma|(\partial G) \leqslant C \hat{\varkappa}(F \cup \partial G)=C\left[\hat{\varkappa}_{G}\left(F^{\prime}\right)+\varkappa(G)+\varkappa(\partial U \backslash \bar{G})\right]
$$

for $F \subset G, F \in \mathcal{F}$, where $\partial G=\bar{G} \backslash G$. Therefore

$$
m(\sigma ; G) \leqslant C[\varkappa(G)+x(\partial U \backslash \bar{G})]-|\sigma|(\partial G)
$$

On the other hand, putting in $(2.20) F=\varnothing$ we get $m(\sigma ; G) \geqslant 0$. Thus

$$
\begin{equation*}
0 \leqslant m(\sigma ; G) \leqslant C \hat{\mathcal{\varkappa}}(\partial G)-|\sigma|(\partial G) . \tag{2.21}
\end{equation*}
$$

To proceed further with the proof we need three simple lemmas. But first introduce the following

Definition 2.9. An open set $G \subset \partial U$ will be called regular if $\hat{\mathcal{x}}(\partial G)<\infty$. The set of all regular sets $G$ will be denoted $\mathcal{G}$.

Lemma 2.2.1. Let $\left\{G_{\nu}\right\}_{1}^{\infty}$ be a sequence of regular sets. Let the following hypotheses hold:
(i) $\bar{G}_{\nu} \subset G_{\nu+1} \quad(\nu=1,2, \ldots) ;$
(ii) $G=\bigcup_{\nu=1}^{\infty} G_{\nu} \in \mathcal{G}$;
(iii) $\hat{\chi}_{G}\left(\partial G_{\nu}\right) \rightarrow 0 \quad(\nu \rightarrow \infty)$.

Then

$$
\begin{equation*}
m\left(\sigma ; G_{\nu}\right) \rightarrow m(\sigma ; G) \quad(\nu \rightarrow \infty) . \tag{2.22}
\end{equation*}
$$

Proof. Let $F_{\nu}=\partial G_{\nu}=\bar{G}_{\nu} \backslash G_{\nu}$. We have

$$
\begin{aligned}
m\left(\sigma ; G_{\nu}\right) & =\sup _{G_{\nu} \sim F G}\left\{|\sigma|(F)-C \hat{\varkappa}_{G_{\nu}}(F)\right\} \\
& =\sup _{G_{\nu} \supset F G \mathcal{F}}\left\{|\sigma|\left(F \cup F_{\nu}\right)-C \hat{\varkappa}_{G}\left(F \cup F_{\nu}\right)\right\}-|\sigma|\left(F_{\nu}\right)+C\left[\hat{\varkappa}_{G}\left(F \cup F_{\nu}\right)-\hat{\varkappa}_{G_{\nu}}(F)\right] \\
& \leqslant m(\sigma ; G)-|\sigma|\left(F_{\nu}\right)+C\left[\hat{\varkappa}_{G}\left(F \cup F_{\nu}\right)-\hat{\varkappa}_{G_{\nu}}(F)\right] \quad\left(\forall G_{\nu} \supset F \in \mathcal{F}\right),
\end{aligned}
$$

because the latter expression in brackets does not depend on $F$ :

$$
\begin{equation*}
\hat{\varkappa}_{G}\left(F \cup F_{\nu}\right)-\hat{\varkappa}_{G_{\nu}}(F)=\varkappa\left(G \backslash\left(F \cup F_{\nu}\right)\right)-\varkappa(G)-\varkappa\left(G_{\nu} \backslash F\right)+\varkappa\left(G_{\nu}\right)=\hat{\varkappa}_{G}\left(F_{\nu}\right), \tag{2.23}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\varlimsup_{\nu \rightarrow \infty} m\left(\sigma ; G_{\nu}\right) \leqslant m(\sigma ; G) . \tag{2.24}
\end{equation*}
$$

On the other hand, in view of (2.20) there is for every $\varepsilon>0$ an $F \in \mathcal{F}, F \subset G$, such that

$$
|\sigma|(F) \geqslant m(\sigma ; G)+C \hat{\varkappa}_{G}(F)-\varepsilon .
$$

If $\nu$ is large enough, then $F \subset G_{\nu}$ and

$$
\begin{equation*}
m\left(\sigma ; G_{\nu}\right) \geqslant|\sigma|(F)-C \hat{\varkappa}_{G_{\nu}}(F) \geqslant m(\sigma ; G)-\varepsilon+C\left[\hat{\varkappa}_{G}(F)-\hat{\varkappa}_{G_{\nu}}(F)\right] . \tag{2.25}
\end{equation*}
$$

Using (2.23) we find

$$
\begin{aligned}
\hat{\varkappa}_{G}(F)-\hat{\varkappa}_{G_{\nu}}(F) & =\left[\hat{\chi}_{G}\left(F \cup F_{\nu}\right)-\hat{\varkappa}_{G_{\nu}}(F)\right]-\left[\hat{\varkappa}_{G}\left(F \cup F_{\nu}\right)-\hat{\varkappa}_{G}(F)\right] \\
& =\hat{\varkappa}_{G}\left(F_{\nu}\right)+\hat{\varkappa}_{G}(F)-\hat{\varkappa}_{G}\left(F \cup F_{\nu}\right) \geqslant 0 ;
\end{aligned}
$$

therefore from (2.25) follows

$$
\lim _{\nu \rightarrow \infty} m\left(\sigma ; G_{\nu}\right) \geqslant m(\sigma ; G),
$$

which together with (2.24) yields (2.22).
Lemma 2.2.2

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} m\left(\sigma, F^{\delta}\right)=|\sigma|(F) \quad(\forall F \in \mathcal{F}) . \tag{2.26}
\end{equation*}
$$

Proof. Let $\left\{I_{\nu}\right\},\left|I_{1}\right| \geqslant\left|I_{2}\right| \geqslant \ldots$, be the complementary arcs of $F$. We have

$$
x\left(F^{\delta} \backslash F\right)=\frac{\delta}{\pi}\left(\log \frac{2 \pi}{\delta}+1\right) \sum_{\left|I_{v}\right| \geqslant 2 \delta} 1+\sum_{\left|I_{v}\right|<2 \delta} \frac{\left|I_{\nu}\right|}{2 \pi}\left(\log \frac{2 \pi}{\left|I_{v}\right|}+1\right) .
$$

Hence

$$
x\left(F \backslash{ }_{\delta} F\right) \rightarrow 0, \hat{x}_{F_{\delta}}(F)=x\left(F_{\delta} \backslash F\right)-\varkappa\left(F_{\delta}\right) \rightarrow 0 \quad(\delta \rightarrow 0),
$$

and

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} m\left(\sigma ; F_{\delta}\right) \geqslant|\sigma|(F) . \tag{2.27}
\end{equation*}
$$

Now we have to show that

$$
\begin{equation*}
\varlimsup_{\delta \rightarrow 0} m\left(\sigma ; F_{\delta}\right) \leqslant|\sigma|(F) . \tag{2.28}
\end{equation*}
$$

Assuming the contrary, a sequence $\delta_{1}>\delta_{2}>\ldots$ and the corresponding sequence $\left\{F_{\nu}\right\}$ of B.-C. sets, $F_{\nu} \subset F^{\delta_{\nu}}$, could be chosen so that

$$
|\sigma|\left(F_{\nu}\right) \geqslant C \hat{x}_{F} \delta_{\nu}\left(F_{v}\right)+|\sigma|(F)+a \quad(a>0 ; \nu=1,2, \ldots),
$$

$C$ being the constant in (2.19). This implies

$$
|\sigma|\left(F_{\nu} \backslash F\right) \geqslant C \hat{\varkappa}_{F} \delta_{\nu}\left(F_{\nu}\right)+a
$$

and consequently

$$
|\sigma|\left(F_{v} \backslash F^{\delta_{\nu^{\prime}}}\right) \geqslant C \hat{\varkappa}_{F_{\nu}} \delta_{\nu}\left(F_{\nu}\right)+\frac{a}{2} \geqslant C \hat{\varkappa}_{F} \delta_{\nu}\left(F_{\nu} \backslash F^{\delta_{v^{\prime}}}\right)+\frac{a}{2}
$$

for sufficiently large $\nu^{\prime}>v$. Taking $S_{\delta}=F_{v_{k}} \backslash F^{\delta_{\nu_{k+1}}}$ with a sufficiently sparse subsequence $\left\{v_{k}\right\}$ and a suitable $\delta_{0}>0$ we can therefore construct a sequence $\left\{S_{k}\right\}$ of disjoint B.-C. sets, all contained in some $F^{\delta_{0}}$, such that

$$
|\sigma|\left(S_{k}\right) \geqslant C \hat{\varkappa}_{F} \delta_{0}\left(S_{k}\right)+\frac{a}{4} .
$$

Hence,

$$
|\sigma|\left(\bigcup_{k=1}^{n} S_{k}\right) \geqslant C \sum_{k=1}^{n} \hat{x}_{F} \delta_{0}\left(S_{k}\right)+\frac{n a}{4} \geqslant C \hat{\varkappa}_{F} \delta_{0}\left(\bigcup_{k=1}^{n} S_{k}\right)+\frac{n a}{4}
$$

and therefore

$$
m\left(\sigma ; F^{\delta_{0}}\right)=\infty
$$

which contradicts (2.21). Thus our lemma is proved.
Lemma 2.2.3. Let $\mathcal{G}_{r}$ denote the set of all open sets $G \subset \partial U$ composed of a finite number of open arcs with rational end points: $G=\bigcup_{\nu=1}^{n} I_{\nu}, I_{\nu}=\left\{\zeta \in \partial U: \alpha_{\nu}<\arg \zeta<\beta_{\nu}\right\}, \alpha_{\nu}$ and $\beta_{\nu}$ rational. Then for every pair $\sigma_{1}, \sigma_{2}$ of $\chi$-singular measures such that $\left|\sigma_{1}\right|>\left|\sigma_{2}\right|$ there is at least one $G \in \mathcal{G}_{r}$ such that

$$
\begin{equation*}
m\left(\sigma_{1} ; G\right)>m\left(\sigma_{2} ; G\right) \tag{2.29}
\end{equation*}
$$

Proof. There is a $F \in \mathcal{F}$ such that $\left|\sigma_{1}\right|(F)>\left|\sigma_{2}\right|(F)$. Lemma 2.2.2 implies that there is a $\delta>0$ such that $m\left(\sigma_{1} ; F^{\delta}\right)>m\left(\sigma_{2} ; F^{\delta}\right)$. Moving slightly the end points of the components of $F^{\delta}$ we can, in view of Lemma 2.2 .1 , replace $F^{\delta}$ by a $G \in \mathcal{G}_{r}$ so that (2.29) should hold.

We are now in a position to complete the proof of Theorem 2.2. Let $\sigma=\sigma_{1} \neq 0$ be a $x$-s.m. Take a $F_{1} \in \mathcal{F}$ such that $|\sigma|\left(F_{1}\right)>0$ and define $\sigma_{2}(F)=\sigma_{1}(F)-\sigma_{1}\left(F \cap F_{1}\right)(\forall F \in \mathcal{F})$. Clearly, $\left|\sigma_{1}\right|>\left|\sigma_{2}\right|$.

We define now $\sigma_{\alpha}, F_{\alpha}$ by induction for all countable transfinite numbers $\alpha$. Assume the $\sigma_{\beta}$ and $F_{\beta}$ have already been defined for all $\beta<\alpha$. If $|\sigma|(F)=\sup _{\beta<\alpha}|\sigma|\left(F_{\beta} \cap F\right)$ $(\forall F \in \mathcal{F})$ set $\sigma_{\alpha}=0$ and $\sigma_{\gamma}=0$ for all $\gamma>\alpha$; if otherwise, take any $S \in \mathcal{F}$ such that $|\sigma|(S)-$ $\sup _{\beta<\alpha}|\sigma|\left(F_{\beta} \cap S\right)>0$ and set $F_{\alpha}=\left(\bigcup_{\beta<\alpha} F_{\beta}\right) \cup S, \sigma_{\alpha}(F)=\sigma(F)-\sigma\left(F \cap F_{\alpha}\right)$. We have thus constructed a decreasing transfinite system of $\varkappa$.s. measures $\left\{\left|\sigma_{\alpha}\right|\right\}$. Since $m\left(\sigma_{\alpha} ; G\right) \leqslant$ $m\left(\sigma_{\beta} ; G\right)$ for $\alpha>\beta, G \in \mathcal{G}_{r}$ and since $\mathcal{G}_{r}$ is countable, there must be a countable transfinite $\gamma$ such that $m\left(\sigma_{\gamma} ; G\right)=0\left(\forall G \in \mathcal{G}_{r}\right)$. Lemma 2.2.3 yields that $\sigma_{\gamma}=0$; therefore 19-772903 Acta mathematica 138. Imprimé le 30 Juin 1977

$$
|\sigma|(F)=\sup _{\beta<\gamma}|\sigma|\left(F_{\beta} \cap F\right) \quad(\forall F \in \mathcal{F}),
$$

which is equivalent to the assertion of Theorem 2.2 because the set

$$
F_{\sigma}=\bigcup_{\beta<\gamma} F_{\beta}
$$

is union of a countable set of B.-C. sets.
Theorem 2.3. Let $\sigma$ be a non-positive $\varkappa$-singular measure and let

$$
\begin{equation*}
0 \geqslant \sigma(F) \geqslant-C \hat{\varkappa}(F) \quad(\forall F \in \mathcal{F}) . \tag{2.30}
\end{equation*}
$$

Then there is a premeasure $\mu$ such that
(i) $\mu$ is $\varkappa$-bounded above and

$$
\begin{equation*}
\|\mu\| \leqslant a C \tag{2.31}
\end{equation*}
$$

a being an absolute constant;
(ii) the $\chi$-singular part of $\mu$ coincides with $\sigma$.

Proof. The proof is broken into a number of steps. First we consider the simplest case when $\sigma$ is concentrated on a finite set of points.

Lemma 2.3.1. Let $F_{0}=\left\{\zeta_{\nu}\right\}_{1}^{n} \subset \partial U, \sigma\left(\left\{\zeta_{\nu}\right\}\right)=-\sigma_{\nu} \leqslant 0(\nu=1,2, \ldots, n)$,

$$
\begin{equation*}
\sum_{\zeta_{\nu} \in F} \sigma_{\nu} \leqslant C \hat{\chi}(F) \quad\left(\forall F \subseteq F_{0}\right) . \tag{2.32}
\end{equation*}
$$

Then a non-negative piecewise constant function $p(\zeta)$ exists defined and continuous on $G=\partial U \backslash F_{0}$ and such that
(i) $\int_{\partial U} p(\zeta)|d \zeta|-\sum_{1}^{n} \sigma_{\nu}=0$;
(ii) $\mu(I)=\int_{I} p(\zeta)|d \zeta|-\sum_{\zeta_{\nu} \in I} \sigma_{v} \leqslant a C \varkappa(I)$
for all open arcs $I \subset \theta U$.
Proof. If $e^{i \theta_{1}}$ and $e^{i \theta_{2}}$ are the end points of $I$ then $\mu(I)$ is linear in $\theta_{i}(i=1,2)$ on every complementary arc of $F_{0}$. On the other hand, $x(I)$ being concave in $\theta_{i}(i=1,2)$, the inequality (2.34) has to be ensured only for those $I$ 's with the end points in $F_{0}$, because it will then hold for all other I's automatically.

Assume that the points $\zeta_{\nu}$ are arranged on $\partial U$ counterclockwise and let $I_{k l}(1 \leqslant k<$ $l \leqslant n+1)$ be the open are between $\zeta_{k}$ and $\zeta_{l}\left(\zeta_{n+1}=\zeta_{1}\right)$. Write those inequalities (2.34) which
correspond to the arcs $I_{k l}$ in the form

$$
\begin{equation*}
\sum_{v=k}^{l} p_{v}\left|I_{v, v+1}\right| \leqslant \sum_{v=k+1}^{l} \sigma_{v}+a C x\left(I_{k, l+1}\right) \quad(1 \leqslant k \leqslant l \leqslant n) ; \tag{2.35}
\end{equation*}
$$

(2.33) will then assume the form

$$
\begin{equation*}
\sum_{\nu=1}^{n} p_{v}\left|I_{v, v+1}\right|=\sum_{v=1}^{n} \sigma_{v} \tag{2.36}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
p_{\nu} \geqslant 0 \quad(\nu=1,2, \ldots, n) \tag{2.37}
\end{equation*}
$$

We shall show that the system composed of (2.35), (2.36) and (2.37) is consistent provided (2.32) holds and $a=1$. By a well-known compatibility criterion for inequalities we have to verify that for every finite system of arcs $\left\{I_{j}\right\}, I_{j}=I_{k j l j}$, and for corresponding positive numbers $\left\{\lambda_{j}\right\}$ such that

$$
\begin{equation*}
\sum_{j} \lambda_{j} \chi_{I_{j}}(\zeta) \geqslant 1 \quad\left(\zeta \in \partial U \backslash F_{0}\right) \tag{2.38}
\end{equation*}
$$

the following inequality holds:

$$
\begin{equation*}
\sum_{\nu=1}^{n} \sigma_{v} \leqslant \sum_{j} \lambda_{i} \sum_{\zeta_{\nu} \varepsilon_{j}} \sigma_{v}+C \sum_{j} \lambda_{j} x\left(I_{j}\right) \tag{2.39}
\end{equation*}
$$

$\chi_{I}(\zeta)$ being the characteristic function of an arc $I$. Clearly, we can confine ourselves to the case when all the $\lambda$, are rational. Moreover, replacing some of the arcs $I_{j}$, by shorter ones or discarding them altogether we can reduce (2.38) to an equality. Multiplying then (2.38) and (2.39) by the common denominator of the $\lambda$, and replacing $\left\{I_{j}\right\}$ by another system of arcs (with some arcs repeated several times, if necessary), we shall give the required result the following form: for any system of open arcs $\left\{I_{j}\right\}$ which have their end points in $F_{0}$, do not contain $\zeta_{1}$ and cover $F_{0}^{c}$ exactly $n$ times,

$$
\sum_{j} \chi_{I_{l}}(\zeta)=n \quad\left(\zeta \in F_{0}^{c}\right)
$$

the inequality holds

$$
\begin{equation*}
n \sum_{\nu=1}^{n} \sigma_{\nu} \leqslant \sum_{j} \sum_{\zeta_{\nu} \varepsilon_{I_{j}}} \sigma_{v}+C \sum_{j} x\left(I_{j}\right) \tag{2.40}
\end{equation*}
$$

or, equivalently,

$$
\begin{equation*}
\sum_{\nu=1}^{n}\left(n-\sum_{\zeta_{\nu} \in I_{j}} 1\right) \sigma_{v} \leqslant C \sum_{j} x\left(I_{j}\right) . \tag{2.41}
\end{equation*}
$$

For $n=1$ this is certainly true because (2.41) is then equivalent to (2.32), $F$ being the set of end points of all the arcs $I_{j}$. The general case is proved by induction which is made possible by the fact that every $n$-covering $\left\{I_{j}\right\}$ of $\partial U \backslash F_{0}$ (with $\zeta_{1}$ not covered at all) can be split up into $n$ simple coverings.

Thus we have proved the existence of a function $p(\zeta)$ which satisfies (2.33) and (2.34) for any $I$ not containing $\zeta_{1}(a=1)$. If $\zeta_{1} \in I$, write (2.34) for the two components of $I \backslash\left\{\zeta_{l}\right\}:$

$$
\mu\left(I_{1}\right) \leqslant C \varkappa\left(I_{1}\right), \quad \mu\left(I_{2}\right) \leqslant C \varkappa\left(I_{2}\right)
$$

Then

$$
\mu(I)=\mu\left(I_{1}\right)+\mu\left(I_{2}\right)-\sigma_{1} \leqslant C\left[\varkappa\left(I_{1}\right)+\varkappa\left(I_{2}\right)\right] \leqslant C\left[\varkappa(I)+\frac{|I| \log 2}{2 \pi}\right] \leqslant a C \varkappa(I)
$$

with $s=1+\log 2$. Thus (2.34) has been proved with $a=1+\log 2$.
Next in the proof of Theorem 2.3 comes the case when $\sigma$ is supported by a B.-C. set $F_{0}$. Let $I_{\nu}(\nu=1,2, \ldots)$ be the components of $F_{0}^{c}$ and let (2.30) hold for any closed $F \in F_{0}$. Consider the closed set

$$
S_{n}=\partial U \backslash \bigcup_{\nu=1}^{n} I_{\nu}=\bigcup_{\nu=1}^{n} J_{\nu}^{(n)} \quad(n \geqslant 1)
$$

which has exactly $n$ components $J_{v}^{(n)}(\nu=1,2, \ldots, n)$ that are either points or closed arcs. Choose in each $J_{\nu}^{(n)}$ one point $\zeta_{\nu}^{(n)} \in F_{0}$ and let $F^{(n)}=\left\{\zeta_{\nu}^{(n)}\right\}_{v=1}^{n}$. Place at each $\zeta_{\nu}^{(n)}$ the mass $-\sigma_{v}^{(n)}=\sigma\left(F_{0} \cap J_{v}^{(n)}\right)$ and apply Lemma 2.3.1. First check condition (2.32). For any subset. $F \subseteq F^{(n)}$ let $M_{F}$ denote the union of all those $J_{\nu}^{(n)}$ that have non-void intersection with $F$ Then

$$
\begin{equation*}
\sum_{\left.\zeta_{\nu}^{(n)}\right)_{F F}} \sigma_{\nu}^{(n)}=-\sigma\left(F_{0} \cap M_{F}\right) \leqslant C \hat{\varkappa}\left(F_{0} \cap M_{F}\right) \tag{2.42}
\end{equation*}
$$

On the other hand,

$$
\hat{x}\left(F_{0} \cap M_{F}\right) \leqslant \hat{\varkappa}(F)+\sum_{\nu=n+1}^{\infty} x\left(I_{\nu}\right)
$$

and since

$$
\lim _{n \rightarrow \infty} \sum_{v=n}^{\infty} x\left(I_{v}\right)=0
$$

we obtain

$$
\sum_{\zeta_{\nu}^{(n)} \in F} \sigma_{\nu}^{(n)} \leqslant(C+\varepsilon) \hat{\varkappa}(F) \quad\left(\forall F \subseteq F^{(n)}\right),
$$

$\varepsilon>0$ being arbitrarily small if $n$ is large enough. Now we can apply Lemma 2.3.1. and find a premeasure (in fact, a measure) $\mu^{(n)}$ with contant non-negative densities between the points $\zeta^{(n)}$ such that

$$
\mu^{(n)}\left(\left\{\zeta_{\nu}^{(n)}\right\}\right)=-\sigma_{v}^{(n)}=\sigma\left(F_{0} \cap J_{v}^{(n)}\right),\left\|\mu^{(n)}\right\|^{+} \leqslant(C+\varepsilon)(1+\log 2) .
$$

Using a Helly-type selection theorem [6] (or just a self-evident diagonal process) we can find a weakly convergent subsequence $\left\{\mu^{\left(n_{s}\right)}\right\}$ such that for every arc $I \subset \partial U$ whose end
points are outside $F_{0}$

$$
\lim _{s \rightarrow \infty} \mu^{\left(n_{s}\right)}(I)=\mu(I)
$$

where $\mu$ is a measure with constant non-negative density on each $I_{\nu}$. Clearly, $\|\mu\|^{+} \leqslant$ $C(1+\log 2)$. It remains to prove that $\mu_{\sigma}=\sigma$.

Let $F \subset F_{0}$ be a closed set. Using Proposition 2.6 we find that

$$
\begin{equation*}
\mu_{\sigma}(F)=\lim _{\delta \rightarrow 0} \mu\left(F^{\delta}\right) \tag{2.43}
\end{equation*}
$$

On the other hand,

$$
\mu\left(F^{\delta}\right)=\lim _{s \rightarrow \infty} \mu^{\left(n_{s}\right)}\left(F^{\delta}\right)
$$

if $\delta$ is such that the end points of the components of $F^{\delta}$ are outside $F_{0}$. We can now estimate $\mu^{(n)}\left(F^{\delta \delta}\right)$ as follows:

$$
-\sum_{\zeta_{\nu}^{(n)} \in F^{\delta}} \sigma_{\nu}^{(n)} \leqslant \mu^{(n)}\left(F^{\delta}\right) \leqslant-\sum_{\zeta_{\nu}^{(n)} \in F^{\delta}} \sigma_{\nu}^{(n)}+(C+\varepsilon)(1+\log 2)\left[\hat{\chi}_{F_{0}^{\delta}}^{\delta}\left(F_{0}\right)+\varkappa\left(F_{0}^{\delta}\right)\right] .
$$

Letting $n \rightarrow \infty$ we get

$$
\sigma\left(F_{0} \cap F^{\delta}\right) \leqslant \mu\left(F^{\delta}\right) \leqslant \sigma\left(F_{0} \cap F^{\delta}\right)+C(1+\log 2)\left[\hat{\varkappa}_{F_{0}^{\delta}}\left(F_{0}\right)+\chi\left(F_{0}^{\delta}\right)\right]
$$

Since $F_{0} \in \mathcal{F}$ we obtain

$$
\lim _{\delta \rightarrow 0} \hat{x}_{F_{0}^{\delta}}\left(F_{0}\right)=0, \quad \lim _{\delta \rightarrow 0} \chi\left(F_{0}^{\delta}\right)=0
$$

and hence, bearing in mind (2.43),

$$
\begin{equation*}
\lim _{\delta \rightarrow 0} \mu\left(F^{\delta}\right)=\mu_{\sigma}(F)=\sigma(F) \quad(\forall F \in \mathcal{F}) . \tag{2.44}
\end{equation*}
$$

To complete the proof of Theorem 2.3 we have to consider the general case when, according to Theorem 2.2, there is a sequence $F_{1} \subseteq F_{2} \subseteq \ldots$ of B.-C. sets such that $\sigma(F)=\lim _{n \rightarrow \infty} \sigma\left(F \cap F_{n}\right)(\forall F \in \mathcal{F})$. For every $F_{n}$ there is a (pre)measure $\mu^{(n)},\left\|\mu^{(n)}\right\|+\leqslant$ $C(1+\log 2)$, which has non-negative piecewise constant density on $F_{n}^{c}$ and whose singular part is

$$
\mu_{\sigma}^{(n)}(F)=\sigma_{n}(F)=\sigma\left(F \cap F_{n}\right) \quad(\forall F \in \mathcal{F}) .
$$

Using again the Helly-type selection theorem [6] we can extract a subsequence $\left\{\mu^{\left(n_{s}\right)}\right\}$ which converges weakly to a premeasure $\mu$. Repeating the same argument we used in proving (2.44) we shall arrive at the following conclusion:

$$
\mu_{\sigma}(F)=\lim _{n \rightarrow \infty} \sigma_{n}(F)=\sigma(F) \quad(\forall F \in \mathcal{F})
$$

Thus Theorem 2.3 has been proved.
We shall later need the following result which can be proved using the same technique:

Corollary 2.3.1. Let $\sigma_{1} \leqslant \sigma_{2} \leqslant \ldots \leqslant 0$ be a sequence of $x$-singular measures and let $\sigma_{1}$ (and consequently all the $\sigma_{\nu}$ ) satisfy condition of the type (2.30). If in addition

$$
\begin{equation*}
\lim _{\nu \rightarrow \infty} \sigma_{\nu}(F)=0 \quad(\forall F \in \mathcal{F}), \tag{2.45}
\end{equation*}
$$

then there is a sequence $\left\{\mu_{\nu}\right\}_{1}^{\infty}$ of premeasures such that
(i) $\left\|\mu_{\nu}\right\|^{+} \leqslant a C$;
(ii) the $x$-singular part of $\mu_{\nu}$ is equal to $\sigma_{\nu}$;
(iii) $\sup _{I \in \mathcal{X}}\left|\mu_{\nu}(I)\right| \rightarrow 0 \quad(\nu \rightarrow \infty)$.

## § 3. An approximation theorem for premeasures

Definition 3.1. A premeasure $\mu$ of bounded $x$-variation is said to be $x$-absolutely continuous below if there is a sequence $\left\{\mu_{\nu}\right\}_{1}^{\infty}$ of premeasures, $\mu_{\nu} \in \varkappa B^{+}$, such that
(i) $\mu+\mu_{\nu} \in \varkappa B^{+},\left\|\mu+\mu_{\nu}\right\|^{+} \leqslant C \quad(\forall \nu)$;
(ii) $\sup _{I \in \mathcal{X}}\left|\left(\mu+\mu_{\nu}\right)(I)\right| \rightarrow 0 \quad(\nu \rightarrow \infty)$.

Theorem 3.1. A premeasure $\mu \in_{\chi V} V$ is $x$-absolutely continuous below if and only if its $\chi$-singular part is non-negative:

$$
\begin{equation*}
\mu_{\sigma} \geqslant 0 \tag{3.3}
\end{equation*}
$$

Proof.
A. Necessity. Let $\mu$ be $x$-absolutely continuous below, i.e. let there be a sequence $\left\{\mu_{\nu}\right\}$ satisfying (3.1) and (3.2). Take an arbitrary set $F \in \mathcal{F}$ and let $\left\{I_{n}\right\}$ be its complementary arcs. We have

$$
\begin{aligned}
-\left(\mu+\mu_{\nu}\right)_{\sigma}(F) & =\sum_{n}\left(\mu+\mu_{\nu}\right)\left(I_{n}\right)=\sum_{n \leqslant N}\left(\mu+\mu_{\nu}\right)\left(I_{n}\right)+\sum_{n>N}\left(\mu+\mu_{\nu}\right)\left(I_{n}\right) \\
& \leqslant \sum_{n \leqslant N}\left(\mu+\mu_{\nu}\right)\left(I_{n}\right)+C \sum_{n>N} x\left(I_{n}\right) .
\end{aligned}
$$

Using (3.2) we get

$$
-\lim _{v \rightarrow \infty}\left(\mu+\mu_{v}\right)_{\sigma}(F) \leqslant C \sum_{n>N} x\left(I_{n}\right) \rightarrow 0 \quad(N \rightarrow \infty),
$$

because $\hat{\varkappa}(F)<\infty$. Thus

$$
\begin{equation*}
\lim _{\overrightarrow{v \rightarrow \infty}}\left(\mu+\mu_{\nu}\right)_{\sigma}\left(F^{\prime}\right) \geqslant 0 \tag{3.4}
\end{equation*}
$$

Since $\mu_{\nu} \in \varkappa B^{+}$its $\chi$-singular part is non-positive; therefore

$$
\begin{equation*}
\left(\mu+\mu_{\nu}\right)_{\sigma}(F) \leqslant \mu(F) \quad(\forall F \in \mathcal{F}) \tag{3.5}
\end{equation*}
$$

From (3.4) and (3.5) follows

$$
\mu_{\sigma}(F) \geqslant 0 \quad(\forall F \in \mathcal{F})
$$

which proves (3.3)
B. Sufficiency. Let $N \geqslant 1$ be entire. Consider the set $\mathcal{L}_{N}$ of half-open arc $I_{k l}=\left\{e^{i \theta}:(2 \pi k) / N \leqslant \theta<(2 \pi l) / N\right\}(0 \leqslant k<l \leqslant N)$; let $\mu\left(I_{k l}\right)=\mu_{k l}$. If $\mu$ is $\psi$-absolutely continuous below then (3.1) and (3.2) imply that the following system of inequalities and equations is consistent:

$$
\left.\begin{array}{l}
x_{k l} \leqslant M \varkappa\left(I_{k l}\right), \\
\mu_{k l}+x_{k l} \leqslant \min \left\{C \varkappa\left(I_{k l}\right), \varepsilon\right\},  \tag{3.6}\\
x_{k l}=\sum_{s=k}^{l-1} x_{s, s+1}, x_{0 N}=0 \quad(0 \leqslant k<l<N)
\end{array}\right\}
$$

for any $\varepsilon>0$ and some $M=M_{8}$. In fact, setting $x_{k l}=\mu_{\nu}\left(I_{k l}\right)$ and writing out all the requirements of Definition 3.1 regarding the intervals $I \in \mathcal{L}_{N}$ as well as all the additivity conditions and $\mu_{\nu}(\partial U)=0$, we obtain (3.6). Conversely, if for any $\varepsilon>0$ and for some $M=M_{\varepsilon}$ (3.6) has solutions for $N=1,2, \ldots$, then $\mu$ is $x$-absolutely continuous below. To prove this we have to form for every solution $\left\{x_{k l}\right\}$ of (3.6) a measure $x$ having constant density $x_{s, s+1}| | I_{s, s+1} \mid$ over every $I_{s, s+1}$. Using then the Helly-type selection theorem for premeasures [6] and effecting transition to the limit with $N \rightarrow \infty$ we shall obtain a premeasure $x$ which meets the following conditions:

$$
x(I) \leqslant M \varkappa(I) ; \quad \mu(I)+x(I) \leqslant \min \{C \nsim(I), \varepsilon\}
$$

for all open arcs $I \subset \partial U$ which do not contain the point $\zeta=1$, the last restriction being easily removed if $(1+\log 2) C, 2 \varepsilon$ is substituted for $C, \varepsilon$ respectively (cf. the proof of Lemma 2.3.1). Consequently, if $\mu$ is not $x$-absolutely continuous below then for every $C>0$ there is an $\varepsilon>0$ such that, however large $M$, (3.5) has no solutions for some $N$. Repeating the argument used in the proof of Lemma 2.3.1 we shall arrive at the conclusion that for such combination of $C, \varepsilon, M$ there is a covering of $\partial U$ by a finite system of disjoint half-closed arcs $\left\{I_{\nu}\right\}$ such that

$$
\sum_{\nu} \min \left\{\mu\left(I_{\nu}\right)+M \varkappa\left(I_{\nu}\right), C \varkappa\left(I_{\nu}\right), \varepsilon\right\}<0
$$

Let $\left\{I_{\nu}^{\prime}\right\}$ be those arcs among $\left\{I_{\nu}\right\}$ for which

$$
\min \left\{\mu\left(I_{\nu}\right)+M \varkappa\left(I_{\nu}\right), C \varkappa\left(I_{\nu}\right), \varepsilon\right\}=\mu\left(I_{\nu}\right)+M \varkappa\left(I_{\nu}\right),
$$

and let $\left\{I_{v}^{\prime \prime}\right\}=\left\{I_{v}\right\} \backslash\left\{I_{v}^{\prime}\right\}$. Clearly, $\mu\left(I_{v}^{\prime}\right)<0$. Setting $F_{M}=U_{v} I_{v}^{\prime}$ we find

$$
\begin{equation*}
\mu\left(F_{M}\right)<-(M-C) \varkappa\left(F_{M}\right)-C \varkappa\left(F_{M}\right)-C \sum_{\left|L_{v}^{\prime}\right|<0} \varkappa\left(I_{v}^{\prime \prime}\right)-\varepsilon \sum_{\left|\left.\right|_{v}\right| \geqslant \delta} 1, \tag{3.7}
\end{equation*}
$$

where $\delta$ is defined by the equation.

$$
\frac{C \delta}{2 \pi}\left(\log \frac{2 \pi}{\delta}+1\right)=\varepsilon .
$$

Put now $C=2 \varkappa \operatorname{Var} \mu$ and let $M \rightarrow \infty$. Bearing in mind the definition of $x$-variation we easily arrive at the following conclusion:
(a) $\left\{I_{\nu}^{\prime \prime}:\left|I_{\nu}^{\prime \prime}\right| \geqslant \delta\right\} \neq \varnothing$ for $M>2 C$;
(b) $\sum x\left(I_{v}^{\prime \prime}\right)=O(\mathrm{l}) \quad(M \rightarrow \infty)$;
(c) $\quad \varkappa\left(F_{M}\right) \rightarrow 0 \quad(M \rightarrow \infty)$;
(d) $\mu\left(F_{M}\right) \leqslant-2 \chi \operatorname{Var} \mu\left[\varkappa\left(F_{M}\right)+\sum_{\left|I_{v}^{\prime \prime}\right|<\delta} \chi\left(I_{\nu}^{\prime \prime}\right)\right]-\varepsilon$.

We shall assume for convenience that $F_{M}$ is a closed set composed of a finite number of closed ares and that $\varkappa\left(F_{M}\right)$ stands for $\varkappa$ (int $F_{M}$ ); the parameter $M$ will be assumed to run through a sequence $M_{1}<M_{2}<\ldots, \lim M_{n}=\infty$. To simplify the notation we shall write $F_{n}$ for $F_{M_{n}}$. Our aim now is to extract a subsequence $\left\{F_{n_{v}}\right\}$ which will converge in some sense (to be specified) to a B.-C. set $F$, and to show using (3.8) that $\mu_{\sigma}$ cannot be nonnegative on $F$. For that we need

Lemma 3.1.1. Let $\left\{F_{n}\right\}$ be a sequence of sets, each one composed of a finite number of closed arcs. Let the following hypotheses hold $(n \rightarrow \infty)$ :
(i) $\left|F_{n}\right| \rightarrow 0$
(ii) $\quad x\left(F_{n}^{c}\right)=O(\mathrm{l})$.

Then there is a subsequence $\left\{F_{n_{\nu}}\right\}$ and a B.-C. set $F$ such that for every $\delta>0$ and some $N=N_{\delta}$
(a) $F_{n_{\nu}} \subset F^{\delta}$,
(b) $F \subset F_{n_{v}}^{\delta}$
for $\nu>N_{\delta}$.
Proof. Let $\left\{I_{k n}\right\}$ be the complementary arcs of $F_{n}$ arranged so that $\left|I_{1 n}\right| \geqslant\left|I_{2 n}\right| \geqslant \ldots$. We show first that $\left|I_{1 n}\right|$ are bounded away from 0 . In fact,

$$
x\left(F_{n}^{c}\right)=\sum_{k} \frac{\left|I_{k n}\right|}{2 \pi}\left(\log \frac{2 \pi}{\left|I_{k n}\right|}+1\right) \geqslant \frac{\left|F_{n}^{c}\right|}{2 \pi}\left(\log \frac{2 \pi}{\left|I_{1 n}\right|}+1\right)
$$

and therefore

$$
\begin{equation*}
\log \frac{2 \pi}{\left|I_{1 n}\right|}+1 \leqslant \frac{2 \pi x\left(F_{n}^{c}\right)}{\left|F_{n}^{c}\right|} . \tag{3.9}
\end{equation*}
$$

Since $\left|F_{n}^{c}\right| \rightarrow 2 \pi$ and $\varkappa\left(F_{n}^{c}\right)=O(1)(3.9)$ shows that $\left|I_{1 n}\right|$ is bounded away from 0 . We can therefore choose a subsequence

$$
\left\{F_{\nu_{n}}\right\}=\left\{F_{n}^{\prime}\right\}
$$

such that

$$
\begin{equation*}
I_{1 n} \rightarrow J_{1} \quad(n \rightarrow \infty), \tag{3.10}
\end{equation*}
$$

where $\left\{I_{k n}^{\prime}\right\}$ are the complementary arcs of $F_{n}^{\prime}, J_{1}$ is some open arc, $\left|J_{1}\right|>0$, and (3.10) means that the end points of $I_{1 n}^{\prime}$ tend to the corresponding end points of $J_{1}$. If $\left|J_{1}\right|=2 \pi$ then $\left\{F_{n}^{\prime}\right\}$ is the required subsequence and $F=J_{1}^{c}$. If $\left|J_{1}\right|<2 \pi$ then the same argument shows that

$$
\begin{equation*}
\log \frac{2 \pi}{\left|I_{2 n}^{\prime}\right|}+1 \leqslant \frac{2 \pi x\left(F_{n}^{\prime c}\right)}{\left|F_{n}^{\prime} c\right|-\left|I_{1 n}^{\prime}\right|} \tag{3.11}
\end{equation*}
$$

and since the denominator of the latter faction tends to $2 \pi-\left|J_{1}\right|>0$ the lengths $\left|I_{2 n}^{\prime}\right|$ must be bounded away from zero. Therefore a subsequence $\left\{F_{n}^{\prime \prime}\right\}=\left\{F_{v_{n}}^{\prime}\right\}$ exists such that $I_{2 n}^{\prime \prime} \rightarrow J_{2}$. Continuing this process we shall either arrive after a finite number of steps at a subsequence $\left\{F_{n}^{(s)}\right.$ such that

$$
I_{k n}^{(s)} \rightarrow J_{k} \quad(n \rightarrow \infty ; k=1,2, \ldots, s)
$$

and $\sum_{k=1}^{s}\left|J_{k}\right|=2 \pi$ in which case $\left\{F_{n}^{(s)}\right\}$ is the required subsequence and $F=\left(\bigcup_{k=1}^{s} J_{k}\right)^{c}$ is a finite set, or the number of steps is infinite. In the latter case

$$
\begin{equation*}
\sum_{k=1}^{\infty}\left|J_{k}\right|=2 \pi . \tag{3.12}
\end{equation*}
$$

In fact, just as (3.9) and (3.11) it is easily seen that

$$
\log \frac{2 \pi}{\left|J_{s}\right|}+1 \leqslant \frac{2 \pi A}{2 \pi-\sum_{k=l}^{s-1}\left|J_{k}\right|},
$$

where $A$ is the upper bound for $\chi\left(F_{n}^{c}\right)$, and that proves (3.12) since clearly $\left|J_{s}\right| \rightarrow 0(s \rightarrow \infty)$. Taking the diagonal subsequence $\left\{F_{n}^{(n)}\right\}_{n=1}^{\infty}$ we get the required result. Thus our lemma is proved.

Now we can continue the proof of Theorem 3.1. As (3.8) shows, the assumption that $\mu$ is not absolutely continuous below implies the existence of a sequence $\left\{F_{n}\right\}_{1}^{\infty}$ of sets, each $F_{n}$ being composed of a finite number of closed arcs, such that
(i) $x\left(F_{n}\right) \rightarrow 0$ and a fortiori $\left|F_{n}\right| \rightarrow 0(n \rightarrow \infty)$;
(ii) $x\left(F_{n}^{c}\right) \leqslant A<\infty \quad(n=1,2, \ldots)$;
(iii) $\mu\left(F_{n}\right) \leqslant-C\left[\varkappa\left(F_{n}\right)+\sum_{\left|I_{k n}\right|<\delta} \varkappa\left(I_{k n}\right)\right]-\varepsilon$,
where $C=2 \chi \operatorname{Var} \mu,\left\{I_{k n}\right\}$ are the complementary ares of $F_{n}$, and $\delta$ and $\varepsilon$ are some positive numbers. Using Lemma 3.1.1 we can form a subsequence $\left\{F_{n_{v}}\right\}$ converging to a B.-C. set $F$ in the sense that for every $\varrho>0 F^{\varrho}$ contains all but a finite number of $F_{n_{\nu}}$ and
is contained in all but a finite number of $F_{n_{v}}^{e}$. Assume for simplicity that $\left\{F_{n}\right\}$ already is such a subsequence. We claim that the $\chi$-singular part $\mu_{\sigma}$ of the premeasure $\mu$ cannot be non-negative on $\boldsymbol{F}$.

If the contrary is true, then $\mu_{\sigma}(S) \geqslant 0(\forall S \subseteq F, S \in \mathcal{F})$ and in particular $\mu_{\sigma}\left(S_{n}\right) \geqslant 0$ with $S_{n}=F_{n} \cap F$. Using Proposition 2.6 we find $\lim _{Q \rightarrow 0} \mu\left(S_{n}^{e}\right) \geqslant 0$.

Therefore we can replace in (3.13) $F_{n}$ by $F_{n} \backslash S^{e} n$ and choose $\varrho_{n}$ so small that (3.13) should still hold though perhaps with a smaller $\varepsilon$ and only for sufficiently large $n$. Thus a sequence of numbers $\varrho_{n} \downarrow 0$ can be chosen as well as a sequence of sets $\left\{F_{n}\right\}_{1}^{\infty}$ (each one composed of a finite number of closed arcs) such that

$$
F_{n} \subset F^{e_{n}} \backslash F^{Q_{n+1}}
$$

and

$$
\begin{equation*}
\mu\left(F_{n}\right) \leqslant-C\left[\varkappa\left(F_{n}\right)+\varkappa\left(G_{n}\right)\right]-\varepsilon, \tag{3.14}
\end{equation*}
$$

where $G_{n}=\left(F^{e_{n}} \backslash F^{e_{n+1}}\right) \backslash F_{n}$.
Let $\mathcal{J}_{n}, \mathcal{Z}_{n}$ and $\mathcal{K}_{n}$ denote the systems of ares $I$ of which $F_{n}, G_{n}$ and $F^{e_{n}}$ are composed respectively; let

$$
S_{n}=\left(\bigcup_{k=1}^{n} \mathcal{J}_{k}\right) \cup\left(\bigcup_{k=1}^{n} \mathcal{Z}_{k}\right) \cup \mathcal{K}_{n+1}
$$

Further let $\mathcal{J}_{0}$ be the system of arcs that form $\partial U \backslash F^{e_{1}}$. Summing (3.14) we get

$$
\begin{aligned}
\sum_{l \in \mathcal{J}_{0}}|\mu(I)|+\sum_{l \in \mathcal{S}_{n}}|\mu(I)| & \geqslant \sum_{\nu=1}^{n}\left|\mu\left(F_{v}\right)\right| \geqslant C\left[\sum_{\nu=1}^{n} x\left(F_{v}\right)+\sum_{v=1}^{n} x\left(G_{v}\right)\right]+n \varepsilon \\
& =C \sum_{l \in S_{n}} x(I)-C \sum_{l \in X_{n}+1} x(I)+n \varepsilon=C\left[\sum_{l \in S_{n} \cup J_{0}} x(I)-\sum_{J \in X_{n}+1} x(I)-\sum_{l \in \mathcal{J}_{0}} x(I)\right]+n \varepsilon
\end{aligned}
$$

Since

$$
\sum_{I \in x_{n+1}} x(I) \rightarrow 0 \quad(n \rightarrow \infty)
$$

we obtain (for large enough $n$ )

$$
\sum_{I \in S_{n} \cup J_{0}}|\mu(I)| \geqslant C \sum_{I \in S_{n} \cup J_{0}} x(I) .
$$

We have arrive therefore at a contradiction, because $S_{n} \cup \mathcal{J}_{0}$ is a system of non-overlapping arcs covering $\partial U$, and $C=2 x \operatorname{Var} \mu$. This contradiction completes the proof of Theorem 3.1.

Corollary 3.1.1. Let $\mu \in_{\chi} V$ and $\mu_{\sigma} \geqslant 0$. Then the function

$$
\begin{equation*}
f(z)=\exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta-z} \mu(|d \zeta|)\right\} \quad(z \in U) \tag{3.15}
\end{equation*}
$$

possesses the following properties:
(i) it is analytic in $U$ and belongs to the class $n=A^{-\infty} / A^{-\infty}$ (cf. section 4);
(ii) $f(0)=1$;
(iii) there is a sequence of functions $\left\{g_{\nu}(z)\right\}_{1}^{\infty}$ belonging to $A^{-\infty}$ such that $h_{p}(z)=f(z) g_{\nu}(z)$ belong to $A^{-N}$ with some $N>0$ and

$$
\begin{equation*}
\left\|\mathbf{l}-h_{\nu}\right\|_{-N} \rightarrow 0 \quad(\nu \rightarrow \infty) . \tag{3.16}
\end{equation*}
$$

Proof. $\mu$ is $x$-absolutely continuous below. Taking $\left\{\mu_{\nu}\right\}$ as in Definition 3.1 and defining

$$
\begin{equation*}
g_{\nu}(z)=\exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta+z} \mu_{\nu}(|d \zeta|)\right\} \quad(z \in U) \tag{3.17}
\end{equation*}
$$

we obtain the required sequence. In fact,

$$
\left|g_{v}(z)\right|=\exp \left\{\int_{\partial U} P(\zeta, z) \mu_{\nu}(|d \zeta|)\right\} \leqslant(1-|z|)^{-\lambda_{2}| | \mu_{\nu}\| \|^{+}}
$$

(see Corollary 2.1.1) so that $g_{\nu} \in A^{-\infty}$. For the same reason $f g_{\nu} \in A^{-\lambda_{2} C}$ where $C$ is the constant in (3.1). We have further

$$
f(z) g_{v}(z)=\exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta-z}\left(\mu+\mu_{v}\right)(|d \zeta|)\right\}=\exp \left\{-\int_{0}^{2 \pi}\left[\frac{d}{d \theta} \frac{\left(e^{1 \theta}+z\right)}{e^{i \theta}-z}\left[\hat{\mu}(\theta)+\hat{\mu}_{\nu}(\theta)\right] d \theta\right\},\right.
$$

and from (3.2) follows easily that $f(z) g_{\nu}(z) \rightarrow 1$ uniformly on compact sets $F \subset U$. Therefore (3.16) holds for any $N>\lambda_{2} C$.

## § 4. Proof of Theorem 1.1

Corollary 3.1.1 implies in particular that an element $f \in A^{-\infty}$ possessing representation (3.15) with $\mu \in \varkappa B^{+}, \mu_{\sigma}=0$, is cyclic, i.e. the closed ideal $I_{f}$ generated by $f$ is $A^{-\infty}$ itself. Clearly, this covers an important special case of Theorem 1.1, provided that equivalence of the two definitions of $\mathcal{F}$ can be proved (cf. section 1). For the reader's convenience we shall give here some results from [6] related to the representation of functions of the classes $A^{-\infty}, n$.

Proposition 4.1. [2]. Every function $f(z) \in A^{-\infty}, f(0) \neq 0$, possesses a unique representation in the form

$$
\begin{equation*}
f(z)=f(0) \tilde{B}_{\alpha}(z) \exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta-z} \mu(|d \zeta|)\right\} \tag{4.1}
\end{equation*}
$$

where $\tilde{B}_{\alpha}(z)$ is the "generalized Blaschke product" associated with $\alpha=\left\{\alpha_{\nu}\right\}=Z_{f}$ :

$$
\begin{equation*}
\tilde{B}_{\alpha}(z)=\prod_{\alpha_{\nu} \in \alpha} \frac{\alpha_{\nu}-z}{1-\bar{\alpha}_{\nu} z} \cdot \frac{\left|\alpha_{\nu}\right|}{\alpha_{\nu}} \exp \left\{\frac{\left(\alpha_{\nu}| | \alpha_{\nu} \mid\right)+z}{\left(\alpha_{\nu}| | \alpha_{\nu} \mid\right)-z} \cdot \log \frac{1}{\left|\alpha_{\nu}\right|}\right\} \tag{4.2}
\end{equation*}
$$

and $\mu \in \varkappa B^{+}$. Moreover,

$$
\begin{equation*}
\operatorname{Sup}\|\mu\|^{+}<\infty \quad\left(\forall f \in A^{-n},\|f\|_{-n} \leqslant C\right) \tag{4.3}
\end{equation*}
$$

for any $n>0, C>0$.
Remark. A corresponding result for the class $n$ holds as well with $\mu \in \varkappa V$ and the quotient of two generalized Blaschke products instead of $\tilde{B}_{\alpha}(z)$ in (4.1).

Let $f(z) \in A^{-n}, f(0) \neq 0, \alpha=\left\{\alpha_{\nu}\right\}=Z_{f}$. Define

Then $\tau_{f}$ is a non-positive $\kappa$-singular measure satisfying condition (1.6) with the constant $C=a n, a$ being an absolute constant. This result follows immediately from the description of $A^{-\infty}$-zero sets (cf condition ( $T_{n}$ ) and ( $T$ ) in [6]).

Definition 4.1. $\tau_{f}$ will be called the Blaschke $x$-singular measure associated with $f$.
Definition 4.2. [6]. Let $f(z) \in A^{-\infty}, f(0) \neq 0$, be represented in the form (4.1). Let $\mu_{\sigma}$ be the $x$-singular part of the premeasure $\mu$ and $\tau$ - be the Blaschke $x$-singular measure. Then

$$
\begin{equation*}
\sigma_{f}=\mu_{\sigma}-\tau_{f} \tag{4.4}
\end{equation*}
$$

will be called the $x$-singular measure associated with $f$. If $f 0)=f^{\prime}(0)=\ldots=f^{(k-1}(0)=0$, $f^{(k)}(0) \neq 0(k \geqslant 1)$ and $f_{1}(z)=z^{-k} f(z)$ then by definition

$$
\sigma_{f}=\sigma_{f_{1}}
$$

Clearly $\sigma_{s} \leqslant 0$ for all $\mathrm{f} \in A^{-\infty}$.
It will be shown later that Definition 4.2 is equivalent to Definition 1.3.
The notion of a $x$-singular measure $\sigma_{1}$ associated with an ideal $0 \neq I \subset A^{-\infty}$ is reduced to $\sigma_{f}$ by means of (1.11).

Definition 4.3. $C^{\infty 0}$ is the linear topological space of all infinitely differentiable functions $F(\zeta)$ on $\partial U$ :

$$
F(\zeta)=\sum_{-\infty}^{\infty} b_{\nu} \zeta^{\nu} \quad\left(b_{\nu}=O\left(|v|^{-k}\right) \forall k>0\right) .
$$

Definition 4.4. $C^{-\infty}$ is the linear topological space of all forma series

$$
f=\sum_{-\infty}^{\infty} a_{\nu} \zeta^{\nu}
$$

where

$$
a_{\nu}=O\left(|\nu|^{k}\right) \text { for some } k=k_{f}>0 .
$$

The spaces $A^{\infty}$ and $A^{-\infty}$ will be thought of as subspaces of $C^{\infty}$ and $C^{-\infty}$ respectively.
The multiplication of elements belonging to $C^{-\infty}$ will be understood as formal multiplication of the corresponding series, whenever this leads to meaningful formulas for the coefficients of the product.

## Proposition 4.2.

(i) $C^{\infty} C^{-\infty} \subseteq C^{-\infty}$;
(ii) $C^{\infty} C^{\infty} \subseteq C^{\infty}$;
(iii) $A^{-\infty} A^{-\infty} \subseteq A^{-\infty}$;
(iv) if $f \in C^{\infty}, g_{\nu} \in C^{-\infty} \quad(\nu=1,2, \ldots)$ and $g_{\nu} \rightarrow g(\nu \rightarrow \infty)$ in the topology of $C^{-\infty}$, then $f g_{\nu} \rightarrow f g$ in $C^{-\infty}$;
(v) if $g_{\nu} \rightarrow g$ and $h_{\nu} \rightarrow h(v \rightarrow \infty)$ in $A^{-\infty}$, then $g_{\nu} h_{\nu} \rightarrow g h$ in $A^{-\infty}$.

The proof is obvious.
Definition 4.5. The annihilator of a closed ideal $I \subseteq A^{-\infty}$ is the subspace $\mathcal{A}_{I}$ of $C^{\infty}$ whose elements $F$ satisfy

Let

$$
\begin{equation*}
F f \in A^{-\infty} \quad(\forall f \in I) \tag{4.6}
\end{equation*}
$$

$$
F_{0}(\zeta)=\sum_{-\infty}^{\infty} b_{v} \zeta^{v}
$$

be some element of $\mathcal{A}_{1}$; then for any $f \in I$,

$$
f(z)=\sum_{0}^{\infty} a_{\nu} z^{n}
$$

(4.6) yields

$$
\begin{equation*}
\sum_{\nu=0}^{\infty} b_{-k-\nu} a_{\nu}=0 \quad(k=1,2, \ldots) \tag{4.7}
\end{equation*}
$$

This shows that $F_{0} \in \mathcal{A}_{I}$ implies $F_{0}+A^{\infty} \subseteq \mathcal{A}_{I}$; in particular, $A^{\infty} \subseteq \mathcal{A}_{I}$. Thus what really matters in Definition 4.5 is the non-analytic part of $F(\zeta)$, i.e. the coefficients $\left\{b_{\nu}\right\}_{-1}^{-\infty}$. It is easily seen that the quotient space $\mathcal{A}_{I} / A^{\infty}$ is isomorphic to the subspace $\mathcal{A}_{l}^{*}$ of $A^{\infty}$ consisting of those functionals $F^{*}$ for which $F^{*}(f)=0(\forall f \in I)$ (see formula (1.4) for the definition of $A^{\infty}$ as the dual of $A^{-\infty}$.

Proposition 4.3. For each closed ideal $I \subseteq A^{-\infty}$

$$
\begin{equation*}
I=\left\{j \in A^{-\infty}: F f \in A^{-\infty} \forall F \in \mathcal{A}_{I}\right\} . \tag{4.8}
\end{equation*}
$$

This is a direct consequence from Definition 4.5 and from the Hahn-Banach theorem for linear topological spaces (see, e.g., [4], chapter 2).

Now we prove some lemmas which will lead eventually to the proof of Theorem 1.1.
Lemma 4.1. Let $0 \neq F \in C^{\infty}$ and $f_{1}, f_{2}, g_{1}, g_{2} \in A^{-\infty}$. If $F f_{1}=g_{1}, F f_{2}=g_{2}$ then $f_{1} g_{2}=f_{2} g_{1}$.
Proof. $\left(F f_{1}\right) f_{2}=F\left(f_{1} f_{2}\right)=g_{1} f_{2} ; \quad\left(F f_{2}\right) f_{1}=F\left(f_{2} f_{1}\right)=g_{2} f_{1} ;$ therefore $g_{1} f_{2}=g_{2} f_{1}$. All the multiplications and transformations are easily justified.

Lemma 4.2. Let $0 \neq F \in C^{\infty}, f_{0} \in A^{-\infty}$ and $F f_{0} \in A^{-\infty}$. Then $F f \in A^{-\infty}$ whenever $f \in A^{-\infty}$ is such that $Z_{f} \supseteq Z_{f_{0}}, \sigma_{f} \leqslant \sigma_{f_{0}}$.

Proof. First take up the case $Z_{f}=Z_{f_{0}}$. Using Proposition 4.1 we can represent $f_{0}$ and $f$ in the form (4.1); then dividing $f_{0}$ by $f$ we obtain

$$
\begin{equation*}
\frac{f_{0}(z)}{f(z)}=\lambda \exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta-z} \mu(|d \zeta|)\right\} \tag{4.9}
\end{equation*}
$$

where $\mu \in_{\varkappa} V$ and $\mu_{\sigma}=\sigma_{f_{0}}-\sigma_{f} \geqslant 0$. Applying Corollary 3.1.1 we can find a sequence $\left\{g_{\nu}\right\}_{1}^{\infty}, g_{\nu} \in A^{-\infty}$, such that

$$
\frac{f_{0} g_{v}}{f} \in A^{-\infty}, \frac{f_{0} g_{v}}{f} \rightarrow 1 \quad(\nu \rightarrow \infty)
$$

in the topology of $A^{-\infty}$. Multiplying by $f$ we get $f_{0} g_{\nu} \rightarrow f(\nu \rightarrow \infty)$. Since by the hypothesis $F f_{0} \in A^{-\infty}$, we find $\left(F f_{0}\right) g_{\nu}=F\left(f_{0} g_{\nu}\right) \in A^{-\infty}$ and therefore, using Proposition 4.2, $F f=$ $\lim _{\nu \rightarrow \infty} F\left(f_{0} g_{\nu}\right) \in A^{-\infty}$.

If $Z_{f} \supset Z_{f_{0}}$ we can construct [6] a function $g \in A^{-\infty}$ such that $Z_{g}=Z_{f} \backslash \boldsymbol{Z}_{f_{0}}, \sigma_{g}=0$. Then $F\left(f_{0} g\right)=\left(F f_{0}\right) g \in A^{-\infty}, Z_{f_{0} g}=Z_{f}, \sigma_{f_{0} g}=\sigma_{f_{0}} \geqslant \sigma_{f}$, and the case $Z_{f} \supset Z_{f_{0}}$ is thus reduced to that already proved.

Lemma 4.3. Let (as in Proposition 1.1) $F \in \mathcal{F}, \sigma_{0}$ be a non-negative Borel measure on $F$ and $\Phi(z)(z \in \bar{U})$ be an outer function belonging to $A^{\infty}$ and vanishing on $F$ together with all its derivatives. Define

$$
\begin{gather*}
I(z)=\exp \left\{-\int_{\partial U} \frac{\zeta+z}{\varrho-z} \sigma_{0}(|d \zeta|)\right\} \quad(z \notin F),  \tag{4.10}\\
\Psi(\zeta)= \begin{cases}\Phi(\zeta) I^{-1}(\zeta) & (\zeta=\partial U \backslash F) \\
0 & (\zeta \in F),\end{cases} \tag{4.11}
\end{gather*}
$$

and

$$
\Psi_{1}(z)= \begin{cases}\Phi(z) I(z) & (z \in \bar{U} \backslash F)  \tag{4.12}\\ 0 & (z \in F)\end{cases}
$$

Then
(i) $\Psi \in C^{\infty}, \quad \Psi_{1} \in A^{\infty}$;
(ii) $\Psi^{\wedge} \Psi_{1}^{\prime}=\Phi^{2} \in A^{\infty}$;
(iii) an element $f \in A^{-\infty}$ has the property $\Psi f \in A^{-\infty}$ if and only if $\sigma_{f} \leqslant-\sigma_{0}$.

Proof. (i) Since $\Phi^{(n)}(z)=O\left[d^{N}(z, F)\right](z \in \bar{U})$ for any $n \geqslant 0, N \geqslant 0$ and $I^{(n)}(z)=O\left[d^{-2 n}(z, F)\right]$ $(z \in U \backslash F), \Psi(\tau)$ and $\Psi_{1}(\tau)$ are infinitely differentiable on $\partial U$. Note that $\Psi(z)=$ $\Phi(z) I^{-1}(z)(z \in U)$ does not belong to $A^{-\infty}$ barring the trivial case $\sigma_{0} \neq 0$.
(ii) Obvious.
(iii) Let $f \in A^{-\infty}$ and $\Psi f=g \in C^{-\infty}$. Multiplying by $\Psi_{1}$ we get $\Phi^{2} f=g \Psi_{1}$. If $g \in A^{-\infty}$ then equating the $x$-singular measures on both sides of the equation we find $\sigma_{f}=\sigma_{0}-\sigma_{0} \leqslant-\sigma_{0}$. Conversely, if $\sigma_{f} \leqslant-\sigma_{0}$ then applying Lemma 4.2 we infer from $\Psi \Psi_{1} \in A^{-\infty}$ that $\Psi f \in A^{-\infty}$ because $\Psi_{1}$ has no zeros and its $x$-singular measure is $-\sigma_{0}$.

Incidentally, Lemma 4.3 proves the equivalence of Definition 1.3 and Definition 4.2, as well as Proposition 1.1.

We are now in a position to complete the proof of Theorem l.1. First prove the second part of the Theorem. Let $\alpha=\left\{\alpha_{\nu}\right\}$ be an $A^{-\infty}$-zero set and let $\sigma_{0}$ be a non-positive $x$-singular measure. By Theorem 2.2 there is a sequence of B.-C. sets $F_{1} \subseteq F_{2} \subseteq \ldots$ and a sequence $\left\{\sigma_{\nu}\right\}_{1}^{\infty}$ of non-positive $\kappa$-singular measures, $\sigma_{\nu}$ being in fact the part of $\sigma_{0}$ supported by $F_{\nu}$, such that for any B.-C. set $F \sigma_{0}(F)=\lim _{\mu \rightarrow \infty} \sigma_{\nu}(F)=\lim _{\nu \rightarrow \infty} \sigma_{0}\left(F \cap F_{\nu}\right)$. Form as in Lemma 4.3 for all $F_{\nu}$ and the corresponding $\sigma_{\nu}$ the function

$$
\begin{equation*}
\Psi_{\nu}(\zeta)=\Phi_{\nu}(\zeta) \exp \left\{-\int_{\partial U} \zeta-z+z \sigma_{\nu}(|d \zeta|)\right\} \tag{4.13}
\end{equation*}
$$

$\Phi_{\nu}(z)$ being an outer function of the class $A^{\infty}$ with the null set $F_{\nu}$, and for every zero $\alpha_{\nu}$ let $\pi_{\nu}(\tau)=\left(\zeta-\alpha_{\nu}\right)^{k_{\nu}}, k_{\nu}$ being the multiplicity of that zero. Then Lemma 4.3 shows that $I\left(\alpha, \sigma_{0}\right)=\left\{f \in A^{-\infty}: \Psi_{\nu} f \in A^{-\infty}, \pi_{\nu} f \in A^{-\infty} \quad \forall \nu \geqslant 1\right\}$.

Therefore $I\left(\alpha, \sigma_{0}\right)$ is a closed ideal in $A^{-\infty}$. Its non-triviality follows from the fact that (by Theorem 2.3) there is a premeasure $\mu \in \varkappa B^{+}$with $\mu_{\sigma}=\sigma_{0}$, and by the results of [6] there is a function $f(z) \in A^{-\infty}$ with $Z_{f}=\alpha, \sigma_{f}=0$; therefore

$$
g(z)=f(z) \exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta-z} \mu(|d \zeta|)\right\}
$$

meets both conditions $Z_{g}=\alpha, \sigma_{g}=\sigma_{0}$.
Take up now the first part of Theorem 1.1. Let $I \neq 0$ be a closed ideal in $A^{-\infty}$ and $\mathcal{A}_{I}$ be its annihilator. Let $F \neq 0$ be a fixed element of $\mathcal{A}_{I}$ :

$$
\begin{equation*}
F f=g_{f} \in A^{-\infty} \quad(\forall f \in I) \tag{4.14}
\end{equation*}
$$

By Lemma 4.1 the function

$$
h(z)=\frac{g_{f}(z)}{f(z)}
$$

does not depend on the choice of $f \in I$. Since $h(z)$ belongs to the class $\eta=A^{-\infty} / A^{-\infty}$ it possesses [6] a unique representation in the form

$$
\begin{equation*}
h(z)=\lambda \frac{\tilde{B}_{\alpha^{\prime}}(z)}{\tilde{B}_{\beta^{\prime}}(z)} \exp \left\{\int_{\partial U} \frac{\zeta+z}{\zeta-z} \mu^{\prime}(|d \zeta|)\right\} \tag{4.15}
\end{equation*}
$$

where $\alpha^{\prime}=\left\{\alpha_{\nu}^{\prime}\right\}$ is the zero set and $\beta^{\prime}=\left\{\beta_{v}^{\prime}\right\}$ is the pole set of $h(z)$ and $\mu^{\prime} \lambda \in \varkappa V$. Therefore $\beta^{\prime} \subseteq Z_{I}, \mu_{\sigma}^{\prime} \leqslant-\sigma_{I}$, i.e. each zero of $f$ which is outside $Z_{I}$ must also be a zero of $g_{f}$, and the part of $\sigma_{f}$ which goes beyond $\sigma_{I}$ must also be a part of $\sigma_{g_{f}}$. Fix now a $f_{0} \in I$ and assume for simplicity that $f_{0}(0) \neq 0$. Let $Z_{f_{0}} \backslash Z_{I}=\left\{z_{\nu}\right\}, \alpha_{f_{0}}-\sigma_{I}=\sigma^{\prime} \leqslant 0$; let further $\sigma^{\prime}$ be concentrated on a set $S=\bigcup_{\nu} F_{\nu}, F_{1} \subseteq F_{2} \subseteq \ldots\left(F_{\nu} \in \mathfrak{F}\right)$ so that $\sigma^{\prime}=$ g.l.b. $\left\{\sigma_{\nu}\right\}$, where $\sigma_{\nu}(F)=\sigma^{\prime}\left(F \cap F_{\nu}\right)$ and $0 \geqslant \sigma_{1} \geqslant \sigma_{2} \geqslant \ldots \geqslant \sigma^{\prime}$. Multiply (4.14) by $\left[(\zeta-z)\left(\zeta-z_{2}\right) \ldots\left(\zeta-z_{\nu}\right)\right]^{-1} \Psi_{\nu}(\zeta)$ where $\Psi_{\nu} \in C^{\infty}$ has the form (4.13) and apply Lemma 4.2; then we arrive at the conclusion that $F f \in A^{-\infty}$ whenever

$$
Z_{f}=Z_{I} \cup\left\{z_{v}\right\}_{n+1}^{\infty}, \quad \sigma_{f}=\sigma_{I}+\left(\sigma^{\prime}-\sigma_{n}\right), \quad n=1,2, \ldots
$$

By use of Theorem 2.3, Corollary 2.3.1 and the technique developed in [6] for constructing function of the class $A^{-\infty}$ with given zero sets, we can form the following functions:
(a) $g \in A^{-\infty}$ such that $Z_{g}=Z_{I}, \sigma_{g}=\sigma_{I}$;
(b) $p_{n} \in A^{-\infty}$ such that $Z_{p_{n}}=\left\{z_{\nu}\right\}_{n+1}^{\infty}, \sigma_{p_{n}}=0$;
(c) $q_{n} \in A^{-\infty}$ such that $Z_{q_{n}}=\varnothing, \sigma_{q_{n}}=\sigma^{\prime}-\sigma_{n}$
and ensure that $p_{n} \rightarrow 1, q_{n} \rightarrow 1$ in the topology of $A^{-\infty}$. We have for all $n \geqslant 1$

$$
F g p_{n} q_{n} \in A^{-\infty} ;
$$

taking the limit when $n \rightarrow \infty$ and observing that $p_{n} q_{n} \rightarrow 1$ we obtain that $F g \in A^{-\infty}$ and therefore by Lemma 4.2

$$
\begin{equation*}
F f \in A^{-\infty} \quad\left(\forall f: Z_{f} \supseteq Z_{I}, \sigma_{f} \approx \sigma_{I}\right) . \tag{4.16}
\end{equation*}
$$

Since $F$ is an arbitrary element of $A_{I}$ this yields

$$
\left\{f \in A^{-\infty}: F f \in A^{-\infty} \forall F \in A_{i}\right\} \supseteq I\left(Z_{I}, \sigma_{I}\right) .
$$

Using (4.8) we find $I \supseteq I\left(Z_{I}, \sigma_{I}\right)$. On the other hand, if $f \in I$ then by the definition of $Z_{I}$ and $\sigma_{I}$ we have

$$
Z_{f} \supseteq Z_{I}, \quad \sigma_{f} \leqslant \sigma_{I}
$$

$$
I=I\left(Z_{I}, \sigma_{I}\right)
$$

and Theorem 1.1 has been proved.

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[^0]:    ${ }^{(1)} A$ is the algebra of all functions continuous in $O$ and analytic in $U$ with sup-norm and pointwise multiplication.
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