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Vladimir Azarin

## Growth Theory of

Subharmonic Functions

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## Growth Theory of Subharmonic Functions

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## Chapter 1

## Preface

This book aims to convert the noble art of constructing an entire function with prescribed asymptotic behavior to a handicraft.

For this you should just consider the limit set that describes the asymptotic behavior of the entire function, i.e., you should consider the set $U[\rho, \sigma]$ of subharmonic functions (that is, $\left\{v\right.$ is subharmonic : $\left.v\left(r e^{i \phi}\right) \leq \sigma r^{\rho}\right\}$ ) and pick out the subset $U$ which characterizes its asymptotic properties.

How to do it? The properties of limit sets are listed in Section 3. All the standard growth characteristics are expressed in terms of limit sets in Sections $3.2,3.3,5.7$. Examples of construction are to be found in Sections 5.4-6.3. So you can use this book as a reference book for construction of entire functions.

Of course, you need some terms. All the terms that we use in this book are listed on pages 249-253.

If you want to study the theory, I recommend that you solve the exercises that are in the text. Most of them are trivial. However, I recommend that you do all of them by the moment that they appear trivial to you.

A few words about the history of this book. It arose from a course of lectures that I gave at Kharkov University in 1977. After some time, under pressure and with active help of Prof. I.V. Ostrovskii, a rotaprint edition (Edition of KhGU) of this course appeared: the first part in 1978, the second one in 1982. Mathematical Reviews did not notice this fact.

Since that time lots of new and important results have been obtained. Some of them were presented in Chapter 3 of the review [GLO].

In 1994, when I started to work in the Bar-Ilan University and obtained a personal computer, my first wish was to study typing on it in English. This was the first impulse for translating this course into English (there are no more than five copies of this book in the world, I believe, one of them being mine). I continued this project while working in Bar-Ilan (1994-2006) but there was not much time for this. And now I have finished.

## Acknowledgements

I am indebted to many people. I start from Prof. I.V. Ostrovskii, who supported this idea for many years, and Prof. A.A. Gol'dberg, who stimulated my mathematical activity all my life by his letters and conversations.

I am indebted to Prof's A. Eremenko and M. Sodin, who, not being my "aspirants," solved a lot of problems connected to limit sets, and also to Dr.'s V. Giner, L. Podoshev and E. Fainberg who worked with me to develop the theory.

I am indebted to Prof's L. Hörmander and R. Sigurdsson who have sent me the preprints of their papers that were not yet published. I am indebted to Prof. I.F. Krasichkov-Ternovskii, who explained to me many years ago the connection between the multiplicator problem and completeness of the exponent system in a convex domain.

I am indebted to Prof's. M.I. Kadec and V.P. Fonf for proving Theorem 4.1.5.2, which is rather far off my speciality.

I am indebted to my coauthors Prof's D. Drasin and P. Poggi-Corradini; I have exploited the results of our joint paper in Section 6.2.

Of course, I am indebted to my late teacher Prof. B.Ya. Levin, who taught me entire and subharmonic functions and gave me the first problems in this area. Actually, the theory of limit sets is a generalization of the theory of functions of completely regular growth.

I am also indebted very much to my grandson Sasha Sodin, who transformed "my English" into English.

## Chapter 2

## Auxiliary Information. Subharmonic Functions

### 2.1 Semicontinuous functions

2.1.1. Let $x \in \mathbb{R}^{m}$ be a point in an $m$-dimensional Euclidean space, $E$ a Borel set and $f(x)$ a function on $E$ such that $f(x) \neq \infty$.

Set

$$
\begin{equation*}
M(f, x, \varepsilon):=\sup \left\{f\left(x^{\prime}\right):\left|x-x^{\prime}\right|<\varepsilon, x^{\prime} \in E\right\} \tag{2.1.1.1}
\end{equation*}
$$

The function

$$
f^{*}(x):=\lim _{\varepsilon \rightarrow 0} M(f, x, \varepsilon)
$$

is called the upper semicontinuous regularization of the function $f(x)$.
In the case of a finite jump, the regularization "raises" the values of the function. However, there is no influence on $f^{*}(x)$, if $f(x)$ tends to $-\infty$ "continuously".
Proposition 2.1.1.1 (Regularization Properties) The following properties hold:

$$
\begin{equation*}
f(x) \leq f^{*}(x) ; \tag{rg1}
\end{equation*}
$$

$$
\begin{equation*}
(\alpha f)^{*}(x)=\alpha f^{*}(x) ; \tag{rg2}
\end{equation*}
$$

$$
\begin{equation*}
\left(f^{*}\right)^{*}(x)=f^{*}(x) ; \tag{rg3}
\end{equation*}
$$

$$
\left(f_{1}+f_{2}\right)^{*}(x) \leq f_{1}^{*}(x)+f_{2}^{*}(x)
$$

$$
\begin{align*}
\left(\max \left(f_{1}, f_{2}\right)\right)^{*}(x) & \leq \max \left(f_{1}^{*}, f_{2}^{*}\right)(x) ;  \tag{rg4}\\
\left(\min \left(f_{1}, f_{2}\right)\right)^{*}(x) & <\min \left(f_{1}^{*}, f_{2}^{*}\right)(x) .
\end{align*}
$$

$$
\left(\min \left(f_{1}, f_{2}\right)\right)^{*}(x) \leq \min \left(f_{1}^{*}, f_{2}^{*}\right)(x)
$$

These properties are obvious corollaries of the definition of $f^{*}(x)$.

Exercise 2.1.1.1 Prove them.
2.1.2 The function $f(x)$ is called upper semicontinuous at a point $x$ if $f^{*}(x)=$ $f(x)$. We denote the class of upper semicontinuous functions on $E$ by $C^{+}(E)$. The function $f(x)$ is called lower semicontinuous if $-f(x)$ is upper semicontinuous (notation $f \in C^{-}(E)$ ).

Examples of semicontinuous functions are given by
Proposition 2.1.2.1 (Semicontinuity of Characteristic Functions of Sets) Let $G \subset$ $\mathbb{R}^{m}$ be an open set. Then its characteristic function $\chi_{G}$ is lower semicontinuous in $\mathbb{R}^{m}$. Let $F$ be a closed set, then $\chi_{F}$ is upper semicontinuous.

The proof is obvious.
Exercise 2.1.2.1 Prove this.
Proposition 2.1.2.2 (Connection with Continuity) If $f \in C^{+} \cap C^{-}$, then $f$ is continuous.

The assertion follows from the equalities

$$
f^{*}(x)=\limsup _{\varepsilon \rightarrow 0}\left\{f\left(x^{\prime}\right):\left|x-x^{\prime}\right|<\varepsilon\right\} ;-(-f)^{*}(x)=\liminf _{\varepsilon \rightarrow 0}\left\{f\left(x^{\prime}\right):\left|x-x^{\prime}\right|<\varepsilon\right\} .
$$

Proposition 2.1.2.3 ( $C^{+}$-Properties) The following holds:
$\left(C^{+} 1\right) \quad f \in C^{+}(E) \Rightarrow \alpha f \in C^{+}(E)$, for $\alpha \geq 0$
$\left(C^{+} 2\right) \quad f_{1}, f_{2} \in C^{+} \Rightarrow f_{1}+f_{2}, \max \left(f_{1}, f_{2}\right), \min \left(f_{1}, f_{2}\right) \in C^{+}$.
These properties follow from the properties of regularization (Proposition 2.1.1.1).

Exercise 2.1.2.2 Prove them.
Let $G$ be an open set. Set $G_{A}:=\{x \in G: f(x)<A\}$.
Theorem 2.1.2.4 (First Criterion of Semicontinuity) One has $f \in C^{+}$if and only if $G_{A}$ is open for all $A \in \mathbb{R}$.

Proof. Let $f(x)=f^{*}(x), x \in G$. Then $\{f(x)<A\} \Longrightarrow\left\{f^{*}(x)<A\right\} \Longrightarrow$ $\{M(f, x, \varepsilon)<A\}$ for all sufficiently small $\varepsilon$. Thus the neighborhood of $x V_{\varepsilon, x}:=$ $\left\{x^{\prime}:\left|x-x^{\prime}\right|<\varepsilon\right\}$ is contained in $G_{A}$.

Conversely, since the set $G_{A}$ is open for $A=f\left(x_{0}\right)+\delta$, we have $f^{*}\left(x_{0}\right) \leq$ $f\left(x_{0}\right)+\delta$ for any $\delta>0$, hence for $\delta=0$. With property $(\operatorname{rg} 1)$ of Proposition 2.1.1.1 this gives $f^{*}\left(x_{0}\right)=f\left(x_{0}\right)$.

Let $F$ be a closed set. Set $F^{A}:=\{x \in F: f(x) \geq A\}$. An obvious corollary of the previous theorem is

Corollary 2.1.2.5 One has $f \in C^{+}$if and only if $F^{A}$ is closed for all $A$.
Exercise 2.1.2.3 Prove the corollary.
We denote compact sets by $K$. Set $M(f, K)=\sup \{f(x): x \in K\}$.
Theorem 2.1.2.6 (Weierstrass) Let $K \subset \mathbb{R}^{m}$ be a compact set and $f \in C^{+}(K)$. Then there exists $x_{0} \in K$ such that $f\left(x_{0}\right)=M(f, K)$.
I.e., $f$ attains its supremum on any compact set.

Proof. Set $K_{n}:=\{x \in K: f(x) \geq M(f, K)-1 / n\}$.
The $K_{n}$ are closed by Corollary 2.1.2.5, nonempty by definition of $M(f, K)$. Their intersection is nonempty and is equal to the set

$$
K_{\max }:=\{x \in K: f(x) \geq M(f, K)\}
$$

It means that there exists $x_{0}$ in $K$ such that $f\left(x_{0}\right) \geq M(f, K)$.
The opposite inequality holds for any $x$ in $K$.

## Exercise 2.1.2.4 Why?

The following theorem shows that the functional $M(f, K)$ is continuous with respect to monotonic convergence of semicontinuous functions.

Proposition 2.1.2.7 (Continuity from the right of $M(f, K)$ ) Let $f_{n} \in C^{+}(K), f_{n} \downarrow$ $f, n=1,2,3 \ldots$

Then $M\left(f_{n}, K\right) \downarrow M(f, K)$.
Proof. It is clear that $\lim _{n \rightarrow \infty} M\left(f_{n}, K\right):=M$ exists.
Set $K_{n}:=\left\{x \in K: f_{n}(x) \geq M\right\}$. The intersection of the closed nonempty sets $K_{n}$ is nonempty and has the following form: $\bigcap_{n} K_{n}=\{x: f(x) \geq M\}$. So $M(f, K) \geq M$.

The opposite inequality is obvious.

## Exercise 2.1.2.5 Why?

In the same way one proves
Proposition 2.1.2.8 (Commutativity of inf and $M(\cdot)$ ) Let

$$
\left\{f_{\alpha} \in C_{+}(K), \alpha \in(0 ; \infty)\right\}
$$

be an arbitrarily decreasing family of semicontinuous functions. Then

$$
\inf _{\alpha} M\left(f_{\alpha}, K\right)=M\left(\inf _{\alpha} f_{\alpha}, K\right)
$$

Exercise 2.1.2.6 Prove this proposition.
Theorem 2.1.2.9 (Second Criterion of Semicontinuity) $f \in C^{+}(K)$ iff there exists a sequence $f_{n}$ of continuous functions such that $f_{n} \downarrow f$.

Sufficiency. Let $f_{n} \in C^{+}(K), f_{n} \downarrow f$. Set $K_{n}^{A}:=\left\{x \in K: f_{n}(x) \geq A\right\}$. This is a sequence of nonempty closed sets. If the set $K^{A}:=\{x: f(x) \geq A\}$ is nonempty, then $K^{A}$ is closed because $\bigcap_{n} K_{n}^{A}=K^{A}$. Hence $f \in C^{+}(K)$ by Corollary 2.1.2.5.

Necessity. Set $f_{n}(x, y):=f(y)-n|x-y|$.
This sequence of functions has the following properties:
a) it decreases monotonically in $n$ and

$$
\lim _{n \rightarrow \infty} f_{n}(x, y)=\left\{\begin{array}{l}
f(x), \text { for } x=y \\
-\infty, \text { for } x \neq y
\end{array}\right.
$$

b) for any fixed $n$ the functions $f_{n}$ are continuous in $x$ uniformly with respect to $y$, because $\left|f_{n}(x, y)-f_{n}\left(x^{\prime}, y\right)\right| \leq n\left|x-x^{\prime}\right|$;
c) $f_{n}$ are upper semicontinuous in $y$.

Proposition 2.1.2.7 and c) imply that

$$
\lim _{n \rightarrow \infty} M_{y}\left(f_{n}(x, y), K\right)=M_{y}\left(\lim _{n \rightarrow \infty} f_{n}(x, y), K\right)
$$

b) implies that the functions $f_{n}(x):=M_{y}\left(f_{n}(x, y), K\right)$ are continuous, and a) implies that they decrease monotonically to $f(x)$.
2.1.3 We will consider a family of upper semicontinuous functions: $\left\{f_{t}: t \in T \subset\right.$ $(0, \infty)\}$. It is easy to prove

Proposition 2.1.3.1 $f_{t} \in C^{+} \Longrightarrow \inf _{t \in T} f_{t}(x) \in C^{+}$.
Exercise 2.1.3.1 Prove this proposition.
Set $f_{T}(x):=\sup _{t \in T} f_{t}(x)$. The function $f_{T}$ is not, generally speaking, upper semicontinuous even if $T$ is countable and $f_{t}$ are continuous. It is not possible to replace $\sup _{t \in T}$ in the definition of $f_{T}$ by $\sup _{t \in T_{0}}$, where $T_{0}$ is a countable set. However, the following theorem holds:

Theorem 2.1.3.2 (Choquet's Lemma) There exists a countable set $T_{0} \subset T$ such that

$$
\left(\sup _{t \in T_{0}} f_{t}\right)^{*}(x)=\left(\sup _{t \in T} f_{t}\right)^{*}(x) .
$$

Proof. Let $\left\{x_{n}\right\}$ be a countable set that is dense in $\mathbb{R}^{m}$ and $\varepsilon_{j} \downarrow 0$. Then the balls

$$
K_{n, j}:=\left\{x:\left|x-x_{n}\right|<\varepsilon_{j}\right\}
$$

cover every point $x \in \mathbb{R}^{m}$ infinitely many times.
Renumbering we obtain a sequence $\left\{K_{l}: l \in \mathbb{N}\right\}$. For any $l$ there exists, by definition of $\sup _{K_{l}}$, such a point $x_{0} \in K_{l}$ that

$$
\begin{equation*}
\sup _{K_{l}} f_{T}(x) \leq f_{T}\left(x_{0}\right)+1 / 2 l \tag{2.1.3.1}
\end{equation*}
$$

By definition of $\sup _{T}$ there exists $t_{l}$ such that

$$
f_{T}\left(x_{0}\right)<f_{t_{l}}\left(x_{0}\right)+1 / 2 l .
$$

Thus

$$
\begin{equation*}
f_{T}\left(x_{0}\right)<\sup \left\{f_{t_{l}}(x): x \in K_{l}\right\}+1 / 2 l . \tag{2.1.3.2}
\end{equation*}
$$

The inequalities (2.1.3.1) and (2.1.3.2) imply that for any $l$ there exists $t_{l}$ such that

$$
\begin{equation*}
\sup \left\{f_{T}(x): x \in K_{l}\right\} \leq \sup \left\{f_{t_{l}}(x): x \in K_{l}\right\}+1 / l \tag{2.1.3.3}
\end{equation*}
$$

Now set $T_{0}=\left\{t_{l}\right\}$. Evidently, $f_{T_{0}}(x) \leq f_{T}(x)$ and thus

$$
\begin{equation*}
f_{T_{0}}^{*}(x) \leq f_{T}^{*}(x) \tag{2.1.3.4}
\end{equation*}
$$

Let us prove the opposite inequality.
Let $x \in \mathbb{R}^{m}$. Choose a subsequence $\left\{K_{l_{j}}\right\}$ that tends to $x$. From (2.1.3.3) we obtain

$$
\begin{align*}
f_{T}^{*}(x) & \leq \limsup _{j \rightarrow \infty} \sup _{x^{\prime} \in K_{l_{j}}} f_{T}\left(x^{\prime}\right) \\
& \leq \limsup _{j \rightarrow \infty} \sup _{x^{\prime} \in K_{l_{j}}} f_{t_{l_{j}}}\left(x^{\prime}\right)  \tag{2.1.3.5}\\
& \leq \limsup _{j \rightarrow \infty} \sup _{x^{\prime} \in K_{l_{j}}} f_{T_{0}}\left(x^{\prime}\right)=f_{T_{0}}^{*}(x) .
\end{align*}
$$

(2.1.3.4) and (2.1.3.5) imply the assertion of the theorem.

### 2.2 Measures and integrals

2.2.1 Let $G$ be an open set in $\mathbb{R}^{m}$ and $\sigma(G)$ a $\sigma$-algebra of Borel sets containing all the compact sets $K \subset G$.

Let $\mu$ be a countably additive nonnegative function on $\sigma(G)$, which is finite on all $K \subset G$. We will call it a measure or a mass distribution.

Let $G_{0}(\mu)$ be the largest open set for which $\mu$ is zero. (It is the union of all the open sets $G^{\prime}$ such that $\mu\left(G^{\prime}\right)=0$.)

The set $\operatorname{supp} \mu:=G \backslash G_{0}(\mu)$ is called the support of $\mu$. It is closed in $G$.
We say that $\mu$ is concentrated on $E \in \sigma(G)$ if $\mu(G \backslash E)=0$.

Theorem 2.2.1.1 (Support) The support of a measure $\mu$ is the smallest closed set on which the measure $\mu$ is concentrated.

Exercise 2.2.1.1 Prove this.
A measure $\mu$ can be concentrated on a non-closed set $E$ and then $E \Subset \operatorname{supp} \mu$.
Example 2.2.1.1 Let $E$ be a countable set dense in $G$. Then supp $\mu=G$ and, of course, $E \neq G$.

The set of all measures on $G$ will be denoted by $\mathcal{M}(G)$.
The measure $\mu_{F}(E):=\mu(E \cap F)$ is called the restriction of $\mu$ onto $F \in \sigma(G)$. It is easy to see that $\mu_{F}$ is concentrated on $F$ and $\operatorname{supp} \mu \subset \bar{F}$.

A countably additive function $\nu$ on $\sigma(G)$ that is finite for all $K \subset G$ is called a charge. We consider only real charges.

Example 2.2.1.2 $\quad \nu:=\mu_{1}-\mu_{2}, \quad \mu_{1}, \mu_{2} \in \mathcal{M}(G)$.
The set of all charges will be denoted $\mathcal{M}^{d}$.
Theorem 2.2.1.2 (Jordan decomposition) Let $\nu \in \mathcal{M}^{d}(G)$. Then there exist two sets $G^{+}, G^{-}$such that
a) $G=G^{+} \cup G^{-}, G^{+} \cap G^{-}=\varnothing$;
b) $\nu(E) \geq 0$ for $E \subset G^{+} ; ~ \nu(E) \leq 0$ for $E \subset G^{-}$.

One can find the proof in [Ha, Ch. VI Sec. 29]
The measures $\nu^{+}:=\nu_{G^{+}}$and $\nu^{-}:=\nu_{G^{-}}$, where $\nu_{G^{+}}, \nu_{G^{-}}$are restrictions of $\nu$ to $G^{+}, G^{-}$, are called the positive and negative, respectively, variations of $\nu$. The measure $|\nu|:=\nu_{+}+\nu_{-}$is called the full variation of $\nu$ or just a variation.

Theorem 2.2.1.3 (Variations) The following holds:

$$
\nu^{+}(E)=\sup _{E^{\prime} \subset E} \nu\left(E^{\prime}\right) ; \nu^{-}(E)=\inf _{E^{\prime} \subset E} \nu\left(E^{\prime}\right) ; \nu=\nu^{+}+\nu^{-} .
$$

The proof is easy enough.
Exercise 2.2.1.2 Prove this.

Example 2.2.1.3 Let $\psi(x)$ be a locally summable function with respect to the Lebesgue measure. Set $\nu(E):=\int_{E} \psi(x) d x$. Then

$$
\nu^{+}(E)=\int_{E} \psi^{+}(x) d x, \nu^{-}(E)=\int_{E} \psi^{-}(x) d x ;|\nu|(E)=\int_{E}|\psi|(x) d x
$$

where

$$
\begin{equation*}
\psi^{+}(x)=\max (0, \psi(x)) ; \psi^{-}(x)=-\min (0, \psi(x)) \tag{2.2.1.1}
\end{equation*}
$$

2.2.2 The function $f(x), x \in G$ is called a Borel function if the set $E^{A}:=\{f(x)>$ $A\}$ belongs to $\sigma(G)$ for any $A \in \mathbb{R}$.

Let $K \Subset G$ be a compact set and $f$ a Borel function. Then the LebesgueStieltjes integrals of the form $\int_{K} f^{+} d \mu, \int_{K} f^{-} d \mu$ with respect to a measure $\mu \in$ $\mathcal{M}(G)$ are defined, and $\int_{K} f d \mu:=\int_{K} f^{+} d \mu-\int_{K} f^{-} d \mu$ is defined if at least one of the terms is finite.

We say that a property holds $\mu$-almost everywhere on $E$ if the set $E_{0}$ of $x$ for which it does not hold satisfies the condition $\mu\left(E_{0}\right)=0$.

We will denote all the compact sets in $G$ as $K$ (sometimes with indexes). The following theorems hold:

## Theorem 2.2.2.1 (Lebesgue) Let $\left\{f_{n}, n \in \mathbb{N}\right\}$ be a sequence of Borel functions

 on $K$ and $g(x) \geq 0$ a function on $K$ that is summable with respect to $\mu$ (i.e., its integral is finite), $\left|f_{n}(x)\right| \leq g(x)$ for $x$ in $K$, and $f_{n} \rightarrow f$ when $n \rightarrow \infty$.Then $\lim _{n \rightarrow \infty} \int_{K} f_{n} d \mu=\int_{K} f d \mu$.
Theorem 2.2.2.2 (B. Levy) Let $f_{n} \downarrow f$ when $n \rightarrow \infty$, and $f$ be a summable function on $K$.

Then $\lim _{n \rightarrow \infty} \int_{K} f_{n} d \mu=\int_{K} f d \mu$.
Theorem 2.2.2.3 (Fatou's Lemma) Let $f_{n}(x) \leq$ const $<\infty$ for $x$ in $K$.
Then $\limsup _{n \rightarrow \infty} \int_{K} f_{n} d \mu \leq \int_{K} \limsup _{n \rightarrow \infty} f_{n} d \mu$.
The proofs can be found in [Ha, Ch. V, Sec. 27].
Let $L(\mu)$ be the space of functions that are summable with respect to $\mu$. We say that $f_{n} \rightarrow f$ in $L(\mu)$ if $f_{n}, f \in L(\mu)$ and

$$
\left\|f_{n}-f\right\|:=\int\left|f_{n}-f\right|(x) d \mu \rightarrow 0
$$

Theorem 2.2.2.4 (Uniqueness in $L(\mu)$ ) Let $f_{n} \rightarrow f$ in $L(\mu)$ and

$$
\int f_{n} \psi d \mu \rightarrow \int g \psi d \mu
$$

for any $\psi$ continuous on $\operatorname{supp} \mu$. Then $\|g-f\|=0$.
For the proof see, e.g., [Hö, Th. 1.2.5].
2.2.3 Let $\phi(x)$ be a Borel function on $G$. The set supp $\phi:=\overline{\{x: \phi(x) \neq 0\}}$ is called the support of $\phi(x)$. A function $\phi$ is called finite in $G$ if supp $\phi \Subset G$.

We say that a sequence $\mu_{n} \in \mathcal{M}$ converges weakly to $\mu \in \mathcal{M}$ if the condition $\int \phi d \mu_{n} \rightarrow \int \phi d \mu$ holds for any continuous function $\phi$.

We will not show the integration domain, because it is always $\operatorname{supp} \phi$.
The weak (it is called also $C^{*}$-) convergence will be denoted as $\xrightarrow{*}$.

Theorem 2.2.3.1 ( $C^{*}$-limits) If $\mu_{n} \xrightarrow{*} \mu$, then for $E \in \sigma(G)$ the following assertions hold:

$$
\begin{aligned}
& \limsup _{n \rightarrow \infty} \mu_{n}(\bar{E}) \leq \mu(\bar{E}) \\
& \liminf _{n \rightarrow \infty} \mu_{n}(\stackrel{\circ}{E}) \geq \mu(\stackrel{\circ}{E})
\end{aligned}
$$

where $\stackrel{\circ}{E}$ is the interior of $E, \bar{E}$ is the closure of $E$.
Proof. Let $\chi_{\bar{E}}$ be the characteristic function of the set $\bar{E}$. It is upper semicontinuous. Thus there exists a decreasing sequence $\varphi_{m}$ of continuous functions finite in $G$ that converges to $\chi_{\bar{E}}$ as $m \rightarrow \infty$. Then we have

$$
\mu_{n}(\bar{E})=\int \chi_{\bar{E}} d \mu_{n} \leq \int \varphi_{m} d \mu_{n}
$$

Passing to the limit as $n \rightarrow \infty$ we obtain

$$
\limsup _{n \rightarrow \infty} \mu_{n}(\bar{E}) \leq \int \varphi_{m} d \mu
$$

Passing to the limit as $m \rightarrow \infty$ we obtain by Theorem 2.2.2.2

$$
\limsup _{n \rightarrow \infty} \mu_{n}(\bar{E}) \leq \int \chi_{\bar{E}} d \mu=\mu(\bar{E})
$$

The proof for $\stackrel{\circ}{E}$ is analogous.
Theorem 2.2.3.2 (Helly) Let $\left\{\mu_{\alpha}: \alpha \in A\right\}$ be a family of measures uniformly bounded on any compact set $K \subset G$, i.e., $\exists C=C(K): \mu_{\alpha}(K) \leq C(K)$, for $K \Subset G$. Then this family is weakly compact, i.e., there exists a sequence $\left\{\alpha_{j}: \alpha_{j} \in A\right\}$ and a measure $\mu$ such that $\mu_{\alpha_{j}} \xrightarrow{*} \mu$.

The proof can be found in [Ha].
A set $E$ is called squarable with respect to measure $\mu$ ( $\mu$-squarable) if $\mu(\partial E)=0$.
Theorem 2.2.3.3. (Squarable Ring) The following holds:
sqr1) if $E_{1}, E_{2}$ are $\mu$-squarable, the sets $E_{1} \cap E_{2}, E_{1} \cup E_{2}, E_{1} \backslash E_{2}$ are $\mu$-squarable;
sqr2) for any couple: an open set $G$ and a compact set $K \subset G$ there exists a $\mu$-squarable set $E$ such that $K \subset E \subset G$.

Proof. The assertion sqr1) follows from

$$
\partial\left(E_{1} \cup E_{2}\right) \bigcup \partial\left(E_{1} \cap E_{2}\right) \bigcup \partial\left(E_{1} \backslash E_{2}\right) \subset \partial E_{1} \cup \partial E_{2}
$$

Let us prove sqr2). Let $K_{t}:=\{x: \exists y \in K:|x-y|<t\}$ be a $t$-neighborhood of the $K$. It is clear that for all the small $t$ we have $K \Subset K_{t} \Subset G$. The function $a(t):=\mu\left(K_{t}\right)$ is monotonic on $t$ and thus has no more than a countable set of jumps.

Let $t$ be a point of continuity of $a(t)$. Then

$$
\mu\left(\partial K_{t}\right) \leq \lim _{\epsilon \rightarrow 0}\left[\mu\left(K_{t+\epsilon}\right)-\mu\left(K_{t-\epsilon}\right)\right]=0
$$

Thus it is possible to set $E:=K_{t}$ for this $t$.
A family $\Phi$ of sets is called a dense ring if the following conditions hold:
dr1) $\forall F_{1}, F_{2} \in \Phi \Longrightarrow F_{1} \cup F_{2}, F_{1} \cap F_{2} \in \Phi$;
dr2) $\forall K, G: K \Subset G \exists F \in \Phi: K \subset F \subset G$.
The previous theorem shows that the class of $\mu$-squarable sets is a dense ring. The following theorem shows how one can extend a measure from a dense ring to the Borel algebra.

Let $\Phi$ be a dense ring and $\Delta(F), F \in \Phi$ a function of a set which satisfies the conditions:
$\Delta 1$ ) monotonicity on $\Phi: F_{1} \subset F_{2} \Longrightarrow \Delta\left(F_{1}\right) \leq \Delta\left(F_{2}\right)$;
$\Delta 2)$ additivity on $\Phi: \Delta\left(F_{1} \cup F_{2}\right) \leq \Delta\left(F_{1}\right)+\Delta\left(F_{2}\right)$ and $\Delta\left(F_{1} \cup F_{2}\right)=\Delta\left(F_{1}\right)+\Delta\left(F_{2}\right)$ if $F_{1} \cap F_{2}=\varnothing$
$\Delta 3$ ) continuity on $\Phi: \forall F \in \Phi$ and $\epsilon>0$ there exists a compact set $K$ and an open set $G \supset K$ such that $\forall F^{\prime} \in \Phi: K \subset F^{\prime} \subset G$ the inequality $\left|\Delta(F)-\Delta\left(F^{\prime}\right)\right|<\epsilon$ holds.

Theorem 2.2.3.4 (N. Bourbaki) There exists a measure $\mu$ such that

$$
\mu(F)=\Delta(F), \forall F \in \Phi
$$

iff the conditions $\Delta 1)-\Delta 3)$ hold.
Theorem 2.2.3.5 (Uniqueness of Measure) Under the conditions $\Delta 1)-\Delta 3$ ) the measure is defined uniquely by the formulae:

$$
\begin{align*}
& \mu(K)=\inf \{\Delta(F): F \in \Phi, F \supset K\}  \tag{2.2.3.1}\\
& \mu(G)=\sup \{\Delta(F): F \in \Phi, F \subset G\}  \tag{2.2.3.2}\\
& \mu(E)=\sup \{\mu(K): K \subset E\}=\inf \{\mu(G): G \supset E\} \tag{2.2.3.3}
\end{align*}
$$

and every $F \in \Phi$ is $\mu$-squarable.
For the proof see $[\mathrm{Bo}, \mathrm{Ch} .4, \operatorname{Sec} 3$, it. 10]. The squarability follows from (2.2.3.3).

The following theorem connects the convergence of measures on any dense ring and on the ring of sets squarable with respect to the limit measure.

Theorem 2.2.3.6 (Set-convergences) If $\mu_{n}(F) \rightarrow \mu(F)$ for all $F$ in a dense ring $\Phi$, then $\mu_{n}(E) \rightarrow \mu(E)$ for any $\mu$-squarable set $E$.

Proof. Suppose $\stackrel{\circ}{E} \neq \varnothing$.
Let $\epsilon>0$. By (2.2.3.3) one can find a compact set $K$ such that

$$
\begin{equation*}
\mu(K)+\epsilon \geq \mu\left(\circ_{E}^{E}\right)=\mu(E) . \tag{2.2.3.4}
\end{equation*}
$$

One can also find an open set $G$ such that

$$
\begin{equation*}
\mu(G)-\epsilon \leq \mu(\bar{E})=\mu(E) \tag{2.2.3.5}
\end{equation*}
$$

By property dr2) of a dense ring one can find $F, F^{\prime} \in \Phi$ such that

$$
K \subset F \subset \stackrel{\circ}{E} \subset E \subset \bar{E} \subset F^{\prime} \subset G
$$

Thus $\mu_{n}(F) \leq \mu_{n}(E) \leq \mu_{n}\left(F^{\prime}\right)$ and hence

$$
\begin{equation*}
\mu(F) \leq \varliminf_{n \rightarrow \infty} \mu_{n}(E) \leq \varlimsup_{n \rightarrow \infty} \mu_{n}(E) \leq \mu\left(F^{\prime}\right) . \tag{2.2.3.6}
\end{equation*}
$$

From (2.2.3.4) and (2.2.3.5) we obtain $0 \leq \mu\left(F^{\prime}\right)-\mu(F) \leq \mu(G)-\mu(K) \leq 2 \epsilon$ for arbitrarily small $\epsilon$. Thus from (2.2.3.6) we obtain

$$
\begin{equation*}
\underline{\lim }_{n \rightarrow \infty} \mu_{n}(E)=\varlimsup_{n \rightarrow \infty} \mu_{n}(E)=\mu(E) . \tag{2.2.3.7}
\end{equation*}
$$

That is to say that $\mu_{n}(E) \rightarrow \mu(E)$.
If $\stackrel{\circ}{E}=\varnothing$, then $\mu(\bar{E})=0$ by the definition of a squarable set. One can show in the same way that $\mu_{n}(E) \rightarrow 0$.

Now we connect the weak convergence to the convergence on squarable sets.
Theorem 2.2.3.7 (Set- and C*-convergences) The conditions

$$
\begin{equation*}
\mu_{n} \xrightarrow{*} \mu \tag{2.2.3.8}
\end{equation*}
$$

and $\mu_{n}(E) \rightarrow \mu(E)$ on $\mu$-squarable sets $E$ are equivalent.
Proof. Sufficiency of (2.2.3.8) follows from Theorem 2.2.3.1.
Exercise 2.2.3.1 Prove this.
Let us prove necessity.
For any compact set one can find a $\mu$-squarable $E$ such that $K \subset E$. Hence $\mu_{n}(K) \leq \mu(E)+1:=C(K)$ when $n$ is big enough.

By Helly's theorem (Theorem 2.2.3.2) there exists a measure $\mu^{\prime}$ and a subsequence $\mu_{n_{j}} \xrightarrow{*} \mu^{\prime}$. By the proved sufficiency, $\mu^{\prime}(E)=\mu(E)$ on a dense ring of the squarable sets. Thus $\mu^{\prime}=\mu$ by Uniqueness Theorem 2.2.3.5. And thus $\mu_{n} \xrightarrow{*} \mu$.

Denote by

$$
\mu_{E}(G):= \begin{cases}\mu(G \cap E) & \text { if } G \cap E \neq \varnothing \\ 0 & \text { if } G \cap E=\varnothing\end{cases}
$$

the restriction of $\mu$ on the set $E$.
Corollary 2.2.3.8 Let $\mu_{n} \xrightarrow{*} \mu$ and $E$ be a squarable set for $\mu$. Then $\left(\mu_{n}\right)_{E} \xrightarrow{*}(\mu)_{E}$.
Indeed, if $E$ is a squarable set for $\mu$ it is a squarable set for $\mu_{E}$. So Theorem 2.2.3.7 implies the corollary.
2.2.4 Let $\sigma\left(\mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}\right)$ be the $\sigma$-algebra of all the Borel sets, $\Phi_{i} \subset \sigma\left(\mathbb{R}^{m_{i}}\right)$, $i=$ 1,2 , be dense rings, $\Phi:=\Phi_{1} \otimes \Phi_{2} \subset \sigma\left(\mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}\right)$ be a ring generated by all the sets of form $F_{1} \times F_{2}, \quad F_{i} \in \Phi_{i}$.

Theorem 2.2.4.1 (Product of Rings) If $\Phi_{i}, i=1,2$ are dense rings, then $\Phi_{1} \otimes \Phi_{2}$ is a dense ring; if they consist of squarable sets, then $\Phi$ consists of squarable sets.

Proof. Let $K \subset G \subset \mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}$. For every point $x \in K$ one can (evidently) find $F_{1} \times F_{2}$ such that $x \subset F_{1} \times F_{2} \subset G$. One can find a finite covering and obtain a finite union $F$ of sets of such form. Thus $F \in \Phi_{1} \otimes \Phi_{2}$ and $F \subset G$.

The second assertion follows from the formula

$$
\partial\left(F_{1} \times F_{2}\right)=\left(\partial F_{1} \times F_{2}\right) \cup\left(F_{1} \times \partial F_{2}\right)
$$

Let $\mu_{i}$ be a measure on $\sigma\left(\mathbb{R}^{m_{i}}\right), i=1,2$, and $\mu:=\mu_{1} \otimes \mu_{2}$ the product of measures, i.e., a measure on $\sigma\left(\mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}\right)$ such that $\mu\left(E_{1} \times E_{2}\right)=\mu_{1}\left(E_{1}\right) \mu_{2}\left(E_{2}\right)$ for all $E_{i} \in \sigma\left(\mathbb{R}^{m_{i}}\right), i=1,2$.

Theorem 2.2.4.2 (Product of Measures) A measure $\mu_{1} \otimes \mu_{2}$ is uniquely defined by its values on $\Phi_{1} \otimes \Phi_{2}$.

The assertion follows from Theorem 2.2.4.1 and Uniqueness Theorem 2.2.3.5.
Theorem 2.2.4.3 (Fubini) Let $f\left(x_{1}, x_{2}\right)$ be a Borel function on $\mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}$. Then

$$
\begin{align*}
\int_{\mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}} f\left(x_{1}, x_{2}\right)\left(\mu_{1} \otimes \mu_{2}\right)\left(d x_{1} d x_{2}\right) & =\int_{\mathbb{R}^{m_{1}}} \mu_{1}\left(d x_{1}\right) \int_{\mathbb{R}^{m_{2}}} f\left(x_{1}, x_{2}\right) \mu_{2}\left(d x_{2}\right) \\
& =\int_{\mathbb{R}^{m_{2}}} \mu_{2}\left(d x_{2}\right) \int_{\mathbb{R}^{m_{1}}} f\left(x_{1}, x_{2}\right) \mu_{1}\left(d x_{1}\right), \tag{2.2.4.1}
\end{align*}
$$

if at least one of the parts of (2.2.4.1) is well defined.
The proof can be found in [Ha, Ch. VII, Sec. 36].

### 2.3 Distributions

2.3.1 Let us consider the set $\mathcal{D}(G)$ of all infinitely differentiable functions $\varphi(x), x \in$ $G \subset \mathbb{R}^{m}$.

It is a linear space because for any constants $c_{1}, c_{2}$,

$$
\begin{equation*}
\varphi_{1}, \varphi_{2} \in \mathcal{D}(G) \Longrightarrow c_{1} \varphi_{1}+c_{2} \varphi_{2} \in \mathcal{D}(G) \tag{D1}
\end{equation*}
$$

It is a topological space with convergence defined by

$$
\varphi_{n} \xrightarrow{\mathcal{D}} \varphi:\left\{\begin{array}{l}
\text { a) } \operatorname{supp} \varphi_{n} \subset K \Subset \mathbb{R}^{m}  \tag{D2}\\
\quad \text { for some compact } K \\
\text { and } \\
\text { b) } \varphi_{n} \rightarrow \varphi \text { uniformly on } K \\
\text { with all their derivatives. }
\end{array}\right.
$$

We consider some examples of functions $\varphi \in \mathcal{D}$. Set

$$
\alpha(t)= \begin{cases}C e^{-\frac{1}{1-t^{2}}}, & \text { for } t \in(-1 ; 1)  \tag{2.3.1.1}\\ 0, & \text { for } t \in(-1 ; 1)\end{cases}
$$

Evidently $\alpha(|x|) \in \mathcal{D}\left(\mathbb{R}^{m}\right)$ and $\operatorname{supp} \alpha \subset\{x:|x| \leq 1\}$.
Exercise 2.3.1.1 Check this.
Let us find $C$ such that

$$
\begin{equation*}
\int \alpha(|x|) d x=\sigma_{m} \int_{0}^{1} \alpha(t) t^{m-1} d t=1 \tag{2.3.1.2}
\end{equation*}
$$

where $\sigma_{m}$ is the area of the unit sphere $\{|x|=1\}$. Set

$$
\begin{equation*}
\alpha_{\varepsilon}(x):=\varepsilon^{-m} \alpha\left(\frac{|x|}{\varepsilon}\right) \tag{2.3.1.3}
\end{equation*}
$$

For any $\varepsilon$ we have $\alpha_{\varepsilon} \in \mathcal{D}$ and $\operatorname{supp} \alpha_{\varepsilon} \subset\{x:|x| \leq \varepsilon\}$.
Let $\psi(y), y \in K \subset G$ be a Lebesgue summable function. Consider the function

$$
\begin{equation*}
\psi_{\varepsilon}(x):=\int_{K} \psi(y) \alpha_{\varepsilon}(x-y) d y \tag{2.3.1.4}
\end{equation*}
$$

The function belongs to $\mathcal{D}(G)$ for $\varepsilon$ small enough and its support is contained in the $\varepsilon$-neighborhood of $K$.
2.3.2 Let $f(x), x \in G \subset \mathbb{R}^{m}$ be a locally summable function in $G$. The formula

$$
\begin{equation*}
\langle f, \varphi\rangle:=\int f(y) \varphi(y) d y, \varphi \in \mathcal{D}(G) \tag{2.3.2.1}
\end{equation*}
$$

defines a linear continuous functional on $\mathcal{D}$, i.e., one that satisfies the conditions

$$
\begin{gather*}
\left\langle f, c_{1} \varphi_{1}+c_{2} \varphi_{2}\right\rangle=c_{1}\left\langle f, \varphi_{1}\right\rangle+c_{2}\left\langle f, \varphi_{2}\right\rangle ; \\
\quad\left(\varphi_{n} \xrightarrow{\mathcal{D}} \varphi\right) \Longrightarrow\left\langle f, \varphi_{n}\right\rangle \rightarrow\langle f, \varphi\rangle .
\end{gather*}
$$

However, (2.3.2.1) does not exhaust all the linear continuous functionals as we will see further. An arbitrarily linear continuous functional on $\mathcal{D}(G)$ is called a Schwartz distribution and the linear topological space of the functionals is denoted as $\mathcal{D}^{\prime}(G)$.

Following are some examples of functionals that do not have the form of (2.3.2.1):

$$
\begin{equation*}
\left\langle\delta_{x}, \varphi\right\rangle:=\varphi(x) \tag{2.3.2.2}
\end{equation*}
$$

This distribution is called the Dirac delta-function. Further,

$$
\begin{equation*}
\left\langle\delta_{x}^{(n)}, \varphi\right\rangle:=(-1)^{n} \varphi^{(n)}(x) \tag{2.3.2.3}
\end{equation*}
$$

This distribution is called the nth derivative of the Dirac delta-function.
Exercise 2.3.2.1 Check that the functionals (2.3.2.2) and (2.3.2.3) are both distributions.

A distribution of the form (2.3.2.1) is called regular.
Theorem 2.3.2.1 (Du Bois Reymond) If two locally summable functions $f_{1}$ and $f_{2}$ define the same distribution, then they coincide almost everywhere.

For the proof see, e.g., [Hö, Thm. 2.1.6].
Note that the converse assertion is obvious.
A distribution $\mu$ is called positive if $\langle\mu, \varphi\rangle \geq 0$ for any $\varphi \in \mathcal{D}(G)$ such that $\varphi(x) \geq 0$ for all $x \in \mathbb{R}^{m}$. We shall write this as $\mu>0$ in $\mathcal{D}^{\prime}$.
Example 2.3.2.1 Let $\mu(E)$ be a measure in $G$. Then the distribution

$$
\begin{equation*}
\langle\mu, \varphi\rangle:=\int \varphi(x) \mu(d x) \tag{2.3.2.4}
\end{equation*}
$$

is positive.
This formula represents all the positive distributions as one can see from
Theorem 2.3.2.2 (Positive Distributions) Let $\mu>0$ in $\mathcal{D}(G)$. Then there exists a unique measure $\mu(E)$ such that the distribution $\mu$ is given by (2.3.2.4).

For the proof see, e.g., [Нö, Thm. 2.1.7].
2.3.3 Let us consider operations on distributions.

A product of a distribution $f$ by an infinitely differentiable function $\alpha(x)$ is defined by

$$
\begin{equation*}
\langle\alpha f, \varphi\rangle:=\langle f, \alpha \varphi\rangle . \tag{2.3.3.1}
\end{equation*}
$$

It is well defined because $\alpha \varphi \in \mathcal{D}$ too.

A sum of distributions $f_{1}$ and $f_{2}$ is defined by

$$
\begin{equation*}
\left\langle f_{1}+f_{2}, \varphi\right\rangle:=\left\langle f_{1}, \varphi\right\rangle+\left\langle f_{2}, \varphi\right\rangle, \tag{2.3.3.2}
\end{equation*}
$$

and the partial derivative $\frac{\partial}{\partial x_{k}}$ is defined by the equality

$$
\begin{equation*}
\left\langle\frac{\partial}{\partial x_{k}} f, \varphi\right\rangle:=\left\langle f,-\frac{\partial}{\partial x_{k}} \varphi\right\rangle . \tag{2.3.3.3}
\end{equation*}
$$

These definitions look reasonable because of the following
Theorem 2.3.3.1 (Operations on Distributions) The sum of regular distributions corresponds to the sum of the functions; the product of a regular distribution by an infinitely differentiable function corresponds to the product of the functions; the derivative of a regular distribution that is generated by a differentiable function corresponds to the derivative of that function.

Proof. We have, for example,

$$
\langle\alpha \cdot(f), \varphi\rangle:=\int f(x)[\alpha(x) \varphi(x)] d x=\int[\alpha(x) f(x)] \varphi(x) d x:=\langle(\alpha f), \varphi\rangle
$$

For the sum we have

$$
\begin{aligned}
\left\langle\left(f_{1}\right)+\left(f_{2}\right), \varphi\right\rangle:=\left\langle f_{1}, \varphi\right\rangle+\left\langle f_{2}, \varphi\right\rangle & =\int f_{1}(x) \varphi(x) d x+\int f_{2}(x) \varphi(x) d x \\
& =\int\left[f_{1}(x)+f_{2}(x)\right] \varphi(x) d x=\left\langle\left(f_{1}+f_{2}\right), \varphi\right\rangle
\end{aligned}
$$

Let $f(x)$ have the derivative $\frac{\partial}{\partial x_{k}} f$. Then

$$
\begin{aligned}
\left\langle\frac{\partial}{\partial x_{k}} f, \varphi\right\rangle:= & \left\langle f,-\frac{\partial}{\partial x_{k}} \varphi\right\rangle \\
= & \int f\left(x_{1}, x_{2}, \ldots, x_{m}\right)\left[-\frac{\partial}{\partial x_{k}} \varphi\left(x_{1}, x_{2}, \ldots, x_{m}\right)\right] d x_{1} d x_{2}, \ldots, d x_{m} \\
= & \int d x_{1}, \ldots, d x_{k-1} d x_{k+1}, \ldots, d x_{m} \\
& \times \int f\left(x_{1}, x_{2}, \ldots, x_{m}\right)\left[-\frac{\partial}{\partial x_{k}} \varphi\left(x_{1}, x_{2}, \ldots, x_{m}\right)\right] d x_{k} .
\end{aligned}
$$

Now we shall integrate by parts and all the substitution will vanish, because $\varphi$ is finite. So we obtain

$$
\left\langle\frac{\partial}{\partial x_{k}} f, \varphi\right\rangle=\int \frac{\partial}{\partial x_{k}} f(x) \varphi(x) d x .
$$

That is to say the derivative of the distribution corresponds to the function derivative.
2.3.4 We say that a sequence of distributions $f_{n}$ converges to a distribution $f$ if

$$
\begin{equation*}
\left\langle f_{n}, \varphi\right\rangle \rightarrow\langle f, \varphi\rangle \forall \varphi \in \mathcal{D}(G) \tag{2.3.4.1}
\end{equation*}
$$

Theorem 2.3.4.1 (Completeness of $\left.\mathcal{D}^{\prime}\right)$ If the sequence of numbers $\left\langle f_{n}, \varphi\right\rangle$ has a limit for every $\varphi \in \mathcal{D}(G)$, then this functional is a linear continuous functional on $\mathcal{D}(G)$, i.e., a distribution.

For the proof see, e.g., [Hö, Thm. 2.1.8].
Differentiation is continuous with respect to convergence of distributions.
Theorem 2.3.4.2 (Continuity of Differential Operators) If $f_{n} \rightarrow f$ in $\mathcal{D}(G)$, then $\frac{\partial}{\partial x_{k}} f_{n} \rightarrow \frac{\partial}{\partial x_{k}} f$.

Proof. Set in (2.3.4.1) $\varphi:=-\frac{\partial}{\partial x_{k}} \varphi$. Then

$$
\left\langle\frac{\partial}{\partial x_{k}} f_{n}, \varphi\right\rangle=\left\langle f_{n},-\frac{\partial}{\partial x_{k}} \varphi\right\rangle \rightarrow\left\langle f,-\frac{\partial}{\partial x_{k}} \varphi\right\rangle=\left\langle\frac{\partial}{\partial x_{k}} f, \varphi\right\rangle .
$$

The following theorem shows that the $\mathcal{D}^{\prime}$-convergence is the weakest of the convergences considered earlier.

Theorem 2.3.4.3 (Connection between Convergences) Let $f_{n}$ be a sequence of Lebesgue summable functions on domain $G$ such that at least one of the following conditions holds:
Cnvr1) $f_{n} \rightarrow f$ uniformly on any compact set $K \Subset G$ and $f$ is a locally summable function;
Cnvr2) $f_{n} \rightarrow f$ on any $K \Subset G$, satisfying the conditions of the Lebesgue theorem (Theorem 2.2.2.1);
Cnvr3) $f_{n} \downarrow f$ monotonically and $f$ is a locally summable function.
Then $f_{n} \rightarrow f$ in $\mathcal{D}^{\prime}(G)$.
Proof. All the assertions are corollaries of Section 2.2 .2 on passing to the limit under an integral.

Let us prove, for example, Cnvr3). Let $f_{n} \downarrow f$. Then

$$
\begin{equation*}
\left\langle f_{n}, \varphi\right\rangle=\int f_{n}(x) \varphi(x) d x=\int f_{n}(x) \varphi^{+}(x) d x-\int f_{n}(x) \varphi^{-}(x) d x \tag{2.3.4.2}
\end{equation*}
$$

where $\varphi^{+}$and $\varphi^{-}$are defined in (2.2.1.1).
Both last integrals in (2.3.4.2) have a limit by the B. Levy theorem (Theorem 2.2.2.2), and thus

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left\langle f_{n}, \varphi\right\rangle=\int f(x) \varphi^{+}(x) d x-\int f(x) \varphi^{-}(x) d x=\int f(x) \varphi(x) d x=\langle f, \varphi\rangle \tag{2.3.4.3}
\end{equation*}
$$

(2.3.4.3) means that $f_{n} \rightarrow f$ in $\mathcal{D}^{\prime}$.

Exercise 2.3.4.1 Prove Cnvr 1) and 2).
Theorem 2.3.4.4 ( $\mathcal{D}^{\prime}$ and $C^{*}$ convergences) Let $\mu_{n}, \mu$ be measures in $G$. The conditions $\mu_{n} \rightarrow \mu$ in $\mathcal{D}^{\prime}(G)$ and $\mu_{n} \xrightarrow{*} \mu$ are equivalent.

It is clear that the first condition is necessary for the second one. The sufficiency holds, because every continuous function can be approximated with functions that belong to $\mathcal{D}$. For more details see, e.g., [Hö, Thm. 2.1.9].

Let $\alpha_{\epsilon}(x)$ be defined as in (2.3.1.3). For any $f \in \mathcal{D}^{\prime}(D)$ we can consider the function $f_{\epsilon}(x):=\left\langle f, \alpha_{\epsilon}(x+\bullet)\right\rangle$. It is called a regularization of the distribution $f$.

Theorem 2.3.4.5.(Properties of Regularizations) The following holds:
reg1) $\quad f_{\epsilon}(x)$ is an infinitely differentiable function in any $K \Subset D$ for sufficiently small $\epsilon$;
reg2) $\quad f_{\epsilon}(x) \rightarrow f$ in $\mathcal{D}^{\prime}(D)$ as $\epsilon \downarrow 0$;
reg3) if $f_{n} \rightarrow f$ in $\mathcal{D}^{\prime}(D),\left(f_{n}\right)_{\epsilon} \rightarrow f_{\epsilon}$ uniformly with all its derivatives on any compact set in $D$.

The property reg1) follows from the formula

$$
\frac{\partial}{\partial x_{j}} f_{\epsilon}=\left\langle f, \frac{\partial}{\partial x_{j}} \alpha_{\epsilon}(x+\bullet)\right\rangle .
$$

The property reg2) follows from the assertion

$$
\phi_{\epsilon}(x):=\int \phi(y) \alpha_{\epsilon}(x+y) d y \rightarrow \phi(x) \text { in } \mathcal{D}(D)
$$

as $\epsilon \downarrow 0$.
For the proof of reg3) see [Hö, Theorems 2.1.8, 4.1.5].
Let us note the following assertion;
Theorem 2.3.4.6 (Continuity $\langle\bullet, \bullet\rangle$ ) The function

$$
\langle f, \phi\rangle: \mathcal{D}^{\prime}(G) \times \mathcal{D}(G) \mapsto \mathbb{R}
$$

is continuous in the appropriate topology.
I.e., $f_{n} \rightarrow f$ in $\mathcal{D}^{\prime}(G)$ and $\phi_{j} \rightarrow \phi$ in $\mathcal{D}(G)$ imply $\left\langle f_{n}, \phi_{j}\right\rangle \rightarrow\langle f, \phi\rangle$.

For the proof see [Hö, Theorem 2.1.8].
2.3.5 Let $G_{1} \subset G$. Then $\mathcal{D}^{\prime}(G) \subset \mathcal{D}^{\prime}\left(G_{1}\right)$, because every functional on $\mathcal{D}(G)$ can be considered as a functional on $\mathcal{D}\left(G_{1}\right)$.

A distribution $f \in \mathcal{D}^{\prime}(G)$ considered as a distribution in $\mathcal{D}^{\prime}\left(G_{1}\right)$ is called the restriction of $f$ to $G_{1}$ and is denoted $\left.f\right|_{G_{1}}$.

Theorem 2.3.5.1 (Sewing Theorem) Let $G_{\alpha} \subset \mathbb{R}^{m}$ be a family of domains and in every of them let there be a distribution $f_{\alpha} \in \mathcal{D}\left(G_{\alpha}\right)$, such that:

If $G_{\alpha_{1}} \cap G_{\alpha_{2}} \neq \varnothing$, the equality

$$
\begin{equation*}
\left.f_{\alpha_{1}}\right|_{G_{\alpha_{1}} \cap G_{\alpha_{2}}}=\left.f_{\alpha_{2}}\right|_{G_{\alpha_{1}} \cap G_{\alpha_{2}}} \tag{2.3.5.1}
\end{equation*}
$$

holds. Then there exists one and only one distribution $f \in \mathcal{D}^{\prime}(G)$ where $G=\bigcup_{\alpha} G_{\alpha}$ such that $\left.f\right|_{G_{\alpha}}=f_{\alpha}$.

In particular, it means that every distribution is defined uniquely by its restriction to a neighborhood of every point.

Let $\mathcal{D}\left(S_{R}\right)$ be a space of infinitely differentiable functions on the sphere $S_{R}:=$ $\{x:|x|=R\}$. The corresponding distribution space is denoted as $\mathcal{D}^{\prime}\left(S_{R}\right)$. The sewing theorem holds for this space in the following form:

Theorem 2.3.5.2 ( $\mathcal{D}^{\prime}$ on Sphere) Let a family of domains $\Omega_{\alpha}$ cover $S_{R}$ and in every of them let there be a distribution $f_{\alpha} \in \mathcal{D}\left(\Omega_{\alpha}\right)$, such that:

If $\Omega_{\alpha_{1}} \cap \Omega_{\alpha_{2}} \neq \varnothing$, the equality

$$
\begin{equation*}
f_{\alpha_{1}}\left|\Omega_{\alpha_{1} \cap \Omega_{\alpha_{2}}}=f_{\alpha_{2}}\right|_{\Omega_{\alpha_{1}} \cap \Omega_{\alpha_{2}}} \tag{2.3.5.2}
\end{equation*}
$$

holds. Then there exists one and only one distribution $f \in \mathcal{D}^{\prime}\left(S_{R}\right)$ such that $\left.f\right|_{\Omega_{\alpha}}=f_{\alpha}$.
2.3.6 Let

$$
\begin{equation*}
L:=\sum_{i, j} \frac{\partial}{\partial x_{i}} a_{i, j}(x) \frac{\partial}{\partial x_{j}}+q(x) \tag{2.3.6.1}
\end{equation*}
$$

be a differential operator of second order with infinitely differentiable coefficients $a_{i, j}, q$.

We will consider only three types of differential operators: a one-dimensional operator with constant coefficients,the Laplace operator and the so-called spherical operator (see Section 2.4).

For all these operators we have the following assertion which follows from the general theory (see, e.g., [Hö, Theorem 11.1.1]):

Theorem 2.3.6.1 (Regularity of Generalized Solution) If the equation $L u=0$ has a solution $u \in \mathcal{D}^{\prime}(G)$, then $u$ is a regular distribution and can be realized as an infinitely differentiable function.

A distribution that satisfies the equation

$$
\begin{equation*}
L u=\delta_{y} \quad \text { in } \mathcal{D}^{\prime}(G), \tag{2.3.6.2}
\end{equation*}
$$

where $\delta_{y}$ is a Dirac delta function (see (2.3.2.2)), is called a fundamental solution of $L$ at the point $y$.

Every differential operator that we are going to consider has a fundamental solution (see, e.g., [Hö, Theorem 10.2.1]).

A restriction of the equation (2.3.6.2) to the domain $G_{y}:=G \backslash y$ is a homogeneous equation $L u=0$ in $\mathcal{D}^{\prime}\left(G_{y}\right)$. Thus we have

Theorem 2.3.6.2 (Regularity of Fundamental Solution) The fundamental solution is an infinitely differentiable function outside the point $y$.
2.3.7 We will need further also the Fourier coefficients for the distribution on the circle.

Let $\mathcal{D}\left(S^{1}\right)$ be a set of all infinitely differentiable functions on the unit circle $S^{1}$. The set of all linear continuous functionals over $\mathcal{D}\left(S^{1}\right)$ with the corresponding topology (see 2.3.2) is the corresponding space of distributions $\mathcal{D}^{\prime}\left(S^{1}\right)$ for which all the previous properties of distributions holds.

The functions $\left\{e^{i k \phi}, k=0, \pm 1, \pm 2, \ldots\right\}$ belong to $\mathcal{D}\left(S^{1}\right)$. The Fourier coefficients of $\nu \in \mathcal{D}^{\prime}\left(S^{1}\right)$ are defined by

$$
\begin{equation*}
\hat{\nu}(k):=\left\langle\nu, e^{-i k \phi}\right\rangle \tag{2.3.7.1}
\end{equation*}
$$

The inverse operator is defined by

$$
\begin{equation*}
\langle\nu, g\rangle=\frac{1}{2 \pi} \sum_{k=-\infty}^{\infty} \hat{\nu}(k)\left\langle g, e^{i k \phi}\right\rangle \tag{2.3.7.2}
\end{equation*}
$$

and the series converges, in any case, for those $\nu$ that are finite derivatives of summable functions, because Fourier coefficients of $g$ decrease faster then every power of $x$.

The convolution of distribution $\nu \in \mathcal{D}^{\prime}\left(S^{1}\right)$ and $g \in \mathcal{D}\left(S^{1}\right)$ is defined by

$$
\begin{equation*}
\nu * g(\phi)=\langle\nu, g(\phi-\bullet)\rangle . \tag{2.3.7.3.}
\end{equation*}
$$

This is a function from $\mathcal{D}\left(S^{1}\right)$.
The convolution of distributions $\nu_{1}, \nu_{2} \in \mathcal{D}^{\prime}\left(S^{1}\right)$ is defined by

$$
\begin{equation*}
\left\langle\nu_{1} * \nu_{2}, g\right\rangle=\nu_{1} *\left(\nu_{2} * g\right) . \tag{2.3.7.4}
\end{equation*}
$$

In spite of the view it is commutative and

$$
\widehat{\nu_{1} * \nu_{2}}(k)=\hat{\nu}_{1}(k) \cdot \hat{\nu}_{2}(k) .
$$

Exercise 2.3.7.1 Count the Fourier coefficients of the functions

$$
\begin{equation*}
G\left(r e^{i \phi}\right)=\log \left|1-r e^{i \phi}\right| \tag{2.3.7.5}
\end{equation*}
$$

for $r>1, r=1, r<1$; the function defined by

$$
\begin{equation*}
\widetilde{\cos \rho}(\phi):=\cos \rho \phi,-\pi<\phi<\pi, \rho \in(0, \infty) \tag{2.3.7.6}
\end{equation*}
$$

and $2 \pi$-periodically extended; the function

$$
\begin{equation*}
\tilde{\phi} \sin p \phi, p \in \mathbb{N} \tag{2.3.7.7}
\end{equation*}
$$

where $\tilde{\phi}$ is the $2 \pi$-periodical extension of the function $f(\phi)=\phi, \phi \in[0,2 \pi)$.

## Exercise 2.3.7.2 Set

$$
\begin{equation*}
P_{p-1}\left(r e^{i \phi}\right):=\Re\left\{\sum_{k=1}^{p-1} \frac{r^{k} e^{i k \phi}}{k}\right\}, \quad p \in \mathbb{N} . \tag{2.3.7.8}
\end{equation*}
$$

Prove that for every distribution $\nu$ :

$$
\begin{equation*}
\left(P_{p-1}\left(r \widehat{e^{\bullet}}\right) * \nu\right)(p)=0 \tag{2.3.7.9}
\end{equation*}
$$

The same for the function

$$
G_{p}\left(r e^{i \phi}\right):=G\left(r e^{i \phi}\right)+P_{p}\left(r e^{i \phi}\right)
$$

for $r<1$.

### 2.4 Harmonic functions

2.4.1 We will denote as $\Delta$ the Laplace operator in $\mathbb{R}^{m}$ :

$$
\Delta:=\frac{\partial^{2}}{\partial x_{1}^{2}}+\cdots+\frac{\partial^{2}}{\partial x_{m}^{2}} .
$$

We introduce in $\mathbb{R}^{m}$ the spherical coordinate system by the formulae:

$$
\begin{aligned}
& x_{1}=r \sin \phi_{0} \sin \phi_{1} \ldots \sin \phi_{m-2} ; \\
& x_{2}=r \cos \phi_{0} \sin \phi_{1} \ldots \sin \phi_{m-2} ; \\
& x_{3}=r \cos \phi_{1} \sin \phi_{2} \ldots \sin \phi_{m-2} ; \\
& \ldots \ldots \ldots \ldots \\
& x_{k}=r \cos \phi_{k-2} \sin \phi_{k-1} \ldots \sin \phi_{m-2} ; \\
& \ldots \ldots \ldots \ldots \\
& x_{m}=r \cos \phi_{m-2},
\end{aligned}
$$

where

$$
0<\phi_{0} \leq 2 \pi ; \quad 0 \leq \phi_{j}<\pi, \quad j=\overline{1, m-2} ; \quad 0<r<\infty
$$

Passing to the coordinates $\left(r, \phi_{0}, \phi_{1}, \ldots, \phi_{m-2}\right)$ in the Laplace operator we obtain

$$
\Delta=\frac{1}{r^{m-1}} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r}+\frac{1}{r^{2}} \Delta_{\boldsymbol{x}^{0}}
$$

The operator $\Delta_{\boldsymbol{x}^{0}}$ is called spherical, and has the form

$$
\Delta_{\boldsymbol{x}^{0}}:=\sum_{i=0}^{m-2} \frac{1}{\Pi} \frac{\partial}{\partial \phi_{i}} \frac{\Pi}{\Pi_{i}} \frac{\partial}{\partial \phi_{i}},
$$

where

$$
\Pi:=\prod_{j=1}^{m-2} \sin ^{j} \phi_{j} ; \Pi_{i}:=\prod_{j=i+1}^{m-2} \sin ^{2} \phi_{j} ; \Pi_{m-2}:=1 .
$$

In particular, for $m=2$, i.e., for the polar coordinates,

$$
\Delta=\frac{1}{r} \frac{\partial}{\partial r} r \frac{\partial}{\partial r}+\frac{1}{r^{2}} \frac{\partial^{2}}{\partial \phi^{2}}
$$

A distribution $H \in \mathcal{D}^{\prime}(G)$ is called harmonic if it satisfies the equation $\Delta H=0$.
The next theorem follows from Theorem 2.3.6.1.
Theorem 2.4.1.1 (Smoothness of harmonic functions) Any harmonic distribution is equivalent to an infinitely differentiable function.

This function, of course, satisfies the same equation and is a harmonic function in the ordinary sense. A direct proof can be found, e.g., in [Ro, Ch. 1, § 2 (1.2.5), p. 60].

Let $f(z), z=x+\imath y$ be a holomorphic function in a domain $G \subset \mathbb{C}$. Then the functions $u(x, y):=\Re f(z)$ and $v(x, y):=\Im f(z)$ are harmonic in $G$. In particular, the functions $r^{n} \cos n \varphi$ and $r^{n} \sin n \varphi$ where $r=|z|, \varphi=\arg z$ are harmonic.

Set

$$
\mathcal{E}_{m}(x):= \begin{cases}-|x|^{2-m}, & \text { for } m \geq 3  \tag{2.4.1.1}\\ \log |z|, & \text { for } m=2\end{cases}
$$

(We will often denote points of the plane as $z$.)
It is easy to check that $\mathcal{E}_{m}(x)$ is a harmonic function for $|x| \neq 0$.
Set

$$
\theta_{m}:= \begin{cases}(m-2) \sigma_{m}, & \text { for } m \geq 3 \\ 2 \pi, & \text { for } m=2\end{cases}
$$

where $\sigma_{m}$ is the surface area of the unit sphere in $\mathbb{R}^{m}$.
Theorem 2.4.1.2 (Fundamental Solution) The function $\mathcal{E}_{m}(x-y)$ satisfies in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$ the equation ${ }^{1}$

$$
\begin{equation*}
\Delta_{x} \mathcal{E}_{m}(x-y)=\theta_{m} \delta(x-y) \tag{2.4.1.2}
\end{equation*}
$$

where $\delta(x)$ is the Dirac $\delta$-function (see 2.3.2).

[^0]Proof. Let us prove the equality (2.4.1.2) for $y=0$. Suppose $\phi \in \mathcal{D}\left(\mathbb{R}^{m}\right)$ and $\operatorname{supp} \phi \subset K \Subset \mathbb{R}^{m}$. We have

$$
\left\langle\Delta \mathcal{E}_{m}, \phi\right\rangle:=\int \mathcal{E}_{m}(x) \Delta \phi(x) d x=\lim _{\epsilon \rightarrow 0} \int_{|x| \geq \epsilon} \mathcal{E}_{m}(x) \Delta \phi(x) d x
$$

Transforming this integral by the Green formula and using the fact that $\phi$ is finite we obtain

$$
\int_{|x| \geq \epsilon} \mathcal{E}_{m}(x) \Delta \phi(x) d x=\int_{|x| \geq \epsilon} \Delta \mathcal{E}_{m}(x) \phi(x) d x+\int_{|x|=\epsilon} \mathcal{E}_{m} \frac{\partial \phi}{\partial n} d s-\int_{|x|=\epsilon} \phi \frac{\partial \mathcal{E}_{m}}{\partial n} d s
$$

where $d s$ is an element of surface area and $\frac{\partial}{\partial n}$ is the differentiation in the direction of the external normal.

Use the harmonicity of $\mathcal{E}_{m}$. Then the first integral is equal to zero. Further we have

$$
\int_{|x|=\epsilon} \mathcal{E}_{m} \frac{\partial \phi}{\partial n} d s=\left.\epsilon\left(\int_{\left|x^{0}\right|=1} \frac{\partial}{\partial r} \phi\left(r x^{0}\right) d s\right)\right|_{r=\epsilon}=O(\epsilon), \text { for } \epsilon \rightarrow 0
$$

For the third term we have

$$
\int_{|x|=\epsilon} \phi \frac{\partial \mathcal{E}_{m}}{\partial n} d s=\frac{m-2}{\epsilon^{m-1}} \epsilon^{m-1} \int_{\left|x^{0}\right|=1} \phi\left(r x^{0}\right) d s=[\phi(0)+o(1)](m-2) \sigma_{m} .
$$

Thus we obtain $\left\langle\Delta \mathcal{E}_{m}, \phi\right\rangle=\phi(0) \theta_{m}$, and this proves (2.4.1.2) for $y=0$.
It is clear that by changing $\phi(x)$ for $\phi(x+y)$ we obtain (2.4.1.2) in the general case.

We will consider now a domain $\Omega$ with a Lipschitz boundary (Lipschitz domain). It means that every part of $\partial \Omega$ can be represented in some local coordinates $\left(x, x^{\prime}\right), x \in \mathbb{R}, x^{\prime} \in \mathbb{R}^{m-1}$ in the form $x=f\left(x^{\prime}\right)$, where $f$ is a Lipschitz function, i.e.,

$$
\left|f\left(x_{1}^{\prime}\right)-f\left(x_{2}^{\prime}\right)\right| \leq M_{\partial \Omega}\left|x_{1}^{\prime}-x_{2}^{\prime}\right|
$$

where $M$ depends only on the whole $\partial \Omega$ and does not depend on this local part.
Let $G(x, y, \Omega)$ be the Green function of a Lipschitz domain $\Omega$.
It is known (see, e.g., [Vl, Ch. V, §28]) that the Green function has the following properties:

$$
\begin{gather*}
G(x, y, \Omega)<0, \text { for }(x, y) \in \Omega \times \Omega ; G(x, y, \Omega)=0 \text { for }(x, y) \in \Omega \times \partial \Omega ;  \tag{g1}\\
G(x, y, \bullet)=G(y, x, \bullet) ;  \tag{g2}\\
G(x, y, \bullet)-\mathcal{E}_{m}(x-y)=H(x, y) \tag{g3}
\end{gather*}
$$

where $H$ is harmonic on $x$ and on $y$ within $\Omega$;

$$
\begin{equation*}
-G\left(x, y, \Omega_{1}\right) \leq-G\left(x, y, \Omega_{2}\right) \text { for } \Omega_{1} \subset \Omega_{2} \tag{g4}
\end{equation*}
$$

From (g3) follows
Theorem 2.4.1.3 (Green Function) The equality

$$
\begin{equation*}
\Delta_{x} G(x, y, \Omega)=\theta_{m} \delta(x-y) \tag{2.4.1.3}
\end{equation*}
$$

holds in $\mathcal{D}^{\prime}(\Omega)$.
Let $f(x)$ be a continuous function on $\partial \Omega$. It is known (see, e.g., [Vl, Ch. V, $\S 29])$ that the function

$$
\begin{equation*}
H(x, f):=\int_{\partial \Omega} f(y) \frac{\partial}{\partial n} G(x, y, \Omega) d s_{y} \tag{2.4.1.4}
\end{equation*}
$$

is the only harmonic function that coincides with $f$ on $\partial \Omega$.
The unique solution of the Poisson equation

$$
\Delta u=p,\left.u\right|_{\partial \Omega}=f
$$

for a continuous function $p$ is given by the formula

$$
\begin{equation*}
u(x, f, p):=\int_{\partial \Omega} f(y) \frac{\partial}{\partial n_{y}} G(x, y, \Omega) d s_{y}+\theta_{m}^{-1} \int_{\Omega} G(x, y, \Omega) p(y) d y \tag{2.4.1.5}
\end{equation*}
$$

Let $D$ be an arbitrarily open domain. We can define a $G(x, y, D)$ in the following way. Consider a sequence $\Omega_{n}$ of a Lipschitz domain such that $\Omega_{n} \uparrow D$. The sequence of the corresponding Green functions $G\left(x, y, \Omega_{n}\right)$ monotonically decreases. If it is bounded from below in some point, it is bounded everywhere while $x \neq y$ (as it follows from Theorem 2.4.1.7). It can be shown that the limit exists for any domain, the boundary of which has positive capacity (see 2.5 and references there). We will mainly use the Green function for the Lipschitz domains.

Let $G\left(x, y, K_{a, R}\right)$ be the Green function of the ball $K_{a, R}:=\{|x-a|<R\}$.

## Theorem 2.4.1.4 (Green Function for a Ball)

$$
\left.G_{( } x, y, K_{a, R}\right)= \begin{cases}-|x-y|^{2-m}-\left(\frac{|y-a|\left|x-y_{a, R}^{*}\right|}{R}\right)^{2-m}, & \text { for } m \geq 3 \\ \log \frac{|\zeta-z| R}{|\zeta-a|\left|z-\zeta_{a . R}^{*}\right|} & \text { for } m=2\end{cases}
$$

where $y_{a . R}^{*}:=a+(y-a)\left(R^{2} /|y-a|^{2}\right)$ is the inversion of $y$ relative to the sphere $\{|x-a|=R\}$.

For the proof see, e.g., [Br, Ch. 6, § 3].
Theorem 2.4.1.5 (Poisson Integral) Let $H$ be a harmonic function in $K_{a, R}$ and continuous in its closure. Then

$$
\begin{equation*}
H(x)=\frac{1}{\sigma_{m} R} \int_{|x-a|=R} H(y) \frac{R^{2}-|x-a|^{2}}{|x-y|^{m}} d s_{y}, x \in K(a, R) . \tag{2.4.1.6}
\end{equation*}
$$

In particular, for $m=2$,

$$
H\left(a+r e^{\imath \phi}\right)=\frac{1}{2 \pi} \int_{0}^{2 \pi} H\left(a+R e^{\imath \psi}\right) \frac{R^{2}-r^{2}}{R^{2}-2 R r \cos (\phi-\psi)+r^{2}} d \psi
$$

This theorem follows from (2.4.1.4).
Theorem 2.4.1.6 (Mean Value) Let $H$ be harmonic in $G \subset \mathbb{R}^{m}$. Then

$$
\begin{equation*}
H(x)=\frac{1}{\sigma_{m} R^{m-1}} \int_{|x-a|=R} H(y) d s_{y} \tag{2.4.1.7}
\end{equation*}
$$

where $x \in G$ and $R$ is taken such that $K(x, R) \Subset G$.
We must only set $a:=x$ in (2.4.1.6). We can rewrite (2.4.1.7) in the form

$$
H(x)=\frac{1}{\sigma_{m}} \int_{|y|=1} H(x+R y) d s_{y}
$$

Theorem 2.4.1.7 (Harnack) Suppose the family $\left.\left\{H_{\alpha}\right), \alpha \in A\right\}$ of harmonic functions in $G$ satisfies the conditions

$$
\begin{align*}
H_{\alpha}(x) & \leq C(K), \text { for } x \in K  \tag{Har1}\\
H_{\alpha}\left(x_{0}\right) & \geq B>-\infty, \text { for } x_{0} \in K \tag{Har2}
\end{align*}
$$

for every compact $K \Subset G$ and $C(K), B$ are constants not depending on $\alpha$.
Then the family is precompact in the uniform topology, i.e., there exists such a sequence $H_{\alpha_{n}}$, and a function $H$ harmonic in the interior of $K$ and continuous in $K$ such that $H_{\alpha_{n}} \rightarrow H$ uniformly in every $K$.

One can prove by using (2.4.1.6) that $\left|\operatorname{grad} H_{\alpha}\right|$ are bounded on every compact set by a constant not depending on $\alpha$. Thus the family is uniformly continuous and thus it is precompact by the Ascoli theorem.

For details see, e.g., [Br, Supplement, § 7].
Theorem 2.4.1.8 (Uniform and $\mathcal{D}^{\prime}$-convergences) Suppose the sequence $H_{n}$ satisfies the conditions of the Harnack theorem and converges to a function $H$ in $\mathcal{D}^{\prime}(G)$. Then $H_{n}$ converges to $H$ uniformly on every compact set $K \Subset G$.

Of course, $H$ is harmonic in $G$.
Proof. By the Harnack theorem the family is precompact. Thus we must only prove the uniqueness of $H$. Suppose there exist two subsequences such that $H_{k}^{1} \rightarrow H^{1}$ and $H_{k}^{2} \rightarrow H^{2}$ uniformly on every compact $K \Subset G$.

By Connection between Convergences (Theorem 2.3.4.3) $H_{k}^{1} \rightarrow H^{1}$ and $H_{k}^{2} \rightarrow H^{2}$ in $\mathcal{D}^{\prime}$. Hence, $H^{1}=H^{2}$ in $\mathcal{D}^{\prime}(G)$. By the De Bois Raimond theorem (Theorem 2.3.2.1) $H^{1}=H^{2}$ almost everywhere and hence everywhere because these functions are continuous.

Let $D$ be a domain with a smooth boundary $\partial D$ and let $F \subset \partial D$. Set

$$
\omega(x, F, D):=\int_{F} \frac{\partial G}{\partial n_{y}}(x, y) d s_{y}
$$

It is called a harmonic measure of $F$ with respect to $D$. A harmonic measure can be defined for an arbitrary domain $D$ by a limit process similar to the one we had for the Green function. In this case the formula (2.4.1.4) has the form

$$
H(x, f):=\int_{\partial D} f(y) d \omega(x, y, D)
$$

However we can not assert that $H(x, f)$ coincides with $f$ in any point $x \in \partial D$. We can only consider it as an operator that maps a function defined on $\partial D$ to a harmonic function in $D$.

By (2.4.1.3) we obtain
Theorem 2.4.1.9 (Two Constants Theorem) Let $H$ be harmonic in $D$ and satisfy the conditions

$$
H(x) \leq A_{1} \text { for } x \in F ; H(x) \leq A_{2} \text { for } x \in \partial D \backslash F
$$

where $A_{1}$ and $A_{2}$ are constants. Then

$$
H(x) \leq A_{1} \omega(x, F, D)+A_{2} \omega(x, \partial D \backslash F, D) \text { for } x \in D
$$

Let $y_{a, R}^{*}$ be the inversion from Green Function for a Ball (Theorem 2.4.1.4). Set $y^{*}:=y_{0,1}^{*}$, i.e., the inversion relative to a unit sphere with the center in the origin. Let $G^{*}:=\left\{y^{*}: y \in G\right\}$ be the inversion of a domain $G$.

Theorem 2.4.1.10 (Kelvin's Transformation) If $H$ is harmonic in $G$, then

$$
\begin{equation*}
H^{*}(y):=|y|^{2-m} H\left(y^{*}\right) \tag{2.4.1.8}
\end{equation*}
$$

is harmonic in $G^{*}$.
For the proof you must honestly compute Laplacian of $H^{*}$. "The computation is straightforward but tedious" ([He, Thm. 2.24]). It is not so tedious if you use the spherical coordinate system.
Exercise 2.4.1.1 Do this.
2.4.2 Denote as $S_{1}:=\left\{x^{0}:\left|x^{0}\right|=1\right\}$ the unit sphere with center in the origin. A function $Y_{\rho}\left(x^{0}\right), x^{0} \in \Omega \subset S_{1}$ is called a spherical function of degree $\rho$ if it satisfies the equation

$$
\begin{equation*}
\Delta_{\boldsymbol{x}^{0}} Y+\rho(\rho+m-2) Y=0 . \tag{2.4.2.1}
\end{equation*}
$$

For $m=2$, (2.4.2.1) gets the form

$$
Y^{\prime \prime}(\theta)+\rho^{2} Y(\theta)=0, \quad \text { i.e., } \quad Y(\theta)=a \cos \rho \theta+b \sin \rho \theta
$$

Spherical functions are obtained if we solve the equation $\Delta H=0$ by the change $H(x)=|x|^{\rho} Y\left(x^{0}\right)$.

Theorem 2.4.2.1 (Sphericality and Harmonicity) The function $Y_{\rho}\left(x^{0}\right)$ is spherical in a domain $\Omega \subset S_{1}$ if and only if the functions $H(x)=|x|^{\rho} Y_{\rho}\left(x^{0}\right)$ and $H^{*}(x)=$ $|x|^{-\rho-m+2} Y_{\rho}\left(x^{0}\right)$ are harmonic in the cone

$$
\begin{equation*}
\operatorname{Con}(\Omega):=\left\{x=r x^{0}: x^{0} \in \Omega, 0<r<\infty\right\} . \tag{2.4.2.2}
\end{equation*}
$$

If $\rho=k, k \geq 0, k \in \mathbb{Z}$, and only in this case, $Y_{k}\left(x^{0}\right)$ is spherical on the whole $S_{1}$, $H(x)$ is a homogeneous harmonic polynomial of degree $k$ and $H^{*}$ is harmonic in $\mathbb{R}^{m} \backslash 0$.

For the proof see, e.g., [Ax, Ch. 5]
The spherical functions of an integer degree $k$ form a finite-dimension space of dimension

$$
\operatorname{dim}(m, k)=\frac{(2 k+m-2)(k+m-3)!}{(m-2)!k!}
$$

In particular, $d(2, k)=2$ for any $k$.
For different $k$ the spherical functions $Y_{k}\left(x^{0}\right)$ are orthogonal on $S_{1}$. In particular, for $m=2$, it means the orthogonality of the trigonometric functions system.

Theorem 2.4.2.2 (Expansion of a Harmonic Function) Let $H(x)$ be a harmonic function in the ball $K_{R}:=\{|x|<R\}$. There exists an orthonormal system of spherical functions $Y_{k}\left(x^{0}\right), k=\overline{0, \infty}$, depending on $H$ such that

$$
\begin{equation*}
H(x)=\sum_{k=0}^{\infty} c_{k} Y_{k}\left(x^{0}\right)|x|^{k}, \text { for }|x|<R \tag{2.4.2.3}
\end{equation*}
$$

For any such system we have

$$
\begin{equation*}
c_{k}=\frac{1}{R^{k}} \int_{S_{1}} H\left(R x^{0}\right) Y_{k}\left(x^{0}\right) d s_{x^{0}} \tag{2.4.2.4}
\end{equation*}
$$

For the proof see, e.g., [Ax, Ch. 10], [TT, Ch. 4, § 10].
Theorem 2.4.2.3 (Liouville) Let $H$ be harmonic in $\mathbb{R}^{m}$ and suppose

$$
\begin{equation*}
\liminf _{R \rightarrow \infty} R^{-\rho} \max _{|x|=R} H(x)<\infty \tag{2.4.2.5}
\end{equation*}
$$

holds. Then $H$ is a polynomial of a degree $q \leq \rho$.
Proof. We can suppose $H(0)=0$ because $H(x)-H(0)$ is harmonic and also satisfies (2.4.2.5). Let $R_{n} \rightarrow \infty$ be a sequence for which

$$
\begin{equation*}
R_{n}^{-\rho} \max _{|x|=R_{n}} H(x) \leq \text { const }<\infty . \tag{2.4.2.6}
\end{equation*}
$$

From (2.4.2.4) we obtain

$$
\begin{equation*}
\left|c_{k}\right| \leq A_{k} R^{-k} \int_{S_{1}}\left|H\left(R x^{0}\right)\right| d s_{x_{0}} \tag{2.4.2.7}
\end{equation*}
$$

where $A_{k}=\max _{S_{1}}\left|Y_{k}\left(x^{0}\right)\right|$.
From the mean value property (Theorem 2.4.1.6)

$$
\int_{S_{1}} H\left(R x^{0}\right) d s_{x_{0}}=H(0) \sigma_{m}=0
$$

Thus

$$
\begin{equation*}
\int_{S_{1}}\left|H\left(R x^{0}\right)\right| d s_{x_{0}}=2 \int_{S_{1}} H^{+}\left(R x^{0}\right) d s_{x_{0}} \leq 2 \sigma_{m} \max _{|x|=R} H(x) . \tag{2.4.2.8}
\end{equation*}
$$

From (2.4.2.8) and (2.4.2.7) we have

$$
\begin{equation*}
\left|c_{k}\right| \leq 2 A_{k} R^{-k} \sigma_{m} \max _{|x|=R} H(x) . \tag{2.4.2.9}
\end{equation*}
$$

Set $R:=R_{n}$ and $k>\rho$. Passing to the limit when $n \rightarrow \infty$, we obtain $c_{k}=0$ for $k>\rho$. Then (2.4.2.3.) implies that $H$ is a harmonic polynomial of degree $q \leq \rho$.

### 2.5 Potentials and capacities

2.5.1 Let $G(x, y . D)$ be the Green function of a Lipschitz domain $D$. We will suppose it is extended as zero outside of $D$.

$$
\Pi(x, \mu, D):=-\int G(x, y, D) \mu(d y)
$$

is called the Green potential of $\mu$ relative to $D$. The domain of integration will always be $\mathbb{R}^{m}$.

Theorem 2.5.1.1 (Green Potential Properties) The following holds:
GPo1) $\Pi(x, \mu, D)$ is lower semicontinuous;
$\mathrm{GPo} 2)$ it is summable over any $(m-1)$-dimensional hyperplane or smooth hypersurface;
GPo3) $\Delta \Pi(\bullet, \mu, D)=-\theta_{m} \mu$ in $\mathcal{D}^{\prime}(D)$;
GPo4) the reciprocity law holds:

$$
\int \Pi\left(x, \mu_{1}, D\right) \mu_{2}(d x)=\int \Pi\left(x, \mu_{2}, D\right) \mu_{1}(d x)
$$

GPo5) semicontinuity in $\mu$ : if $\mu_{n} \rightarrow \mu$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$, then

$$
\liminf _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \geq \Pi(x, \mu, D)
$$

GPo6) continuity in $\mu$ in $\mathcal{D}^{\prime}:$ if $\mu_{n} \rightarrow \mu$, then $\Pi\left(\bullet, \mu_{n}, D\right) \rightarrow \Pi(\bullet, \mu, D)$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$ and in $\mathcal{D}^{\prime}\left(S_{R}\right)$, where $S_{R}$ is the sphere $\{|x|=R\}$.

Proof. Let us prove GPo1). Let $N>0$. Set $G_{N}(x, y):=\max (G(x, y),-N)$, a truncation of the function $G(x, y)$.

The functions $G_{N}$ are continuous in $\mathbb{R}^{m} \times \mathbb{R}^{m}$ and $G_{N}(x, y) \downarrow G(x, y)$ for every $(x, y)$ when $N \rightarrow \infty$. Set

$$
\Pi_{N}(x, \mu, D):=-\int G_{N}(x, y, D) \mu(d y)
$$

The functions $\Pi_{N}$ are continuous and $\Pi_{N}(x, \bullet) \uparrow \Pi(x, \bullet)$ by the B . Levy theorem (Theorem 2.2.2.2). Then $\Pi_{N}(x, \bullet)$ is lower semicontinuous by the Second Criterion of semicontinuity (Theorem 2.1.2.9).

Let us prove GPo5). From Theorem 2.3.4.4 ( $\mathcal{D}^{\prime}$ and $C^{*}$ convergences)

$$
\lim _{n \rightarrow \infty} \Pi_{N}\left(x, \mu_{n}, D\right)=\Pi_{N}(x, \mu, D)
$$

Further $\Pi\left(x, \mu_{n}, D\right) \geq \Pi_{N}\left(x, \mu_{n}, D\right)$, hence

$$
\liminf _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \geq \Pi_{N}(x, \mu, D)
$$

Passing to the limit while $N \rightarrow \infty$, we obtain GPo5).
The assertion GPo2) follows from the local summability of the function $|x|^{2-m}$ that can be checked directly.

Let us prove GPo3). For $\phi \in \mathcal{D}(D)$ we have

$$
\begin{aligned}
\langle\Delta \Pi, \phi\rangle:=\langle\Pi, \Delta \phi\rangle & =-\int \mu(d y) \int G(x, y, D) \Delta \phi(x) d x \\
& =-\int\left\langle\Delta_{x} G(\bullet, y, D), \phi\right\rangle \mu(d y)=-\theta_{m} \int \phi(y) \mu(d y) \\
& =-\theta_{m}\langle\mu, \phi\rangle
\end{aligned}
$$

since

$$
\left\langle\Delta_{x} G(\bullet, y, D), \phi\right\rangle=\theta_{m} \phi(y)
$$

by Theorem 2.4.1.3. The property GPo4) follows from the symmetry of $G(x, y, \bullet)$ (property (g2)).

Let us prove GPo6). Note that integral $\int|x|^{m-1} d x$ converges locally in $\mathbb{R}^{m}$ and in $\mathbb{R}^{m-1}$. From this one can obtain by some simple estimates that functions
$\Psi(y):=\int G(x, y, D) \psi(x) d x$ while $\psi \in \mathcal{D}\left(\mathbb{R}^{m}\right)$ and $\Theta(y):=\int_{S_{R}} G(x, y, D) \theta(x) d s_{x}$ while $\theta \in \mathcal{D}\left(S_{R}\right)$ are continuous on $y \in \mathbb{R}^{m}$.

Now we have

$$
\left\langle\Pi\left(\bullet, \mu_{n}, D\right), \psi\right\rangle=\int \Psi(y) \mu_{n}(d y) \rightarrow \int \Psi(y) \mu(d y)=\langle\Pi(\bullet, \mu, D), \psi\rangle
$$

Thus the first assertion in GPo6) is proved. The second one can be proved in the same way.

Set $\nu:=\mu_{1}-\mu_{2}$, and let $\Pi(x, \nu, D):=\Pi\left(x, \mu_{1}, D\right)-\Pi\left(x, \mu_{2}, D\right)$ be a potential of this charge. Consider the boundary problem of the form

$$
\begin{equation*}
\Delta u=\mu_{1}-\mu_{2}, \text { in } \mathcal{D}^{\prime}(D),\left.\quad u\right|_{\partial D}=f \tag{2.5.1.2}
\end{equation*}
$$

where $f$ is a continuous function.
Theorem 2.5.1.2 (Solution of Poisson Equation) The solution of the boundary problem (2.5.1.2) is given by the formula

$$
u(x)=H(x, f)-\theta_{m}^{-1} \Pi(x, \nu, D)
$$

where $H(x, f)$ is the harmonic function from (2.4.1.4).
Proof. Since $\left.\Pi(x, \nu, D)\right|_{\partial D}=0$, the function $u(x)$ satisfies the boundary condition. Using GPo3) we obtain

$$
\Delta u=\Delta H-\left[\theta_{m}\right]^{-1} \Delta \Pi=\mu_{1}-\mu_{2} .
$$

A potential of the form

$$
\Pi(x, \mu):=\int \frac{\mu(d y)}{|x-y|^{m-2}}
$$

is called a Newton potential. It is the Green potential for $D=\mathbb{R}^{m}$. The potential

$$
\Pi(z, \mu)=-\int \log |z-\zeta| \mu(d \zeta)
$$

is called logarithmic.
2.5.2 Let $K \Subset D$. The quantity

$$
\begin{equation*}
\operatorname{cap}_{G}(K, D):=\sup \mu(K) \tag{2.5.2.1}
\end{equation*}
$$

where the supremum is taken over all mass distributions $\mu$ for which the following conditions are satisfied:

$$
\begin{align*}
\Pi(x, \mu, D) & \leq 1  \tag{2.5.2.2}\\
\operatorname{supp} \mu & \subset K \tag{2.5.2.3}
\end{align*}
$$

is called the Green capacity of the compact set $K$ relative to the domain $D$.

Theorem 2.5.2.1 (Properties of $\operatorname{cap}_{G}$ ) For $\mathbf{c a p}_{G}$ the following properties hold:
capG1) monotonicity with respect to $K: K_{1} \subset K_{2}$ implies $\operatorname{cap}_{G}\left(K_{1}, D\right) \leq$ $\operatorname{cap}_{G}\left(K_{2}, D\right)$.
capG2) monotonicity with respect to $D: K \Subset D_{1} \subset D_{2}$ implies $\operatorname{cap}_{G}\left(K, D_{1}\right) \geq$ $\operatorname{cap}_{G}\left(K, D_{2}\right)$
capG3) subadditivity with respect to $K$ :

$$
\operatorname{cap}_{G}\left(K_{1} \cup K_{2}, D\right) \leq \operatorname{cap}_{G}\left(K_{1}, D\right)+\operatorname{cap}_{G}\left(K_{2}, D\right)
$$

Proof. The set of all mass distributions that satisfy (2.5.2.2) for $K=K_{1}$ is not less than the analogous set for $K=K_{2}$. Thus capG1) holds.

By the Green function property (g3) (see § 2.4.1) $-G\left(x, y, D_{1}\right) \leq-G\left(x, y, D_{2}\right)$. Thus the set of all $\mu$ that satisfy (2.5.2.2) for $D=D_{1}$ is wider than for $D=D_{2}$. Hence capG2) holds.

Let $\operatorname{supp} \mu \subset K_{1} \cup K_{2}$ and let $\mu_{1}, \mu_{2}$ be the restrictions of $\mu$ to $K_{1}, K_{2}$ respectively.

If $\mu$ satisfies (2.5.2.2) for $K:=K_{1} \cup K_{2}$ then $\mu_{1}, \mu_{2}$ satisfy (2.5.2.2) for $K:=K_{1}, K_{2}$ respectively.

From the inequality

$$
\mu\left(K_{1} \cup K_{2}\right) \leq \mu\left(K_{1}\right)+\mu\left(K_{2}\right)
$$

we obtain that

$$
\mu\left(K_{1} \cup K_{2}\right) \leq \operatorname{cap}_{G}\left(K_{1}, D\right)+\operatorname{cap}_{G}\left(K_{2}, D\right)
$$

for any $\mu$ with supp $\mu \subset K_{1} \cup K_{2}$. Thus capG3) holds.
The equivalent definition of the Green capacity is given by
Theorem 2.5.2.2 (Dual Property) The following holds:

$$
\begin{equation*}
\operatorname{cap}_{G}(K, D)=\left[\inf _{\mu} \sup _{x \in D} \Pi(x, \mu, D)\right]^{-1} \tag{2.5.2.4}
\end{equation*}
$$

where the infimum is taken over all mass distributions $\mu$ such that $\mu(K)=1$.
For the proof see, e.g., [La, Ch. $2, \S 4$ it. 18]. For $D=\mathbb{R}^{m}, m \geq 3$, the Green capacity is called Wiener capacity $\left(\boldsymbol{c a p}_{m}(K)\right)$. It has the following properties in addition to those of the Green capacity:
capW1) invariance with respect to translations and rotations, i.e.,

$$
\operatorname{cap}_{m}\left(V\left(K+x_{0}\right)\right)=\operatorname{cap}_{m}(K),
$$

where $V K$ and $K+x_{0}$ are the rotation and the translation of $K$ respectively.
The presence of these properties brings the notion of capacity closer to the notion of measure. Thus it is natural to extend the capacity to the Borel algebra of sets.

The Wiener capacity of an open set is defined as

$$
\operatorname{cap}_{m}(D):=\sup _{K} \operatorname{cap}_{m}(K)
$$

where the supremum is taken over all compact $K \Subset D$.
The outer and inner capacity of any set $E$ can be defined by the equalities

$$
\overline{\operatorname{cap}}_{m}(E):=\inf _{D \supset E} \operatorname{cap}_{m}(D) ; \underline{\operatorname{cap}}_{m}(E):=\sup _{K \subset E} \operatorname{cap}_{m}(K) .
$$

A set $E$ is called capacible if $\overline{\mathbf{c a p}}_{m}(E)=\underline{\mathbf{c a p}}_{m}(E)$.
Theorem 2.5.2.3 (Choquet) Every set E belonging to the Borel ring is capacible.
For the proof see, e.g., [La, Ch2, Thm. 2.8].
Sets which have "small size" are sets of zero capacity. We emphasize the following properties of these sets:
$\operatorname{capZ1)}$ If $\operatorname{cap}_{m}\left(E^{j}\right)=0, j=1,2, \ldots$ then $\operatorname{cap}_{m}\left(\cup_{1}^{\infty} E^{j}\right)=0$;
capZ2) Having the property of zero capacity does not depend on the type of capacity: Green, Wiener or logarithmic capacity that we define below.
Example 2.5.2.1 Using Theorem 2.5.2.2 we obtain that any point has zero capacity, because for every mass distribution concentrated in the point the potential is equal to infinity. The same holds for any set of zero $m-2$ Hausdorff measure (see 2.5.4).

Example 2.5.2.2 Any ( $m-1$ )-hyperplane or smooth hypersurface has positive capacity, because the potential with masses uniformly distributed over the surface is bounded.

The Wiener 2-capacity can be defined naturally only for sets with diameter less then 1 , because the logarithmic potential is positive only when this condition holds.

Instead, one can use the logarithmic capacity which is defined by the formulae

$$
\begin{equation*}
\operatorname{cap}_{l}(K):=\exp \left[-\mathbf{c a p}_{2}(K)\right] \tag{2.5.2.5}
\end{equation*}
$$

for $K \subset\{|z|<1\}$ and

$$
\operatorname{cap}_{l}(K):=t^{-1} \mathbf{c a p}_{l}(t K)
$$

for any other bounded $K$, where $t$ is chosen in such a way that $t K \subset\{|z|<1\}$.
One can check that this definition is correct, i.e., it does not depend on $t$.

### 2.5.3

Theorem 2.5.3.1 (Balayage; sweeping) Let $D$ be a domain such that $\partial D \Subset \mathbb{R}^{m}$, and supp $\mu \Subset D$. Then there exists a mass distribution $\mu_{b}$ such that for $m \geq 3$, or for $m=2$ and for $D$ which is a bounded domain, the following holds:
bal1) $\Pi\left(x, \mu_{b}\right)<\Pi(x, \mu)$ for $x \in D$;
bal2) $\Pi\left(x, \mu_{b}\right)=\Pi(x, \mu)$ for $x \notin \bar{D}$;
bal3) $\operatorname{supp} \mu_{b} \subset \partial D$;
bal4) $\quad \mu_{b}(\partial D)=\mu(D)$.
If $m=2$ and the domain is unbounded, a potential of the form

$$
\hat{\Pi}(z, \mu):=-\int \log |1-z / \zeta| \mu(d \zeta)
$$

satisfies all the properties.
Proof. We will prove this theorem when $\partial D$ is smooth enough. For $y \in D, x \in$ $\mathbb{R}^{m} \backslash \bar{D}$ the function $|x-y|^{2-m}$ is a harmonic function of $y$ on $D$.

Since $|x-y|^{2-m} \rightarrow 0$ as $y \rightarrow \infty$ we can apply the Poisson formula (2.4.1.4) even if $D$ is unbounded. Thus

$$
\begin{equation*}
|x-y|^{2-m}=\int_{\partial D}\left|x-y^{\prime}\right|^{2-m} \frac{\partial G}{\partial n_{y^{\prime}}}\left(y, y^{\prime}\right) d s_{y^{\prime}} \tag{2.5.3.1}
\end{equation*}
$$

where $G$ is the Green function of $D$. From this we have

$$
\int_{D}|x-y|^{2-m} \mu(d y)=\int_{\partial D}\left|x-y^{\prime}\right|^{2-m} d s_{y^{\prime}}\left(\int_{D} \frac{\partial G}{\partial n_{y^{\prime}}}\left(y, y^{\prime}\right) \mu(d y)\right) .
$$

The inner integral is nonnegative, because $\frac{\partial G}{\partial n}>0$ for $y^{\prime} \in \partial D$. Let us denote

$$
\mu_{b}\left(d y^{\prime}\right):=\left(\int_{D} \frac{\partial G}{\partial n}\left(y, y^{\prime}\right) \mu(d y)\right) d s_{y^{\prime}}
$$

Then we obtain the properties bal2) and bal3).
The potential $\Pi\left(x, \mu_{b}\right)$ is harmonic in $D$. Thus the function

$$
u(x):=\Pi\left(x, \mu_{b}\right)-\Pi(x, \mu)
$$

is a subharmonic function (see Theorem 2.6.4.1). Every subharmonic function satisfies the maximum principle (see Theorem 2.6.1.2), i.e.,

$$
u(x)<\sup _{y \in \partial D} u(y)=0
$$

Thus the property bal1) is fulfilled. To prove bal4) we can write the identity

$$
\int_{\partial G} \mu_{b}\left(d y^{\prime}\right)=\int_{G} \mu(d y) \int_{\partial G} \frac{\partial G}{\partial n_{y^{\prime}}}\left(y, y^{\prime}\right) d y^{\prime}
$$

The inner integral is equal to 1 identically, because the function $\equiv 1, y \in G$ is harmonic. Thus bal4) is true.

Consider now the special case when $m=2$, and $D$ is an unbounded domain. Since $\log |1-z / \zeta| \rightarrow 0$ when $\zeta \rightarrow \infty$, we obtain an equality like (2.5.3.1). Repeating the previous reasoning we obtain the last assertion for $D$ with a smooth boundary.

Exercise 2.5.3.1 Check this in detail.
For the general case see [La, Ch. 4, §1]; [Ca, Ch. 3, Thm. 4].

Pay attention that the swept potential $\Pi\left(x, \mu_{b}\right)$ is also a solution of the Dirichlet problem in the domain $D$ and the boundary function $f(x)=\Pi(x, \mu)$ in the following sense:

Theorem 2.5.3.2 (Wiener) The equality bal2) holds in the points $x \in \partial D$ which can be reached by the top of a cone placed outside D. For $m=2$ it can fail only for isolated points.

For the proof see $[\mathrm{He}],[\mathrm{La}, \mathrm{Ch} .4, \S 1, \mathrm{Thm} .4 .3$.$] .$
The points of $\partial D$ where the equality bal2) does not hold are called irregular.
Theorem 2.5.3.3 (Kellogg's Lemma) The set of all the irregular points of $\partial D$ has zero capacity.

For the proof see, e.g., [He], [La, Ch. 4, § 2, it. 10].
One can often compute the capacity using the following
Theorem 2.5.3.4 (Equilibrium distribution) For any compact $K$ with $\mathbf{c a p}_{m}(K)>0$ there exists a mass distribution $\lambda_{K}$ such that the following holds:
eq1) $\Pi(x, \lambda)=1, x \in \bar{D} \backslash E, \operatorname{cap}_{m}(E)=0$;
eq2) $\operatorname{supp} \lambda_{K} \subset \partial K$;
eq3) $\lambda_{K}(\partial K)=\boldsymbol{c a p}_{m}(K)$.
For the proof see [He], [La, Ch. 2, §1, it. 3, Thm. 2.3].
Let us note that the set $E$ in the previous theorem is a set of irregular points.
The mass distribution $\lambda_{K}$ is called equilibrium distribution, and the corresponding potential is called equilibrium potential.
2.5.4 Let $h(x), x \geq 0$ be a positive continuous, monotonically increasing function which satisfies the condition $h(0)=0$. Let $\left\{K_{j}^{\epsilon}\right\}$ be a family of balls such that their diameters $d_{j}:=d\left(K_{j}^{\epsilon}\right)$ are no bigger then $\epsilon$. Let us denote

$$
m_{h}(E, \epsilon):=\inf \sum h\left(\frac{1}{2} d\left(K_{j}^{\epsilon}\right)\right)
$$

where the infimum is taken over all coverings of the set $E$ by the families $\left\{K_{j}^{\epsilon}\right\}$.
The quantity

$$
m_{h}(E):=\lim _{\epsilon \rightarrow 0} m_{h}(E, \epsilon)
$$

is called $h$-Hausdorff measure [Ca, Ch. II].
Theorem 2.5.4.1 (Properties of $m_{h}$ ) The following properties hold:
h1) monotonicity:

$$
E_{1} \subset E_{2} \Longrightarrow m_{h}\left(E_{1}\right) \leq m_{h}\left(E_{2}\right)
$$

h2) countable additivity:

$$
m_{h}\left(\cup E_{j}\right)=\sum m_{h}\left(E_{j}\right) ; E_{j} \cap E_{i}=\varnothing, \text { for } i \neq j ; \quad E_{j} \in \sigma\left(\mathbb{R}^{m}\right)
$$

We will quote two conditions (necessary and sufficient) that connect the $h$ measure to the capacity (see, [La, Ch. 3, §4, it. 9, 10].

Theorem 2.5.4.2 Let $\boldsymbol{\operatorname { c a p }} E=0$. Then $m_{h}(E)=0$ for all $h$ such that

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

Theorem 2.5.4.3 Let $h(r)=r^{m-2}$ for $m \geq 3$ and $h(r)=(\log 1 / r)^{-1}$ for $m=2$. If the $h$-measure of a set $E$ is finite, then $\mathbf{c a p}_{m}(E)=0$.

Side by side with the Hausdorff measure the Carleson measure (see, $[\mathrm{Ca}$, Ch. II], is often considered. It is defined by

$$
m_{h}^{C}(E):=\inf \sum h\left(0.5 d_{j}\right)
$$

where the infimum is taken over all coverings of the set $E$ with balls of radii $0.5 d_{j}$. The inequality $m_{h}^{C}(E) \leq m_{h}(E)$ obviously holds. Let $\beta-\operatorname{mes}_{C} E$ be the Carleson measure for $h=r^{\beta}$. The following assertion connects the $\beta-\operatorname{mes}_{C}$ to capacity.

Theorem 2.5.4.4 The following inequalities hold:

$$
\begin{aligned}
\beta-\operatorname{mes}_{C} E & \leq N(m)\left(\operatorname{cap}_{m}(E)\right)^{\beta / m-2}, \text { for } m \geq 3, \beta>m-2 \\
\beta-\operatorname{mes}_{C}(E) & \leq 18 \operatorname{cap}_{l}(E), \text { for } m=2, \beta>0
\end{aligned}
$$

where $N$ depends only on the dimension of the space.
For the proof see [La, Ch. III, $\S 4$, it. 10, Cor. 2].
2.5.5 Now we will formulate an analog of the Luzin theorem for potentials.

Theorem 2.5.5.1 Let $\operatorname{supp} \mu=K$ and let the potential $\Pi(x, \mu)$ be bounded on $K$. Then for any $\delta>0$ there exists a compact set $K^{\prime} \subset K$ such that $\mu\left(K \backslash K^{\prime}\right)<\delta$ and the potential $\Pi\left(x, \mu^{\prime}\right)$ of the measure $\mu^{\prime}:=\left.\mu\right|_{K}$ (the restriction of $\mu$ to $K$ ) is continuous.

For the proof see, e.g., [La, Ch. 3, § 2, it. 3, Thm. 3.6].
Let us prove the following assertion:
Theorem 2.5.5.2 Let $\mathbf{c a p} K>0$. Then for arbitrarily small $\epsilon>0$ there exists a measure $\mu$ such that $\operatorname{supp} \mu \subset K$, the potential $\Pi(x, \mu)$ is continuous and $\mu(K)>$ $\boldsymbol{\operatorname { c a p }}(K)-\epsilon$.

Proof. Consider the equilibrium distribution $\lambda_{K}$ on $K$. Its potential is bounded by Theorem 2.5.3.4. By Theorem 2.5.5.1 we can find a mass distribution $\mu$ such that $\Pi(x, \mu)$ is continuous, supp $\mu \subset K$ and $\mu(K)>\lambda_{K}(K)-\epsilon=\boldsymbol{\operatorname { c a p }}(K)-\epsilon$.

### 2.6 Subharmonic functions

2.6.1 Let $u(x), x \in D \subset \mathbb{R}^{m}$ be a measurable function bounded from above which can be $-\infty$ on a set of no more than zero measure.

Let us denote as

$$
\begin{equation*}
\mathcal{M}(x, r, u):=\frac{1}{\sigma_{m} r^{m-1}} \int_{S_{x, r}} u(y) d s_{y} \tag{2.6.1.1}
\end{equation*}
$$

the mean value of $u(x)$ on the sphere $S_{x, r}:=\{y:|y-x|=r\}$.
The function $\mathcal{M}(x, r, u)$ is defined if $S_{x, r} \subset D$, but it can be $-\infty$ a priori.
A function $u(x)$ is called subharmonic if it is upper semicontinuous, $\not \equiv-\infty$, and for any $x \in D$ there exists $\epsilon=\epsilon(x)$ such that the inequality

$$
\begin{equation*}
u(x) \leq \mathcal{M}(x, r, u) \tag{2.6.1.2}
\end{equation*}
$$

holds for all $r<\epsilon$.
The class of functions subharmonic in $D$ will be denoted as $S H(D)$.
Example 2.6.1.1 The function

$$
u(x):=-|x|^{2-m}, x \in \mathbb{R}^{m}
$$

belongs to $S H\left(\mathbb{R}^{m}\right)$ for $m \geq 3$, and the function

$$
u(z):=\log |z|, z \in \mathbb{R}^{2}
$$

is subharmonic in $\mathbb{R}^{2}$.
Example 2.6.1.2 Let $f(z)$ be a holomorphic function in a plane domain $D$. Then $\log |f(z)| \in S H(D)$.

Example 2.6.1.3 Let $f=f\left(z_{1}, z_{2}, \ldots, z_{n}\right)$ be a holomorphic function of $z=$ $\left(z_{1}, \ldots, z_{n}\right)$. Then $u\left(x_{1}, y_{1}, \ldots, x_{n}, y_{n}\right):=\log \left|f\left(x_{1}+i y_{1}, \ldots, x_{n}+i y_{n}\right)\right|$ is subharmonic in every pair $\left(x_{j}, y_{j}\right)$, and, as we can see later, in all the variables.

Example 2.6.1.4 Every harmonic function is subharmonic, as follows from Theorem 2.4.1.6. (Mean Value).

Theorem 2.6.1.1 (Elementary Properties) The following holds:
sh1) if $u \in S H(D)$, then $C u \in S H(D)$ for any constant $C \geq 0$;
sh2) if $u_{1}, u_{2} \in S H(D)$, then $u_{1}+u_{2}, \max \left[u_{1}, u_{2}\right] \in S H(D)$;
sh3) suppose $u_{n} \in S H(D), n=1,2, \ldots$, and the sequence converges to $u$ monotonically decreasing or uniformly on every compact set in $D$. Then $u \in$ $S H(D)$;
sh4) suppose $u(x, y) \in S H\left(D_{1}\right)$ for all $y \in D_{2}$, and be upper semicontinuous in $D_{1} \times D_{2}$. Let $\mu$ be a measure in $D_{2}$ such that $\mu\left(D_{2}\right)<\infty$. Then the function $u(x):=\int u(x, y) \mu(d y)$ is subharmonic in $D_{1}$.
sh5) let $V \in S O(m)$ be an orthogonal transformation of the space $\mathbb{R}^{m}$ and $u \in$ $S H\left(\mathbb{R}^{m}\right)$. Then $u(V \bullet) \in S H\left(\mathbb{R}^{m}\right)$.

All the assertions follow directly from the definition of subharmonic functions, properties of semicontinuous functions and properties of the Lebesgue integral. For a detailed proof see, e.g., [HK, Ch. 2].

Theorem 2.6.1.2 (Maximum Principle) Let $u \in S H(D), G \subset \mathbb{R}^{m}$ and $u(x) \not \equiv$ const. Then the inequality

$$
u(x)<\sup _{x^{\prime} \in \partial D} \limsup _{y \rightarrow x^{\prime}, y \in D} u(y), x \in D
$$

holds.
I.e., the maximum is not attained inside the domain.

The assertion follows from (2.6.1.2) and the upper semicontinuity of $u(x)$. For details see [HK, Ch. 2].

Let $K \Subset D$ be a compact set with nonempty interior $\stackrel{\circ}{K}$, and let $f_{n}$ be a decreasing sequence of functions continuous in $K$ that tends to $u \in S H(D)$. Such a sequence exists by Theorem 2.1.2.9. (The second criterion of semicontinuity).

Consider a sequence $\left\{H\left(x, u_{n}\right)\right\}$ of functions which are harmonic in $\stackrel{\circ}{K}$ and $\left.H\right|_{\partial K}=f_{n}$. The sequence converges monotonically to a function $H(x)$ harmonic in $\stackrel{\circ}{K}$ by Theorem 2.3.4.3. (Connection between convergences), Theorem 2.4.1.8. (Uniform and $\mathcal{D}^{\prime}$-convergences) and Theorem 2.6.1.2. The limit depends only on $u$ as one can see, i.e., it does not depend on the sequence $f_{n}$. This harmonic function $H(x):=H(x, u, K)$ is called the least harmonic majorant of $u$ in $K$.

This name is justified because of the following
Theorem 2.6.1.3. (Least Harmonic Majorant) Let $u \in S H(D)$. Then for any $K \Subset$ $D, u(x) \leq H(x, u, K), x \in K$. If $h(x)$ is harmonic in $\stackrel{\circ}{K}$ and satisfies the condition $h(x) \geq u(x), x \in \stackrel{\circ}{K}$, then $H(x, u, K) \leq h(x), x \in \stackrel{\circ}{K}$.

For the proof see [HK, Ch. 3].
2.6.2 Let us study properties of the mean values of subharmonic functions. Let $\mathcal{M}(x, r . u)$ be defined by (2.6.1.1) and $\mathcal{N}(x, r, u)$ by

$$
\mathcal{N}(x, r, u):=\frac{1}{\omega_{m} r^{m}} \int_{K_{x, r}} u(y) d y
$$

where $\omega_{m}$ is the volume of the ball $K_{0,1}$.

Theorem 2.6.2.1 (Properties of Mean Values) The following holds:
me1) $\mathcal{M}(x, r, u)$ and $\mathcal{N}(x, r, u)$ non-decreases in $r$ monotonically;
me2) $u(x) \leq \mathcal{N}(x, \bullet) \leq \mathcal{M}(x, \bullet)$;
me3) $\lim _{r \rightarrow 0} \mathcal{M}(x, r, u)=\lim _{r \rightarrow 0} \mathcal{N}(x, r, u)=u(x)$.
Proof. For simplicity let us prove me1) for $m=2$. We have

$$
\mathcal{M}\left(z_{0},|z|, u\right)=\frac{1}{2 \pi} \int_{0}^{2 \pi} u\left(z_{0}+z e^{i \phi}\right) d \phi
$$

Since $u(z, \phi):=u\left(z_{0}+z e^{i \phi}\right)$ is a family of subharmonic functions that satisfies the condition sh4) of Theorem 2.6.1.1, $\mathcal{M}\left(z_{0},|z|, u\right)$ is subharmonic in $z$ on any $K_{0, r}$. By Maximum Principle (Theorem 2.6.1.2) we have

$$
\mathcal{M}\left(z_{0}, r_{1}, u\right)=\max _{S_{0}, r_{1}} \mathcal{M}\left(z_{0},|z|, u\right) \leq \max _{S_{0, r_{2}}} \mathcal{M}\left(z_{0},|z|, u\right)=\mathcal{M}\left(z_{0}, r_{2}, u\right)
$$

for $r_{1}<r_{2}$.
Monotonicity of $\mathcal{N}(x, r, u)$ follows from the equality

$$
\begin{equation*}
\mathcal{N}(x, r, u)=m \int_{0}^{1} s^{m-1} \mathcal{M}(x, r s . u) d s \tag{2.6.2.1}
\end{equation*}
$$

and monotonicity of $\mathcal{M}(x, r, u)$.
The property me2) follows now from the definition of a subharmonic function and (2.6.2.1).

Let us prove me3). Let $M(u, x, r)$ be defined by (2.1.1.1). We have

$$
\mathcal{M}(x, r, u) \leq M(u, x, r) \quad \text { and } \quad M(u, x, r) \rightarrow u(x)
$$

because of upper semicontinuity of $u(x)$. Thus me2) implies me3).
It is clear from me2) that a subharmonic function is locally summable. From me3) we have the corollary

Theorem 2.6.2.2 (Uniqueness of subharmonic function) If $u, v \in S H(D)$ and $u=$ $v$ almost everywhere, then $u \equiv v$.

Let $\alpha(t)$ be defined by the equality (2.3.1.1), $\alpha_{\epsilon}(x)$ by (2.3.1.3). For a Borel set $E$ let

$$
E^{\epsilon}:=\{x: \exists y \in E:|x-y|<\epsilon\} .
$$

This is the $\epsilon$-extension of $E$; this is, of course, an open set. For an open set $D$ we set

$$
D^{-\epsilon}:=\bigcup_{E^{\epsilon} \subset D} E^{\epsilon} .
$$

This is the maximal set such that its $\epsilon$-extension is a subset of $D$. One can see that $D^{-\epsilon}$ is not empty for small $\epsilon$ and $D^{-\epsilon} \uparrow D$ when $\epsilon \downarrow 0$. Therefore for any $D_{1} \Subset D$ there exists $\epsilon$ such that $D_{1} \Subset D^{-\epsilon}$.

For $u \in S H(D)$ set

$$
\begin{equation*}
u_{\epsilon}(x):=\int u(x+y) \alpha_{\epsilon}(y) d y \tag{2.6.2.2}
\end{equation*}
$$

which is defined in $D^{-\epsilon}$.
Theorem 2.6.2.3 (Smooth Approximation) The following holds:
ap1) $\quad u_{\epsilon}$ is an infinitely differentiable subharmonic function in any open set $D_{1} \subset$ $D^{-\epsilon}$;
ap2) $\quad u_{\epsilon} \downarrow u(x)$ while $\epsilon \downarrow 0$ for all $x \in D$.
Proof. The property ap1) follows from sh4) (Theorem 2.6.1.1) and the following equality that one can obtain from (2.6.2.2):

$$
\begin{equation*}
u_{\epsilon}(x)=\int u(y) \alpha_{\epsilon}(x-y) d y \tag{2.6.2.3}
\end{equation*}
$$

Exercise 2.6.2.1 Prove this.
Let us prove ap2). From (2.6.2.2) we obtain

$$
\begin{equation*}
u_{\epsilon}(x)=\int_{0}^{1} \alpha(s) s^{m-1} \mathcal{M}(x, \epsilon s, u) d s \tag{2.6.2.4}
\end{equation*}
$$

It follows from the property me1) (Theorem 2.6.2.1) that $u_{\epsilon_{1}} \leq u_{\epsilon_{2}}$ while $\epsilon_{1}<\epsilon_{2}$. Now we pass to the limit in (2.6.2.4). Using me3) we have $\mathcal{M}(x, \epsilon s, u) \downarrow u(x)$. We can pass to the limit under the integral because of Theorem 2.2.2.2. Thus

$$
\lim _{\epsilon \downarrow 0} u_{\epsilon}(x)=\int_{0}^{1} \alpha(s) s^{m-1} u(x) d s=u(x) .
$$

Theorem 2.6.2.4 (Symmetry of $u_{\epsilon}$ ) If $u(x)$ depends only on $|x|$ then $u_{\epsilon}$ depends only on $|x|$.

Proof. Let $V \in S O(m)$ be a rotation of $\mathbb{R}^{m}$. Then

$$
u_{\epsilon}(V x)=\int u(y) \alpha_{\epsilon}(V x-y) d y
$$

Set $y=V y^{\prime}$ and change the variables. We obtain

$$
u_{\epsilon}(V x)=\int u\left(V y^{\prime}\right) \alpha_{\epsilon}\left(V\left(x-y^{\prime}\right)\right) d y
$$

Since $\alpha_{\epsilon}=\alpha_{\epsilon}(|x|)$ and $u=u(|x|), \alpha_{\epsilon}(V y)=\alpha_{\epsilon}(y)$ and $u(V y)=u(y)$. Thus $u_{\epsilon}(V x)=u_{\epsilon}(x)$ for any $V$ and thus $u_{\epsilon}(x)=u_{\epsilon}(|x|)$.
2.6.3 Since a subharmonic function is locally summable and defined uniquely by its values almost everywhere, every $u \in S H(D)$ corresponds to a (unique) distribution

$$
\langle u, \phi\rangle:=\int u(x) \phi(x) d x, \phi \in \mathcal{D}^{\prime}
$$

Theorem 2.6.3.1 (Necessary Differential Condition for Subharmonicity) If $u \in$ $S H(D)$, then $\Delta u$ is a positive distribution in $\mathcal{D}^{\prime}(D)$.

Proof. Suppose to begin that $u(x)$ has second continuous derivatives. By using (2.4.1.5) and (2.4.1.6) we can represent $u(x)$ in the form

$$
\begin{equation*}
u(x)=\mathcal{M}(x, r, u)+\int_{K_{x, r}} G\left(x, y, K_{x, r}\right) \Delta u(y) d y \tag{2.6.3.1}
\end{equation*}
$$

where $G$ is negative for all $r$. Suppose $\Delta u(x)<0$. Then it is negative in $K_{x, r}$ for some $r$. Thus the integral in (2.6.3.1) is positive and we obtain that $u(x)-$ $\mathcal{M}(x, r, u)>0$. This contradicts the subharmonicity of $u(x)$.

Now suppose $u(x)$ is an arbitrarily subharmonic function. Then $\Delta u_{\epsilon}(x) \geq 0$ for every $x \in D$ when $\epsilon$ is small enough.For each $x$ there is a neighborhood $D_{x}$ such that every $u_{\epsilon}$ defines a distribution from $\mathcal{D}^{\prime}\left(D_{x}\right)$. Hence $\Delta u_{\epsilon}(x)$ defines a positive distribution from $\mathcal{D}^{\prime}\left(D_{x}\right)$. Passing to the limit in $u_{\epsilon}$ when $\epsilon \downarrow 0$ we obtain in $\mathcal{D}^{\prime}\left(D_{x}\right)$ a distribution that is defined by function $u(x)$. Since the Laplace operator is continuous in any $\mathcal{D}^{\prime}$ (Theorem 2.3.4.2), $\Delta u>0$ in $\mathcal{D}^{\prime}\left(D_{x}\right)$. From Theorem 2.3.5.1 we obtain that $\Delta u$ is a positive distribution in $\mathcal{D}^{\prime}(D)$.

The distribution $\Delta u$ can be realized as a measure by Theorem 2.3.2.2. The measure $\left(\theta_{m}\right)^{-1} \Delta u$ is called the Riesz measure of the subharmonic function $u$.

Theorem 2.6.3.2 (Subharmonicity and Convexity) Let $u(|x|)$ be subharmonic in $x$ on $K_{0, R}$. Then $u(r)$ is convex with respect to $-r^{2-m}$ for $m \geq 3$ and with respect to $\log r$ for $m=2$.

Proof. By Theorem 2.6.2.4, $u_{\epsilon}(x)$ depends on $|x|$ only, i.e., $u_{\epsilon}(x)=u_{\epsilon}(|x|)$, and the function $u_{\epsilon}(r)$ is smooth. Passing to the spherical coordinates we obtain

$$
\Delta u_{\epsilon}=\frac{1}{r^{m-1}} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r} u_{\epsilon}(r) \geq 0
$$

By changing variables, $r=e^{v}$ for $m=2$ or $r=(-v)^{\frac{1}{2-m}}$ for $m \geq 3$, we obtain $\left[u_{\epsilon}(r(v)]^{\prime \prime} \geq 0\right.$, i.e., $u_{\epsilon}(r(v))$ is convex in $v$.

Passing to the limit on $\epsilon \downarrow 0$ we obtain that $u(r(v))$ is convex too, as a monotonic limit of convex functions.
2.6.4 Now we will consider the connection between subharmonicity and potentials.

Theorem 2.6.4.1 (Subharmonicity of $-\Pi)-\Pi(x, \mu, D) \in S H(D)$
It is because of GPo1) and GPo3) (Theorem 2.5.1.1).

The following theorem is inverse to Theorem 2.6.3.1.
Theorem 2.6.4.2 (Sufficient Differential Condition of Subharmonicity) Let $\Delta u \in$ $\mathcal{D}^{\prime}(D)$ be a positive distribution. Then there exists $u_{1} \in S H(D)$ that realizes $u$.

Proof. Set $\mu:=\theta_{m}^{-1} \Delta u$. Let $\Omega_{1} \Subset \Omega \Subset D$ and $\Pi\left(x, \mu_{\Omega}\right)$ be the Newtonian (or logarithmic) potential of $\left.\mu\right|_{\Omega}$. By GPo5) (Theorem 2.5.1.1) the sum $H:=u+\Pi$ is a harmonic distribution in $\mathcal{D}^{\prime}\left(\Omega_{1}\right)$. Hence there exists a "natural" harmonic function $H_{1}$ that realizes $H$ (Theorem 2.4.1.1). Thus the function $u_{1}:=H_{1}-\Pi \in S H\left(\Omega_{1}\right)$ and realizes $u$ in $\mathcal{D}^{\prime}(\Omega)$. Since $\Omega$ and $\Omega_{1}$ can be chosen such that a neighborhood of any $x \in D$ belongs to $\Omega_{1}$, the assertion holds for $D$.

By the way, we showed in this theorem that every subharmonic function can be represented inside its domain of subharmonicity as a difference of a harmonic function and a Newton potential. Thus all the smooth properties of a subharmonic function depend on the smooth properties of the potential only because any harmonic function is infinitely differentiable.

The following representation determines the harmonic function completely.
Theorem 2.6.4.3 (F. Riesz representation) Let $u \in S H(D)$ and let $K$ be a compact Lipschitz subdomain of $D$. Then

$$
u(x)=H(x, u, K)-\Pi\left(x, \mu_{u}, K\right)
$$

where $\mu_{u}$ is the Riesz measure of $u$ and $H(x, u, K)$ the least subharmonic majorant.
Proof. We can prove as above that the function $H(x):=u(x)+\Pi\left(x, \mu_{u}, K\right)$ is harmonic in $\stackrel{\circ}{K}$. Since $H(x) \geq u(x)$ we have $H(x) \geq H(x, u, K)$. So we need the reverse inequality.

Let us write the same equality for $u_{\epsilon}$ that is smooth.

$$
u_{\epsilon}:=H\left(x, u_{\epsilon}\right)-\Pi\left(x, \mu_{u_{\epsilon}}, K\right) .
$$

Passing to the limit as $\epsilon \downarrow 0$ we obtain

$$
u(x)=H(x, u, K)-\lim _{\epsilon \downarrow 0} \Pi\left(x, \mu_{u_{\epsilon}}, K\right)
$$

and the potentials converge because other summands converge. By Gpo5)

$$
\lim _{\epsilon \downarrow 0} \Pi\left(x, \mu_{u_{\epsilon}}, K\right) \geq \Pi\left(x, \mu_{u}, K\right)
$$

Hence $H(x) \leq H(x, u, K)$.
2.6.5 In this item we will consider subharmonic functions in the ball $K_{R}:=K_{0, R}$ which are harmonic in some neighborhood of the origin and write $u \in S H(R)$.

Set

$$
\begin{align*}
M(r, u) & :=\max \{u(x):|x|=r\} \\
\mu(r, u) & :=\mu_{u}\left(K_{r}\right) \\
\mathcal{M}(r, u) & :=\mathcal{M}(0, r, u) \\
N(r, u) & :=A(m) \int_{0}^{r} \frac{\mu(t, u)}{t^{m-1}} d t, \text { where } A(m)=\max (1, m-2) . \tag{2.6.5.1}
\end{align*}
$$

Theorem 2.6.5.1 (Jensen-Privalov) For $u \in S H(R)$,

$$
\begin{equation*}
\mathcal{M}(r, u)-u(0)=N(r, u), \text { for } 0<r<R \tag{2.6.5.2}
\end{equation*}
$$

Proof. By Theorem 2.6.4.3 we have

$$
u(x)=\frac{1}{\sigma_{m} r} \int_{|y|=r} u(y) \frac{r^{2}-|x|^{2}}{|x-y|^{m}} d s_{y}+\int_{K_{r}} G\left(x, y, K_{r}\right) \mu(d y) .
$$

For $x=0$ we obtain

$$
u(0)= \begin{cases}-\int_{0}^{r}\left(\frac{1}{t^{m-2}}-\frac{1}{r^{m-2}}\right) \mu(d t, u)+\mathcal{M}(r, u), & \text { for } m \geq 3 \\ -\int_{0}^{r} \log \frac{r}{t} \mu(d t, u)+\mathcal{M}(r, u), & \text { for } m=2\end{cases}
$$

Integrating by parts gives

$$
u(0)-\mathcal{M}(r, u)= \begin{cases}-\left.\mu(t, u)\left(\frac{1}{t^{m-2}}-\frac{1}{r^{m-2}}\right)\right|_{0} ^{r}+(m-2) \int_{0}^{r} \frac{\mu(t, u)}{t^{m-1}} d t, & \text { for } m \geq 3  \tag{2.6.5.3}\\ -\left.\mu(t, u) \log \frac{r}{t}\right|_{0} ^{r}+\int_{0}^{r} \frac{\mu(t, u)}{t} d t, & \text { for } m=2\end{cases}
$$

We have $\mu(t, u)=0$ for small $t$ because of harmonicity of $u(x)$. Thus (2.6.5.3) implies (2.6.5.2).

Theorem 2.6.5.2 (Convexity of $M(r, u)$ and $\mathcal{M}(r, u)$ ) These functions increase monotonically and are convex with respect to $\log r$ for $m=2$ and $-r^{2-m}$ for $m \geq 3$.

Proof. Consider the case $m=2$. Set $M(z):=\max _{\phi} u\left(z e^{i \phi}\right)$. One can see that $M(r)=M(r, u)$.

Let $u$ be a continuous subharmonic function. Then $M(z)$ is subharmonic (Theorem 2.6.1.1, sh5) and continuous because the family $\left\{u_{\phi}(z):=u\left(z e^{i \phi}\right)\right\}$ is uniformly continuous. The function $M(z)$ depends only on $|z|$. Thus it is convex with respect to $\log r$ by Theorem 2.6.3.2.

Let $u(z)$ be an arbitrarily subharmonic function and $u_{\epsilon} \downarrow u$ while $\epsilon \downarrow 0$. Then $M\left(r, u_{\epsilon}\right) \downarrow M(r, u)$ by Proposition 2.1.2.7 and is convex with respect to $\log r$ by sh3), Theorem 2.6.1.1.

If $m \geq 3$ you should consider the function $M(x):=\max _{|y|=|x|} u\left(V_{y} x\right)$ where $V_{y}$ is a rotation of $\mathbb{R}^{m}$ transferring $x$ into $y$.

The convexity of $\mathcal{M}(r, u)$ is proved analogously.

Exercise 2.6.5.1 Prove it.
The monotonicity of $M(r, u)$ follows from the Maximum Principle (Theorem 2.6.1.2). The monotonicity of $\mathcal{M}(r, u)$ was proved in Theorem 2.6.2.1.

The following classical assertion is a direct corollary of Theorem 2.6.5.2.
Theorem 2.6.5.3 (Three Circles Theorem of Hadamard) Let $f(z)$ be a holomorphic function in the disc $K_{R}$ and let $M_{f}(r)$ be its maximum on the circle $\{|z|=r\}$. Then

$$
M_{f}(r) \leq\left(\left[M_{f}\left(r_{1}\right)\right]^{\log \frac{r_{2}}{r}}\left[M_{f}\left(r_{2}\right)\right]^{\log \frac{r}{r_{1}}}\right)^{\frac{1}{\log \frac{v_{2}}{r_{1}}}}
$$

for $0<r_{1} \leq r \leq r_{2}<R$.
For the proof you should write down the condition of convexity with respect to $\log r$ of the function $\log M_{f}(r)$ which is the maximum of the subharmonic function $\log |f(z)|$.

Exercise 2.6.5.2 Do this.
For details see [PS, Part I, Sec. III, Ch. 6, Problem 304].

### 2.7 Sequences of subharmonic functions

2.7.1 We will formulate the following analogue for the Montel theorem of normal families of holomorphic functions.

The family

$$
\begin{equation*}
\left\{u_{\alpha}, \alpha \in A\right\} \subset S H(D) \tag{2.7.1.1}
\end{equation*}
$$

is called precompact in $\mathcal{D}^{\prime}(D)$ if, for any sequence $\left\{\alpha_{n}, n=1,2, \ldots\right\} \subset A$, there exists a subsequence $\alpha_{n_{j}}, j=1,2, \ldots$ and a function $u \in S H(D)$ such that $u_{\alpha_{n_{j}}} \rightarrow$ $u$ in $\mathcal{D}^{\prime}(D)$.

Example 2.7.1.1 $u_{\alpha}:=\log |z-\alpha|,|\alpha|<1$ form a precompact family.
Example 2.7.1.2 $u_{\alpha}:=\log \left|f_{\alpha}\right|$ where $\left\{f_{\alpha}\right\}$ is a family of holomorphic functions bounded in a domain $D$ form a precompact family.

A criterion of precompactness is given by
Theorem 2.7.1.1 (Precompactness in $\mathcal{D}^{\prime}$ ) A family (2.7.1.1) is precompact iff the following conditions hold:
comp1) for any compact set $K \subset D$ a constant $C(K)$ exists such that

$$
\begin{equation*}
u_{\alpha}(x) \leq C(K) \tag{2.7.1.2}
\end{equation*}
$$

for all $\alpha \in A$ and $x \in K$;
comp2) there exists a compact set $K_{1} \Subset D$ such that

$$
\begin{equation*}
\inf _{\alpha \in A} \max \left\{u_{\alpha}(x): x \in K_{1}\right\}>-\infty \tag{2.7.1.3}
\end{equation*}
$$

For the proof see [Hö, Thm. 4.1.9].
Theorem 2.7.1.2 Let $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}\left(K_{R}\right)$. Then $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}\left(S_{r}\right)$ for any $r<R$.
Proof. We have $\mu_{n} \rightarrow \mu$. Let us choose $R_{1}$ such that $r<R_{1}<R$. Then

$$
u_{n}(x)=H\left(x, u_{n}, K_{R_{1}}\right)-\Pi\left(x, \mu_{n}, K_{R_{1}}\right)
$$

by the F. Riesz theorem (Theorem 2.6.4.3).
Now, we have $\Pi\left(x, \mu_{n}, K_{R_{1}}\right) \rightarrow \Pi\left(x, \mu, K_{R_{1}}\right)$ in $\mathcal{D}^{\prime}\left(R_{1}\right)$ by GPo6), Theorem 2.5.1.1. Thus $H\left(x, u_{n}, K_{R_{1}}\right) \rightarrow H\left(x, u, K_{R_{1}}\right)$ in $\mathcal{D}^{\prime}\left(R_{1}\right)$.

By Theorem 2.4.1.8, $H\left(x, u_{n}, K_{R_{1}}\right) \rightarrow H\left(x, u, K_{R_{1}}\right)$ uniformly on any compact set in $K_{R_{1}}$, in particular, on $S_{r}$. Hence $H\left(x, u_{n}, K_{R_{1}}\right) \rightarrow H\left(x, u, K_{R_{1}}\right)$ in $\mathcal{D}^{\prime}\left(S_{r}\right)$. Also $\Pi\left(x, \mu_{n}, K_{R_{1}}\right) \rightarrow \Pi\left(x, \mu, K_{R_{1}}\right)$ in $\mathcal{D}^{\prime}\left(S_{r}\right)$ by GPo6), Theorem 2.5.1.1. Hence, $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}\left(S_{r}\right)$.

We say that a sequence $f_{n}$ of locally summable functions converges in $L_{\text {loc }}$ to a locally summable function $f$ if for any $x \in D$ there exists a neighborhood $V \ni x$ such that $\int_{V}\left|f_{n}-f\right| d x \rightarrow 0$.

Theorem 2.7.1.3 (Compactness in $L_{\mathrm{loc}}$ ) Under conditions of Theorem 2.7.1.1 the family (2.7.1.1) is precompact in $L_{\mathrm{loc}}$.

For the proof see [Hö, Thm. 4.1.9].
Theorem 2.7.1.4 Let $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}\left(K_{R}\right)$. Then $u_{n}^{+} \rightarrow u^{+}$in $\mathcal{D}^{\prime}\left(K_{R}\right)$.
This is because $u_{n}^{+}(x) \leq M, x \in K$, for all compact sets $K \Subset K_{R}$.
2.7.2 The following theorem shows that a subharmonic function is much more "flexible" than a harmonic or analytic function.

Theorem 2.7.2.1 Let $D \Subset \mathbb{R}^{m}$ be a Lipschitz domain and let $u \in S H(D)$ satisfy the condition $u(x)<C$ for $x \in D$. Then for any closed domain $D_{1} \Subset D$ there exists a function $\tilde{u}(x):=\tilde{u}\left(x, D_{1}\right)$ such that:
ext1) $u(x)=\tilde{u}(x)$ for $x \in D_{1}$;
ext2) $\quad \tilde{u}(x)=C$ for $x \in \partial D$;
ext3) $\quad \tilde{u} \in S H(D)$ and is harmonic in $D \backslash \bar{D}_{1}$;
ext4) $u(x) \leq \tilde{u}(x)$ for $x \in D$.
The function $\tilde{u}$ is defined uniquely.

Proof. We can suppose without loss of generality that $C=0$, because we can consider the function $u-C$.

Let $u(x)$ be continuous in $\bar{D}_{1}$. Consider a harmonic function $H(x)$ which is zero on $\partial D$ and $u(x)$ on $\partial D_{1}$. We have $H(x) \geq u(x)$ for $x \in D \backslash D_{1}$ because of Theorem 2.6.1.3. Set

$$
\tilde{u}(x)= \begin{cases}H(x), & x \in D \backslash D_{1} \\ u(x), & x \in D_{1}\end{cases}
$$

The function $\tilde{u}(x)$ is subharmonic in $D$. For $x \notin \partial D_{1}$ it is obvious, and for $x \in \partial D_{1}$ it follows from

$$
u(x)=\tilde{u}(x) \leq \mathcal{M}(x, r, u) \leq \mathcal{M}(x, r, \tilde{u})
$$

for $r$ small enough.
It is easy to check that all the assertions of the theorem are fulfilled for the function $\tilde{u}$.

Exercise 2.7.1.1 Check this.
Let $u(x)$ be an arbitrarily subharmonic function. Consider the family $u_{\epsilon}$ of smooth subharmonic functions that converges to $u(x)$ decreasing monotonically in a neighborhood of $\bar{D}_{1}$. The sequence $\widetilde{\left(u_{\epsilon}\right)}$ converges monotonically to a subharmonic function that has all the properties ext1)-ext4).

Theorem 2.7.2.2 (Continuity of $\tilde{\boldsymbol{\varphi}}$ ) Let $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}(D)$ and $u_{n}(x)<0$ in $D$. Then for any $K \Subset D$ with a smooth boundary $\partial K \widetilde{u_{n}}(\bullet, K) \rightarrow \tilde{u}(\bullet, K)$ in $\mathcal{D}^{\prime}(D)$.

For proving, we need the following auxiliary statement:
Theorem 2.7.2.3 Let $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}(D)$. Then for any smooth surface $S \Subset D$ and any function $g(x)$ continuous in a neighborhood of $S$ the assertion

$$
\begin{equation*}
\int_{S} u_{n}(x) g(x) d s_{x} \rightarrow \int_{S} u(x) g(x) d s_{x} \tag{2.7.2.1}
\end{equation*}
$$

holds.
Proof. Since $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}(D)$ also the Riesz measures of the functions converge. Hence $\mu_{n}(K) \leq C(K)$ for some $K \ni S$. Thus, for the sequence of potentials $\Pi\left(x, \mu_{n}\right)$, we have

$$
\int_{S} \Pi\left(x, \mu_{n}\right) g(x) d s_{x}=\int \mu_{n}(d y) \int_{S} \frac{g(x) d s_{x}}{|x-y|^{m-2}}
$$

The inner integral is a continuous function of $y$ as can be seen by simple estimates. Thus the assertion (2.7.2.1) holds for potentials. Now, one can represent $u_{n}$ in the form

$$
u_{n}(x)=H_{n}(x)-\Pi\left(x, \mu_{n}\right)
$$

in $K$. The sequence $H_{n}$ convergences in $\mathcal{D}^{\prime}$ and, hence, uniformly on $S$. Thus (2.7.2.1) holds for every $u_{n}$.

Proof of Theorem 2.7.2.2. Let $\phi \in \mathcal{D}(D)$ and $\operatorname{supp} \phi \subset \stackrel{\circ}{K}$. Then

$$
\left\langle\widetilde{u_{n}}, \phi\right\rangle=\left\langle u_{n}, \phi\right\rangle \rightarrow\langle u, \phi\rangle=\langle\tilde{u}, \phi\rangle .
$$

Let $x \in D \backslash K$. Then

$$
\widetilde{u}_{n}(x)=\int_{\partial K} \frac{\partial G}{\partial n_{y}}(x, y) u_{n}(y) d s_{y}
$$

By Theorem 2.7.2.3, $\widetilde{u_{n}}(x) \rightarrow \tilde{u}(x)$ for $x \in D \backslash K$. The sequence $\widetilde{u}_{n}$ is precompact in $\mathcal{D}^{\prime}(D)$. Thus every limit $u_{0}$ of the $\widetilde{u}_{n}$ coincides with $\tilde{u}(x)$ in $\stackrel{\circ}{K}$ and in $D \backslash K$. Hence, $u_{0} \equiv \tilde{u}$ in $\mathcal{D}^{\prime}(D)$.
2.7.3 The property sh2), Theorem 2.6.1.1, shows that the maximum of any finite number of subharmonic functions is a subharmonic function too. However, it is not so if the number is not finite.
Example 2.7.3.1 Set $u_{n}(z)=\frac{1}{n} \log |z|, n=1,2 \ldots$ The functions $u_{n} \in S H\left(K_{1}\right)$. Taking the supremum in $n$ we obtain

$$
u(z)=: \sup _{n} u_{n}(z)= \begin{cases}0, & \text { for } z \neq 0 \\ -\infty & \text { for } z=0\end{cases}
$$

The function is not semicontinuous, thus it is not subharmonic. However, it differs from a subharmonic function on a set of zero capacity. The following theorem shows that this holds in general.

Theorem 2.7.3.1 (H. Cartan) Let a family $\left\{u_{\alpha} \in S H(D), \alpha \in A\right\}$ be bounded from above and $u(x):=\sup _{\alpha \in A} u_{\alpha}(x)$. Then $u^{*} \in S H(D)$ and the set $E:=\{x:$ $\left.u^{*}(x)>u(x)\right\}$ is a zero capacity set.

For proving this theorem we need an auxiliary assertion
Theorem 2.7.3.2 Let $\Pi\left(x, \mu_{n}, D\right)$ be a monotonically decreasing sequence of Green potentials and $\operatorname{supp} \mu_{n} \subset K \Subset D$. Then there exists a measure $\mu$ such that the inequality

$$
\lim _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \geq \Pi(x, \mu, D)
$$

holds for all $x \in D$ with equality outside some set of zero capacity.
Proof. The sequence $\Pi\left(x, \mu_{n}, D\right)$ converges monotonically and thus in $\mathcal{D}^{\prime}$ (Theorem 2.3.4.3). Then $\mu_{n} \rightarrow \mu$ in $\mathcal{D}^{\prime}$ (Theorem 2.2.4.2.) and thus in $C^{*}$ - topology (Theorem 2.3.4.4). By GPo5) (Theorem 2.5.1.1) we have

$$
\lim _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \geq \Pi(x, \mu, D)
$$

Suppose that the strict inequality holds on some set $E$ of a positive capacity. By Theorem 2.5.2.3 one can find a compact set $K \subset E$ such that $\mathbf{c a p}(K)>0$. Then there exists a measure $\nu$ concentrated on $E$ such that its potential $\Pi(x, \nu, D)$ is
continuous (Theorem 2.5.5.2). Thus we have

$$
\begin{aligned}
\int \Pi(x, \mu, D) \nu(d x)<\int \lim _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \nu(d x) & =\lim _{n \rightarrow \infty} \int \Pi\left(x, \mu_{n}, D\right) \nu(d x) \\
\lim _{n \rightarrow \infty} \int \Pi(x, \nu, D) \mu_{n}(d x)=\int \Pi(x, \nu, D) \mu(d x) & =\int \Pi(x, \mu, D) \nu(d x)
\end{aligned}
$$

The equalities use Theorem 2.2.2.2 (B. Levy), reciprocity law (GPo4), Theorem 2.5.1.1, $C^{*}$-convergence of $\mu_{n}$ and once more the reciprocity law, respectively. So we have a contradiction.

Proof of Theorem 2.7.3.1. Suppose that $u_{n}(x) \uparrow u(x)$. We can suppose also that $u_{n}<0$. For any domain $G \Subset D$ the sequence $\tilde{u}_{n}(x) \rightarrow u(x)$ for $x \in G$ (see Theorem 2.7.2.1), because $u_{n}(x)=\tilde{u}_{n}(x)$ for $x \in G$. Since $\tilde{u}_{n}=\Pi\left(x, \tilde{\mu}_{n}, D\right)$ for $x \in D, \tilde{u}(x)=\Pi(x, \tilde{\mu}, D)=u(x)$ for $x \in G$ and coincides with $\lim _{n \rightarrow \infty} u_{n}(x)$ outside some set $E_{G}$ of zero capacity. Consider a sequence of domains $G_{n}$ that exhaust $D$. Then $u(x)=\lim _{n \rightarrow \infty} u_{n}(x)$ outside the set $E:=\cup_{n=1}^{\infty} E_{G_{n}}$ which has zero capacity by capZ1) (see item 2.5.2).

Now let $\left\{u_{n}, n=1,2 \ldots\right\}$ be a general countable set that satisfies the conditions of the theorem. Then the sequence $v_{n}:=\max \left\{u_{k}: k=1,2 \ldots, n\right\} \in S H(D)$ and $v_{n} \uparrow u$. Applying the previous reasoning we obtain the assertion of the theorem also in this case.

Let $\left\{u_{\alpha}, \alpha \in A\right\}$ be an arbitrary set satisfying the condition of the theorem. By Theorem 2.1.3.2 (Choquet's Lemma) one can find a countable set $A_{0} \subset A$ such that

$$
\left(\sup _{A_{0}} u_{\alpha}\right)^{*}=\left(\sup _{A} u_{\alpha}\right)^{*} .
$$

Since $\sup _{A_{0}} u_{\alpha} \leq \sup _{A} u_{\alpha}$, we have

$$
E:=\left\{x:\left(\sup _{A} u_{\alpha}\right)^{*}>\sup _{A} u_{\alpha}\right\} \subset E_{0}:=\left\{x:\left(\sup _{A_{0}} u_{\alpha}\right)^{*}>\sup _{A_{0}} u_{\alpha}\right\} .
$$

Thus $\boldsymbol{c a p}(E) \leq \boldsymbol{c a p}\left(E_{0}\right)=0$.
Corollary of Theorem 2.7.3.1 is
Theorem 2.7.3.3 (H. Cartan + ) Let $\left\{u_{t}, t \in(0 ; \infty)\right\} \subset S H(D)$ be a bounded from above family, and $v:=\lim \sup _{t \rightarrow \infty} u_{t}$. Then $v^{*} \in S H(D)$ and the set $E:=\{x:$ $\left.v^{*}(x)>v(x)\right\}$ has zero capacity.

Proof. Set $u_{n}:=\sup _{t \geq n} u_{t}, E_{n}:=\left\{x:\left(u_{n}\right)^{*}>u_{n}\right\}, E:=\cup E_{n}$. Since $\operatorname{cap}\left(E_{n}\right)=$ $0, \operatorname{cap} E=0$ too.

Let $x \notin E$. Then

$$
v(x)=\lim _{n \rightarrow \infty} \sup _{t \geq n} u_{t}(x)=\lim _{n \rightarrow \infty}\left(u_{n}\right)^{*}(x)
$$

The function

$$
v^{*}:=\lim _{n \rightarrow \infty}\left(u_{n}\right)^{*}(x)
$$

is the upper semicontinuous regularization of $v(x)$ for all $x \in D$.
In spite of Example 2.7.3.1 we have
Theorem 2.7.3.4 (Sigurdsson's Lemma) [Si] Let $S \subset S H(D)$ be compact in $\mathcal{D}^{\prime}$. Then

$$
v(x):=\sup \{u(x): u \in S\}
$$

is upper semicontinuous
and, hence, subharmonic.
Proof. Note that

$$
u_{\epsilon}(x)=\langle u, \alpha(x-\bullet)\rangle
$$

(see (2.6.2.3), (2.3.2.1)); and it is continuous in $(u, x)$ with respect to the product topology on $\left(S H(D) \cap \mathcal{D}^{\prime}\right) \times \mathbb{R}^{m}$ (Theorem 2.3.4.6).

Let $x_{0} \in D, a \in \mathbb{R}$ and assume that $v\left(x_{0}\right)<a$. We have to prove that there exists a neighborhood $X$ of $x_{0}$ such that

$$
\begin{equation*}
v(x)<a, x \in X \tag{2.7.3.1}
\end{equation*}
$$

We choose $\delta>0$ such that $v\left(x_{0}\right)<a-\delta$. If $u^{0} \in S H(D)$ and $\epsilon$ is chosen sufficiently small, then

$$
u^{0}\left(x_{0}\right) \leq u_{\epsilon}^{0}\left(x_{0}\right)<a-\delta
$$

by Theorem 2.6.2.3 (Smooth Approximation).
Since $u_{\epsilon}(x)$ is continuous, there exists an open neighborhood $U_{0}$ of $u^{0}$ in $S H(D)$ and an open neighborhood $X_{0}$ of $x_{0}$ such that

$$
u_{\epsilon}(x)<a-\delta, u \in U_{0}, x \in X_{0} .
$$

The property ap2) (Theorem 2.6.2.3) implies

$$
\begin{equation*}
u(x)<a-\delta, u \in U_{0}, x \in X_{0} \tag{2.7.3.2}
\end{equation*}
$$

Since $u^{0}$ is arbitrary and $S$ is compact, there exists a finite covering $U_{1}, U_{2}, \ldots, U_{n}$ of $S$ and open neighborhoods $X_{1}, X_{2}, \ldots, X_{n}$ of $x_{0}$ such that (2.7.3.2) holds for all $(u, x): u \in U_{j}, x \in X_{j}, j=1, \ldots, n$. Set $X:=\cap_{j} X_{j}$. Then (2.7.3.1) holds.
2.7.4 Now we are going to connect $\mathcal{D}^{\prime}$-convergence to convergence outside a zero capacity set, the so-called quasi-everywhere convergence.

Theorem 2.7.4.1 ( $\mathcal{D}^{\prime}$ and Quasi-everywhere Convergence) Let $u_{n}, u \in S H(D)$ and $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}(D)$. Then $u(x)=\lim \sup _{n \rightarrow \infty} u_{n}(x)$ quasi-everywhere and $u(x)=\left(\limsup \sup _{n \rightarrow \infty} u_{n}(x)\right)^{*}$ everywhere in $D$.

For the proof we need the following assertion in the spirit Theorem 2.7.3.2.
Theorem 2.7.4.2 Let $\mu_{n} \rightarrow \mu$ in $\mathcal{D}^{\prime}(D)$ and $\operatorname{supp} \mu_{n} \subset K \Subset D$. Then

$$
\liminf _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \geq \Pi(\mu, D)
$$

with equality quasi-everywhere.
Proof. The inequality was in GPo5), Theorem 2.5.1.1.
Suppose the set

$$
E:=\left\{x: \liminf _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right)>\Pi(x, \mu, D)\right.
$$

has a positive capacity. By Theorem 2.5.2.3 one can find a compact set $K \subset E$ such that $\boldsymbol{\operatorname { c a p }}(K)>0$. By Theorem 2.5.5.2 one can find a measure $\nu$ concentrated on $K$ with continuous potential. As in the proof of Theorem 2.7.3.2 we have

$$
\begin{aligned}
& \int \Pi(x, \mu, D) \nu(d x)<\int \liminf _{n \rightarrow \infty} \Pi\left(x, \mu_{n}, D\right) \nu(d x) \leq \liminf _{n \rightarrow \infty} \int \Pi\left(x, \mu_{n}, D\right) \nu(d x) \\
= & \liminf _{n \rightarrow \infty} \int \Pi(x, \nu, D) \mu_{n}(d x)=\int \Pi(x, \nu, D) \mu(d x)=\int \Pi(x, \mu, D) \nu(d x) .
\end{aligned}
$$

The second inequality uses Theorem 2.2.2.3 (Fatou's Lemma). The equalities use the reciprocity law (GPo4), Theorem 2.5.1.1, $C^{*}$-convergence of $\mu_{n}$ and once more the reciprocity law, respectively. So we have a contradiction.

Proof of Theorem 2.7.4.1. Let $D_{1} \Subset D$. Then the sequence $u_{n}$ is bounded in $D_{1}$ by Theorem 2.7.1.1. We can assume that $u_{n}(x)<0$ for $x \in D_{1}$.

For any domain $G \Subset D_{1}$ the sequence $\tilde{u}_{n}(x, G) \rightarrow u(x)$ in $\mathcal{D}^{\prime}\left(D_{1}\right)$ by Theorem 2.7.2.2. We also have the equality $\tilde{u}_{n}=-\Pi\left(x, \tilde{\mu}_{n}, D_{1}\right)$. Thus $\tilde{\mu}_{n} \rightarrow \tilde{\mu}$ in $\mathcal{D}^{\prime}\left(D_{1}\right)$. By Theorem 2.7.4.2, $\liminf _{n \rightarrow \infty} \Pi\left(x, \tilde{\mu}_{n}, D_{1}\right)=\Pi\left(x, \tilde{\mu}, D_{1}\right)$ quasi-everywhere in $D_{1}$. Hence

$$
\begin{equation*}
\limsup _{n \rightarrow \infty} u_{n}=u \tag{2.7.4.1}
\end{equation*}
$$

quasi-everywhere in $G$ because $u_{n}(x)=\tilde{u}_{n}(x)$ for $x \in G$.
Consider a sequence of domains $G_{n}$ that exhaust $D$. Then (2.7.4.1) holds outside a set $E_{n}$ of zero capacity and (2.7.4.1) holds in $D$ outside the set $E:=\cup_{n=1}^{\infty} E_{n}$ which has zero capacity by capZ1) (see item 2.5.2), i.e., quasi-everywhere.
2.7.5 Now we connect the convergence of subharmonic functions in $\mathcal{D}^{\prime}$ to the convergence relative to the Carleson measure (see 2.5.4).

We say that a sequence of functions $u_{n}$ converges to a function $u$ relative to the $\alpha$-Carleson measure if the sets $E_{n}:=\left\{x:\left|u_{n}(x)-u(x)\right|>\epsilon\right\}$ possess the property

$$
\begin{equation*}
\alpha-\operatorname{mes}_{C} E_{n} \rightarrow 0 \tag{2.7.5.1}
\end{equation*}
$$

Theorem 2.7.5.1 ( $\mathcal{D}^{\prime}$ and $\alpha-$ mes $_{C}$ Convergences) Let $u_{n}, u \in S H(D)$ and $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}(D)$. Then for an every $\alpha>0$ and every domain $G \Subset D u_{n} \rightarrow u$ relative to the ( $\alpha+m-2$ )-Carleson measure.

For proving this theorem we need some auxiliary definitions and assertions.
Let $\mu$ be a measure in $\mathbb{R}^{m}$. We will call a point $x \in \mathbb{R}^{m}\left(\alpha, \alpha^{\prime}, \epsilon\right)$-normal with respect to the measure $\mu,\left(\alpha<\alpha^{\prime}\right)$ if the inequality

$$
\mu_{x}(t):=\mu\left(K_{x, t}\right)<\epsilon^{-\alpha^{\prime}} t^{\alpha+m-2}
$$

holds for all $t<\epsilon$.
Theorem 2.7.5.2 In any $\left(\alpha, \alpha^{\prime}, \epsilon\right)$-normal point the following inequality holds:

$$
\begin{aligned}
& -\int_{K_{z, \epsilon}}[\log |z-\zeta|-\log \epsilon] d \mu_{\zeta} \leq C \epsilon^{\alpha-\alpha^{\prime}}, \text { for } m=2 \\
& \int_{K_{x, t}}\left[|x-y|^{2-m}-\epsilon^{2-m}\right] d \mu_{y} \leq C \epsilon^{\alpha-\alpha^{\prime}}, \text { for } m \geq 3
\end{aligned}
$$

while $C=C(\alpha, m)$ depends on $\alpha$ and $m$ only.
Proof. Let us consider the case $m=2$. We have

$$
\int_{K_{z, \epsilon}} \log \frac{\epsilon}{|z-\zeta|} d \mu_{\zeta}=\int_{0}^{\epsilon} \log \frac{\epsilon}{t} d \mu_{z}(t)
$$

Integrating by parts we obtain

$$
\begin{aligned}
\int_{K_{z, \epsilon}} \log \frac{\epsilon}{|z-\zeta|} d \mu_{\zeta} & =\left.\log \frac{\epsilon}{t} \mu_{z}(t)\right|_{0} ^{\epsilon}+\int_{0}^{\epsilon} \frac{\mu_{z}(t)}{t} d t \\
& \leq \epsilon^{-\alpha^{\prime}} \int_{0}^{\epsilon} t^{\alpha-1} d t=\frac{1}{\alpha} \epsilon^{\alpha-\alpha^{\prime}}
\end{aligned}
$$

Let us consider the case $m \geq 3$. We have

$$
\begin{aligned}
\int_{K_{x, t}}\left[|x-y|^{2-m}-\epsilon^{2-m}\right] d \mu_{y} & =\int_{0}^{\epsilon}\left(t^{2-m}-\epsilon^{2-m}\right) d \mu_{x}(t) \\
& =\left.\left(t^{2-m}-\epsilon^{2-m}\right) \mu_{x}(t)\right|_{0} ^{\epsilon}+(m-2) \int_{0}^{\epsilon} \frac{\mu_{x}(t)}{t^{m-2}} d t \\
& \leq \frac{m-2}{\epsilon^{\alpha^{\prime}}} \int_{0}^{\epsilon} t^{\alpha-1} d t=\frac{m-2}{\alpha} \epsilon^{\alpha-\alpha^{\prime}}
\end{aligned}
$$

Theorem 2.7.5.3 (Ahlfors-Landkof Lemma) Let a set $E \subset \mathbb{R}^{m}$ be covered by balls with bounded radii such that every point is a center of a ball. Then there exists an at most countable subcovering of the same set with maximal multiplicity $\mathrm{cr}=\mathrm{cr}(\mathrm{m})$.
I.e., every point of $E$ is covered no more than $c r$ times. Let us note that $c r(2)=6$.

For the proof see [La, Ch. III, § 4, Lem. 3.2].

Theorem 2.7.5.4 Let $K \Subset D$. The set $E:=E\left(\alpha, \alpha^{\prime}, \epsilon, \mu\right)$ of points that belong to $K$ and are not $\left(\alpha, \alpha^{\prime}, \epsilon\right)$-normal with respect to $\mu$ satisfies the condition

$$
\begin{equation*}
(\alpha+m-2)-\operatorname{mes}_{C} E \leq c r(m) \epsilon^{\alpha^{\prime}} \mu\left(K^{\epsilon}\right) \tag{2.7.5.2}
\end{equation*}
$$

where $K^{\epsilon}$ is the $2 \epsilon$-extension of $K$.
Proof. Let $x \in E$. Then there exists $t_{x}$ such that

$$
\mu_{x}\left(t_{x}\right) \geq t_{x}^{\alpha+m-2} \epsilon^{-\alpha^{\prime}}
$$

Thus every point of $E$ is covered by a ball $K_{x, t_{x}}$. By the Ahlfors-Landkof lemma (Theorem 2.7.5.3) one can find a no more than $\operatorname{cr}(m)$-multiple subcovering $\left\{K_{x_{j}, t_{x_{j}}}\right\}$. Then we have

$$
\sum_{j} t_{x_{j}}^{\alpha+m-2} \leq c r(m) \epsilon^{\alpha^{\prime}} \mu\left(K^{\epsilon}\right)
$$

By definition of the Carleson measure we obtain (2.7.5.2).
Theorem 2.7.5.5 Let $\mu_{n} \rightarrow \mu$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$ and $\operatorname{supp} \mu_{n} \subset K \Subset \mathbb{R}^{m}$. Then for every $\alpha>0$ and $G \Subset \mathbb{R}^{m}, \Pi\left(x, \mu_{n}\right) \rightarrow \Pi(x, \mu)$ relative to the $(\alpha+m-2)$-Carleson measure.

Proof. Let $m=2$. Set

$$
\log _{\epsilon}|z-\zeta|= \begin{cases}\log |z-\zeta|, & \text { for }|z-\zeta|>\epsilon \\ \log \epsilon, & \text { for }|z-\zeta| \leq \epsilon\end{cases}
$$

This function is continuous for $(z, \zeta) \in K \times K$.
Set $\nu_{n}:=\mu_{n}-\mu$. Then we have

$$
\begin{gathered}
-\int \log |z-\zeta| \mu_{n}(d \zeta)+\int \log |z-\zeta| \mu(d \zeta)=-\int \log |z-\zeta| \nu_{n}(d \zeta) \\
\quad=-\int \log _{\epsilon}|z-\zeta| \nu_{n}(d \zeta)-\int_{K_{z, \epsilon}}[\log |z-\zeta|-\log \epsilon] \nu_{n}(d \zeta)
\end{gathered}
$$

The function $\log _{\epsilon}|z-\zeta|$ is continuous in $\zeta$ uniformly over $z \in K$. Thus the sequence

$$
\Pi_{\epsilon}(z):=\int \log _{\epsilon}|z-\zeta| \nu_{n}(d \zeta)
$$

converges uniformly to zero on $K$. Suppose now that

$$
z \notin E\left(\alpha, \alpha^{\prime}, \epsilon, \mu\right) \cup E\left(\alpha, \alpha^{\prime}, \epsilon, \mu_{n}\right)
$$

i.e., it is an $\left(\alpha, \alpha^{\prime}, \epsilon\right)$ - normal point for $\mu$ and $\mu_{n}$. By Theorem 2.7.5.2 we have

$$
\int_{K_{z, \epsilon}}[\log |z-\zeta|-\log \epsilon] \nu_{n}(d \zeta)<2 C \epsilon^{\alpha-\alpha^{\prime}}
$$

Thus for sufficiently large $n>n_{0}(\epsilon)$,

$$
\left|\Pi\left(z, \mu_{n}\right)-\Pi(z, \mu)\right|=\left|\int \log \right| z-\zeta\left|\mu_{n}(d \zeta)-\int \log \right| z-\zeta|\mu(d \zeta)|<\delta=\delta(\epsilon)
$$

while $z \notin E\left(\alpha, \alpha^{\prime}, \epsilon, \mu\right) \cup E\left(\alpha, \alpha^{\prime}, \epsilon, \mu_{n}\right):=E_{n}(\epsilon)$.
By Theorem 2.7.5.3 the Carleson measure of $E_{n}(\epsilon)$ satisfies the inequality

$$
\alpha-\operatorname{mes}_{C} E_{n}(\epsilon) \leq c r(m) \epsilon^{\alpha^{\prime}}\left[\mu(K)+\mu_{n}(K)\right] \leq C \epsilon^{\alpha^{\prime}}:=\gamma(\epsilon)
$$

where $C=C(K)$ does not depend on $n$ because $\mu_{n}(K)$ are bounded uniformly.
Hence, for any $\epsilon>0$ the set

$$
E_{n}^{\prime}(\epsilon):=\left\{z:\left|\Pi\left(z, \mu_{n}\right)-\Pi(z, \mu)\right|>\delta(\epsilon)\right\}
$$

satisfies the condition

$$
\begin{equation*}
\alpha-\operatorname{mes}_{C} E_{n}^{\prime}(\epsilon) \leq \gamma(\epsilon) \tag{2.7.5.3}
\end{equation*}
$$

while $n>n_{0}=n_{0}(\epsilon)$.
Let us show that $\Pi\left(z, \mu_{n}\right) \rightarrow \Pi(z, \mu)$ relative to $\alpha-\operatorname{mes}_{C}$ on $K$. Let $\gamma_{0}, \delta_{0}$ be arbitrarily small. One can find $\epsilon$ such that $\delta(\epsilon)<\delta_{0}, \gamma(\epsilon)<\gamma_{0}$. One can find $n_{0}=n_{0}(\epsilon)$ such that (2.7.5.3) is fulfilled. Now the set

$$
E_{n, \delta_{0}}:=\left\{z:\left|\Pi\left(z, \mu_{n}\right)-\Pi(z, \mu)\right|>\delta_{0}\right\}
$$

is contained in $E_{n}^{\prime}(\epsilon)$. Thus $\alpha-\operatorname{mes}_{C} E_{n, \delta_{0}}<\gamma_{0}$ and this implies the convergence relative to $\alpha-\operatorname{mes}_{C}$. An analogous reasoning works for $m \geq 3$.

Proof of Theorem 2.7.5.1. Let $u_{n} \rightarrow u$ in $\mathcal{D}^{\prime}$. One can assume that $u_{n}, u$ are potentials on any compact set (Theorem 2.7.2.2). Hence, by Theorem 2.7.5.5 it converges relative $(\alpha+m-2)-\operatorname{mes}_{C}$.

### 2.8 Scale of growth. Growth characteristics of subharmonic functions

2.8.1 Let $A$ be a class of nondecreasing functions $a(r), r \in(0, \infty)$ such that $a(r) \geq 0$ and $a(r) \rightarrow \infty$ when $r \rightarrow \infty$. The quantity

$$
\begin{equation*}
\rho[a]:=\limsup _{r \rightarrow \infty} \frac{\log a(r)}{\log r} \tag{2.8.1.1}
\end{equation*}
$$

is called the order of $a(r)$.

Suppose $\rho:=\rho[a]<\infty$. The number

$$
\begin{equation*}
\sigma[a]:=\limsup _{r \rightarrow \infty} \frac{a(r)}{r^{\rho}} \tag{2.8.1.2}
\end{equation*}
$$

is called the type number.
If $\sigma[a]=0$, we say $a(r)$ has minimal type. If $0<\sigma[a]<\infty, a(r)$ has normal type. If $\sigma[a]=\infty$, it has maximal type.

Example 2.8.1.1 Set $a(r):=\sigma_{0} r^{\rho_{0}}$. Then $\rho[a]=\rho_{0}, \sigma[a]=\sigma_{0}$.
Example 2.8.1.2 Set $a(r):=(\log r)^{-1} r^{\rho_{0}}$. Then $\rho[a]=\rho_{0}, \sigma[a]=0$.
Example 2.8.1.3 Set $a(r):=(\log r) r^{\rho_{0}}$. Then $\rho[a]=\rho_{0}, \sigma[a]=\infty$.
Theorem 2.8.1.1 (Convergence Exponent) The following equality holds:

$$
\begin{equation*}
\rho[a]=\inf \left\{\lambda: \int^{\infty} \frac{a(r) d r}{r^{\lambda+1}}<\infty\right\} \tag{2.8.1.3}
\end{equation*}
$$

If the integral converges for $\lambda=\rho[a], a(r)$ has minimal type.
Exercise 2.8.1.1 Prove this.
For the proof see, e.g., [HK, §4.2].
Example 2.8.1.4 Let $r_{j}, j=1,2, \ldots$ be a nondecreasing sequence of positive numbers. Let us concentrate the unit mass in every point $r_{j}$ and define a mass distribution

$$
n(E):=\left\{\text { the number points of the sequence }\left\{r_{j}\right\} \text { in } E\right\}, E \subset \mathbb{R} .
$$

Then

$$
\begin{equation*}
\int_{0}^{\infty} \frac{d n}{r^{\lambda}}=\sum_{1}^{\infty} \frac{1}{r_{j}^{\lambda}} \tag{2.8.1.4}
\end{equation*}
$$

The infimum of $\lambda$ for which the series in (2.8.1.4) converges is usually called the convergence exponent for the sequence $\left\{r_{j}\right\}$ [PS, Part I, Sec. 1, Ch. III, § 2]. Integrating by parts one can transform the integral in (2.8.1.4) to an integral of the form (2.8.1.3) where $a(r)=n((-\infty, r))$. Theorem 2.8.1.1 shows that the convergence exponent coincides with the order of this $a(r)$.

A function $\rho(r)$ is called a proximate order with respect to order $\rho$ if
po1) $\quad \rho(r) \geq 0$,
po2) $\lim _{r \rightarrow \infty} \rho(r)=\rho$,
po3) $\quad \rho(r)$ has a continuous derivative on $(0, \infty)$,
po4) $\quad \lim _{r \rightarrow \infty} r \log r \rho^{\prime}(r)=0$.

Two proximate orders $\rho_{1}(r)$ and $\rho_{2}(r)$ are called equivalent, if

$$
\begin{equation*}
\rho_{1}(r)-\rho_{2}(r)=o\left(\frac{1}{\log r}\right) . \tag{2.8.1.5}
\end{equation*}
$$

For $a \in A$ set

$$
\begin{equation*}
\sigma[a, \rho(r)]:=\limsup _{r \rightarrow \infty} \frac{a(r)}{r^{\rho(r)}} \tag{2.8.1.6}
\end{equation*}
$$

It is called a type number with respect to a proximate order $\rho(r)$. It is clear that this type number is the same for equivalent proximate orders.

Theorem 2.8.1.2 (Proper Proximate Order) Let $a \in A$ and $\rho[a]=\rho<\infty$. Then there exists a proximate order $\rho(r)$ such that

$$
\begin{equation*}
0<\sigma[a, \rho(r)]<\infty \tag{2.8.1.7}
\end{equation*}
$$

For the proof see [Le, Ch. 1, Sec. 12, Thm. 16].
If a proximate order satisfies the condition (2.8.1.7), we will call it the proper proximate order of $a(r)$ (p.p.o.). The function $r^{\rho(r)}$ inherits a lot of useful properties of the power function $r^{\rho}$.

Theorem 2.8.1.3 (Properties of P.O) The following holds:
ppo1) the function $V(r):=r^{\rho(r)}$ increases monotonically for sufficiently large values of $r$.
ppo2) for $q<\rho+1$,

$$
\int_{1}^{r} t^{\rho(t)-q} d t \sim \frac{r^{\rho(r)+1-q}}{\rho+1-q}
$$

and for $q>\rho+1$,

$$
\int_{r}^{\infty} t^{\rho(t)-q} d t \sim \frac{r^{\rho(r)+1-q}}{q-\rho-1}
$$

as $r \rightarrow \infty$.
ppo3) the function $L(r):=r^{\rho(r)-\rho}$ satisfies the condition

$$
\forall \delta>0, L(k r) / L(r) \rightarrow 1
$$

when $r \rightarrow \infty$ uniformly for $k \in\left[\frac{1}{\delta}, \delta\right]$.
Exercise 2.8.1.2 Prove these properties.
For the proof see, e.g., [Le, Ch. 2, Sec. 12]. The following assertion allows us to replace any p.o. with a smooth one.

Theorem 2.8.1.4 (Smooth P.O) Let $\rho(r)$ be an arbitrary p.o. There exists an infinitely differentiable equivalent p.o. $\rho_{1}(r)$ such that

$$
\begin{equation*}
r^{k} \log r \rho_{1}^{(k)}(r) \rightarrow 0, k=1,2, \ldots \tag{2.8.1.8}
\end{equation*}
$$

when $r \rightarrow \infty$.

Proof. Let $\alpha_{\epsilon}$ be defined by (2.3.1.3). Set $\epsilon:=0.5, p o(x):=\rho\left(e^{x}\right)$ and

$$
p o_{1}(x):=p o(n)+[p o(n+1)-p o(n)] \int_{n}^{x} \alpha_{0.5}(t+0.5) d t
$$

for $x \in[n, n+1)$. The function $p o_{1}(x)$ is continuous and infinitely differentiable due to properties of $\alpha_{\epsilon}$ and $p o_{1}(n)=p o(n)$ for $n=1,2, \ldots$ By property po3) of p.o. we have

$$
(n+1)|p o(n+1)-p o(n)| \leq \frac{n+1}{n} \max _{y \in[n, n+1]}\left|y \cdot p o^{\prime}(y)\right| \rightarrow 0
$$

as $n \rightarrow \infty$. Thus

$$
\max _{y \in[n, n+1]}\left|y \cdot p o_{1}^{(k)}(y)\right| \leq \text { const } \cdot(n+1)|p o(n+1)-p o(n)| \rightarrow 0
$$

as $n \rightarrow \infty$.
So $\rho_{1}(r):=p o_{1}(\log r)$ is a p.o. that satisfies (2.8.1.8). Let us show that it is equivalent to $\rho(r)$. Indeed

$$
\begin{aligned}
\left|p o(x)-p o_{1}(x)\right| & =\left|\int_{n}^{x}\left[p o(y)-p o_{1}(y)\right]^{\prime} y \frac{d y}{y}\right| \\
& \leq \max _{y \in[n, n+1]}\left[\left|y \cdot p o^{\prime}(y)\right|+\left|y \cdot p o_{1}^{\prime}(y)\right|\right] \log \frac{n+1}{n}=o\left(\frac{1}{x}\right)
\end{aligned}
$$

when $x \in[n, n+1]$ and $n \rightarrow \infty$.
We will further need (in 2.9.3) the following assertion.
Theorem 2.8.1.5. (A.A. Gol'dberg) Let $\rho(r) \rightarrow \rho$ be a p.o., and let $f(t)$ be a function that is locally summable on $(0, \infty)$ and such that

$$
\begin{equation*}
\lim _{t \rightarrow 0} t^{\rho+\delta} f(t)=\lim _{t \rightarrow \infty} t^{\rho+1+\gamma} f(t)=0 \tag{2.8.1.9}
\end{equation*}
$$

for some $0<\delta, \gamma<1$.
Then

$$
\begin{align*}
\lim _{r \rightarrow \infty} r^{-\rho(r)} \int_{c r}^{x}(r t)^{\rho(r t)} f(t) d t & =\int_{0}^{x} t^{\rho} f(t) d t \\
\lim _{r \rightarrow \infty} r^{-\rho(r)} \int_{x}^{\infty}(r t)^{\rho(r t)} f(t) d t & =\int_{x}^{\infty} t^{\rho} f(t) d t \tag{2.8.1.10}
\end{align*}
$$

for any $c>0$ and any $x \in(0, \infty)$.

Proof. Set

$$
I(r):=\int_{c r^{-1}}^{\infty} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} f(t) d t
$$

It will be enough to prove that

$$
\begin{equation*}
\lim _{r \rightarrow \infty} I(r)=\int_{0}^{\infty} t^{\rho} f(t) d t \tag{2.8.1.11}
\end{equation*}
$$

because both functions

$$
f_{0}(t, x):= \begin{cases}f(t), & \text { for } t \in(0, x) \\ 0 & \text { for } t \in[x, \infty)\end{cases}
$$

and $f_{\infty}(t, x):=f(t)-f_{0}(t, x)$ also satisfy the condition of the theorem.
Let us represent the integral as the following sum:

$$
\begin{equation*}
I(r):=\int_{c r^{-1}}^{\infty} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} f(t) d t=I_{1}(r, \epsilon)+I_{2}(r, \epsilon)+I_{3}(r, \epsilon), \tag{2.8.1.12}
\end{equation*}
$$

where

$$
\begin{aligned}
& I_{1}(r, \epsilon):=\int_{c r^{-1}}^{\epsilon} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} f(t) d t \\
& I_{2}(r, \epsilon):=\int_{\epsilon}^{\epsilon^{-1}} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} f(t) d t \\
& I_{3}(r, \epsilon):=\int_{\epsilon^{-1}}^{\infty} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} f(t) d t .
\end{aligned}
$$

We can represent $I_{2}(r, \epsilon)$ in the form

$$
I_{2}(r, \epsilon)=\int_{\epsilon}^{\epsilon^{-1}} \frac{L(r t)}{L(r)} t^{\rho} f(t) d t
$$

By ppo3) (Theorem 2.8.1.3),

$$
\begin{equation*}
\lim _{r \rightarrow \infty} I_{2}(r, \epsilon)=\int_{\epsilon}^{\epsilon^{-1}} t^{\rho} f(t) d t \tag{2.8.1.13}
\end{equation*}
$$

Let us estimate the "tails". From (2.8.1.9) we have

$$
|f(t)| \leq C t^{-\rho-\delta}
$$

for $0<t \leq \epsilon$ where $C$ does not depend on $\epsilon$ and

$$
|f(t)| \leq C t^{-\rho-1-\gamma}
$$

for $t \geq \epsilon^{-1}$. We have

$$
\left|I_{1}(r, \epsilon)\right| \leq C \int_{c r-1}^{\epsilon} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} t^{-\rho-\delta} d t:=C J_{1}(r, \epsilon)
$$

and

$$
\begin{equation*}
\limsup _{r \rightarrow \infty}\left|I_{1}(r, \epsilon)\right| \leq C \lim _{r \rightarrow \infty} J_{1}(r, \epsilon) \tag{2.8.1.14}
\end{equation*}
$$

Let us calculate the last limit.We perform the change $x=t r$ :

$$
J_{1}(r, \epsilon)=r^{-\rho(r)+\rho+\delta-1} \int_{c}^{\epsilon r} t^{-\rho(x)-(\rho+\delta)} d x
$$

Now we use ppo2) for $q=\rho+\delta$ and ppo3):

$$
\begin{aligned}
\lim _{r \rightarrow \infty} J_{1}(r, \epsilon) & =\frac{1}{1-\delta} \lim _{r \rightarrow \infty} \frac{(\epsilon r)^{\rho(\epsilon r)-(\rho+\delta)+1}}{r^{\rho(r)-(\rho+\delta)+1}} \\
& =\frac{\epsilon^{1-\delta}}{1-\delta} \lim _{r \rightarrow \infty} \frac{L(\epsilon r)}{L(r)}=\frac{\epsilon^{1-\delta}}{1-\delta} .
\end{aligned}
$$

Substituting in (2.8.1.14) we obtain

$$
\begin{equation*}
\limsup _{r \rightarrow \infty}\left|I_{1}(r, \epsilon)\right| \leq C \frac{\epsilon^{1-\delta}}{1-\delta} \tag{2.8.1.15}
\end{equation*}
$$

Analogously one can obtain

$$
\begin{equation*}
\limsup _{r \rightarrow \infty}\left|I_{3}(r, \epsilon)\right| \leq C \frac{\epsilon^{\gamma}}{\gamma} \tag{2.8.1.16}
\end{equation*}
$$

Using (2.8.1.13), (2.8.1.15) and (2.8.1.16) one can pass to the limit in (2.8.1.12) as $r \rightarrow \infty$, then let $\epsilon \rightarrow 0$ and obtain (2.8.1.11).
2.8.2 Let

$$
\begin{equation*}
u(x):=u_{1}(x)-u_{2}(x) \tag{2.8.2.1}
\end{equation*}
$$

where $u_{1}, u_{2} \in S H\left(\mathbb{R}^{m}\right), u_{1}(0)>-\infty, u_{2}(0)=0$ and $\mu_{1}:=\mu_{u_{1}}, \mu_{2}:=\mu_{u_{2}}$ are concentrated on disjoint sets.

Let $m=2, u_{j}(z):=\log \left|f_{j}(z)\right|, j=1,2$ where $f_{j}(z), j=1,2$ are entire functions. Then the function $u(z)=\log |f(z)|$, where $f(z):=f_{1}(z) / f_{2}(z)$, is meromorphic. The condition for masses means that $f_{1}$ and $f_{2}$ have no common zeros, $u_{2}(0)=0$ corresponds to $f_{2}(0)=1$ and $u_{1}(0)>-\infty$ means $f_{1}(0) \neq 0$.

The class of such functions is denoted as $\delta S H\left(\mathbb{R}^{m}\right)$. In spite of the standardization conditions the representation (2.8.2.1) is not unique. However for any pair of representations $u_{1}-u_{2}$ and $u_{1}^{\prime}-u_{2}^{\prime}$,

$$
\begin{equation*}
u_{j}(x)-u_{j}^{\prime}(x)=H_{j}(x), j=1,2 \tag{2.8.2.2}
\end{equation*}
$$

where $H_{j}$ are harmonic and $H_{2}(0)=0$.

Really, from the equality $u_{1}-u_{2}=u_{1}^{\prime}-u_{2}^{\prime}$ we obtain $\mu_{1}-\mu_{2}=\mu_{1}^{\prime}-\mu_{2}^{\prime}$. Using Theorem 2.2.1.2 (Jordan decomposition) we obtain $\mu_{1}=\mu_{1}^{\prime}, \mu_{2}=\mu_{2}^{\prime}$. Thus (2.8.2.2) holds. Obviously $H_{2}(0)=0$.

Set

$$
\begin{equation*}
T(r, u):=\frac{1}{\sigma_{m}} \int_{|y|=1} \max \left(u_{1}, u_{2}\right)(r y) d y \tag{2.8.2.3}
\end{equation*}
$$

where $\sigma_{m}$ is the surface square of the unit sphere. It is called the Nevanlinna characteristic of $u \in \delta S H\left(\mathbb{R}^{m}\right)$.

The Nevanlinna characteristic does not depend on the representation (2.8.2.1). Indeed,

$$
\begin{aligned}
\int_{|y|=1} \max \left(u_{1}, u_{2}\right)(r y) d y & =\int_{|y|=1}\left[\left(u_{1}-u_{2}\right)^{+}(r y)-u_{2}(r y)\right] d y \\
& =\int_{|y|=1}\left[\left(u_{1}^{\prime}-u_{2}^{\prime}\right)^{+}(r y)-u_{2}^{\prime}(r y)+H_{2}(r x)\right] d y \\
& =\int_{|y|=1}\left[\max \left(u_{1}^{\prime}, u_{2}^{\prime}\right)(r y)+H_{2}(r x)\right] d y \\
& =\int_{|y|=1} \max \left(u_{1}^{\prime}, u_{2}^{\prime}\right)(r y) d y+H_{2}(0) \\
& =\int_{|y|=1} \max \left(u_{1}^{\prime}, u_{2}^{\prime}\right)(r y) d y
\end{aligned}
$$

Note also that the class $\delta S H\left(\mathbb{R}^{m}\right)$ is linear.
Actually, let $u \in \delta S H\left(\mathbb{R}^{m}\right)$. Then $\lambda u \in \delta S H\left(\mathbb{R}^{m}\right)$ for $\lambda>0$. The function $-u \in \delta S H\left(\mathbb{R}^{m}\right)$, since

$$
-u(x)=\left[u_{2}(x)-u_{1}(0)\right]-\left[u_{1}(x)-u_{1}(0)\right] .
$$

Let us show that $u_{1}+u_{2} \in \delta S H\left(\mathbb{R}^{m}\right)$ if $u, v \in \delta S H\left(\mathbb{R}^{m}\right)$.
Set $\nu:=\nu_{u}+\nu_{v}$, where $\nu_{u}, \nu_{v}$ are the corresponding charges. By Theorem 2.2.1.2 (Jordan decomposition) $\nu=\nu^{+}-\nu^{-}$, where $\nu^{+}, \nu^{-}$are measures concentrated on disjoint sets.

Let $u_{1}$ be a subharmonic function in $\mathbb{R}^{m}$ the mass distribution of which coincides with $\nu^{+} .{ }^{2}$ Then $u_{2}:=u_{1}-(u+v)$ is a subharmonic function with the mass distribution $\nu^{-}$. Hence $u(x)+v(x)=\left[u_{1}(x)-u_{2}(0)\right]-\left[u_{2}(x)-u_{2}(0)\right]$.

Theorem 2.8.2.1 (Properties $T(r, u)$ ) The following holds:
t1) $T(r, u)$ increases monotonically and is convex with respect to $-r^{m-2}$ for $m=$ 2 and with respect to $\log r$ for $m=2$.

[^1]t2) For $u \in S H\left(\mathbb{R}^{m}\right)$, (i.e., $\left.u_{2} \equiv 0\right)$
$$
T(r, u)=\frac{1}{\sigma_{m}} \int_{|y|=1} u^{+}(r y) d y
$$
t3) $T(r, u)=T(r,-u)-u_{1}(0)$.
t4) $T\left(r, u+u^{\prime}\right) \leq T(r, u)+T\left(r, u^{\prime}\right), T(r, \lambda u)=\lambda T(r, u)$ for $\lambda>0$.

Proof. Since $v(x):=\max \left(u_{1}, u_{2}\right)(x)$ is subharmonic, t 1$)$ follows from Theorem 2.6.5.2 (Convexity of $M(r, u)$ and $\mathcal{M}(r, u))$.

The property t2) is obvious, t3) follows from the equality $-u(x)=u_{2}(x)-$ $\left[u_{1}(x)-u_{1}(0)\right]-u_{1}(0)$.

The properties $t 4$ ) follow from the properties of maximum and $t 3$ ).
Set $\rho_{T}[u]:=\rho[a]($ see, (2.8.1.1)) where $a(r):=T(r, u)$. It is called the order of $u(x)$ with respect to $T(r)$.

Theorem 2.8.2.2 ( $\rho_{T}$-property) For $u_{1}, u_{2} \in \delta S H\left(\mathbb{R}^{m}\right)$ the following inequality holds:

$$
\begin{equation*}
\rho_{T}\left[u_{1}+u_{2}\right] \leq \max \left(\rho_{T}\left[u_{1}\right], \rho_{T}\left[u_{2}\right]\right) \tag{2.8.2.4}
\end{equation*}
$$

Equality in (2.8.2.4) is attained if $\rho_{T}\left[u_{1}\right] \neq \rho_{T}\left[u_{2}\right]$.

Proof. Set $u:=u_{1}+u_{2}$. From t3) and t4)

$$
T(r, u) \leq T\left(r, u_{1}\right)+T\left(r, u_{2}\right)+O(1) \leq 2 \max \left[T\left(r, u_{1}\right), T\left(r, u_{2}\right)\right]+O(1)
$$

From the definition of $\rho_{T}$ we obtain (2.8.2.4).
Suppose, for example, $\rho_{T}\left[u_{1}\right]>\rho_{T}\left[u_{2}\right]$. Let us show that $\rho_{T}[u]=\rho_{T}\left[u_{1}\right]$. From the equality $u_{1}=u+\left(-u_{2}\right)$ we obtain $\rho_{T}\left[u_{1}\right] \leq \max \left(\rho_{T}[u], \rho_{T}\left[u_{2}\right]\right.$ If $\rho_{T}[u]<$ $\rho_{T}\left[u_{1}\right]$, then from the previous inequality we would have the contradiction $\rho_{T}\left[u_{1}\right]<$ $\rho_{T}\left[u_{1}\right]$.

Let us define $\sigma_{T}[u]$ by (2.8.1.2) while $\rho:=\rho_{T}[u]$. Set also $\sigma_{T}[u, \rho(r)]:=$ $\sigma[a, \rho(r)]($ see $(2.8 .1 .6))$, where $a(r):=T(r, u)$.

The characteristics $\rho_{T}[u], \sigma_{T}[u], \sigma_{T}[u, \rho(r)]$ are defined for $u \in \delta S H\left(\mathbb{R}^{m}\right)$. For the class of subharmonic functions we have the inclusion $S H\left(\mathbb{R}^{m}\right) \subset \delta S H\left(\mathbb{R}^{m}\right)$ and, of course, all these characteristics can be applied to a subharmonic function. However, for the class $S H\left(\mathbb{R}^{m}\right)$ the standard characteristic of growth is $M(r, u)$ that we can not apply to a $\delta$-subharmonic function $u \in \delta S H\left(\mathbb{R}^{m}\right)$. Thus for $u \in S H\left(\mathbb{R}^{m}\right)$ we define new characteristics $\rho_{M}[u], \sigma_{M}[u], \sigma_{M}[u, \rho(r)]$ in the same way by replacing $T(r, u)$ for $M(r, u)$. The following theorem shows that there is not a big difference between characteristics with respect to $T$ and $M$ for $u \in S H\left(\mathbb{R}^{m}\right)$.

Theorem 2.8.2.3 ( $T$ and $M$-characteristics) Let $u \in S H\left(\mathbb{R}^{m}\right)$ and $\rho(r)(\rightarrow \rho)$ any p.o. Then
$\rho \mathrm{MT1}) \quad \rho_{T}[u]$ and $\rho_{M}[u]$ are finite simultaneously and $\rho_{T}[u]=\rho_{M}[u]:=\rho[u]$
$\rho \mathrm{MT} 2)$ there exists $A:=A(\rho, m)$ such that

$$
A \sigma_{M}[u, \rho(r)] \leq \sigma_{T}[u, \rho(r)] \leq \sigma_{M}[u, \rho(r)]
$$

In particular, the last property means that the types with respect to $T(r)$ and $M(r)$ for the same p.o. are minimal, normal or maximal at the same time.

Proof. From t2), Theorem 2.8.2.1 we have $T(r, u) \leq M(r, u)$ for $u \in S H\left(\mathbb{R}^{m}\right)$. Thus $\rho_{T}[u] \leq \rho_{M}[u]$, proving the second part of $\rho \mathrm{MT} 2$ ).

Let $H(x)$ be the least harmonic majorant of $u(x)$ in the ball $K_{2 R}$. By the Poisson formula (Theorem 2.4.1.5) and Theorem 2.6.1.3,

$$
\begin{align*}
M(R, u) \leq M(R, H) & =\max _{|x|=R} \frac{1}{\sigma_{m} 2 R} \int_{|y|=2 R} u(y) \frac{\left(4 R^{2}-|x|^{2}\right)}{|x-y|^{m}} d s_{y} \\
& \leq \frac{2^{m-2}}{\sigma_{m}} \int_{|y|=1}|u(2 R y)| d s_{y}  \tag{2.8.2.5}\\
& =2^{m-2}[T(2 R, u)+T(2 R,-u)]=2^{m-2}[2 T(2 R, u)-u(0)] .
\end{align*}
$$

From here one can obtain $\rho_{T}[u] \geq \rho_{M}[u]$. The left side of $\rho$ MT2) with $A(\rho, m):=$ $2^{-\rho-m+2}$ follows from the properties of p.o.

Exercise 2.8.2.1 Prove the first inequality from $\rho$ MT2).
2.8.3 Let $\mu$ be a mass distribution (measure) in $\mathbb{R}^{m}\left(\mu \in \mathcal{M}\left(\mathbb{R}^{m}\right)\right)$. The characteristic

$$
\rho[\mu]:=\rho[a]-m+2
$$

for $a(r):=\mu\left(K_{r}\right)$ (see (2.8.1.1)) is called the convergence exponent of $\mu$, and

$$
\bar{\Delta}[\mu]:=\sigma[a]
$$

for the same $a$ (see (2.8.1.2)) is called the upper density of $\mu$.
The least integer number $p$ for which the integral

$$
\begin{equation*}
\int^{\infty} \frac{\mu(t)}{t^{p+m}} d t \tag{2.8.3.1}
\end{equation*}
$$

converges is called the genus of $\mu$ and is denoted $p[\mu]$.
Theorem 2.8.3.1 (Convergence Exponent and Genus) The following holds:
ceg1) $p[\mu] \leq \rho[\mu] \leq p[\mu]+1$,
ceg2) for $\rho[\mu]=p[\mu]+1, \bar{\Delta}[\mu]=0$.

Proof. From Theorem 2.8.1.1 (Convergence Exponent) we have $\rho[\mu]+1+m-2 \leq$ $p[\mu]+m$. Thus $\rho[\mu] \leq p[\mu]+1$. The same theorem implies $\rho[\mu]+m-2+1 \geq$ $p[\mu]+m-1$. Thus $p[\mu] \leq \rho[\mu]$, and ceg1) is proved.

Let $\rho(\mu)=p[\mu]+1$. Then the integral (2.8.3.1) converges for $p[\mu]=\rho[\mu]-1$. We use the inequality

$$
\int_{r}^{\infty} \frac{\mu(t)}{t^{\rho[\mu]+m-1}} d t \geq \mu(r) \int_{r}^{\infty} \frac{d t}{t^{\rho[\mu]+m-1}} d t=\frac{\mu(r)}{r^{\rho[\mu]+m-2}}(\rho[\mu]+m-2)^{-1}
$$

Since the left side of the inequality tends to zero we obtain

$$
\bar{\Delta}[\mu]=\lim _{r \rightarrow \infty} \frac{\mu(r)}{r^{\rho[\mu]+m-2}}=0
$$

Set

$$
\begin{equation*}
\bar{\Delta}[\mu, \rho(r)]:=\sigma[a, \rho(r)+m-2] \tag{2.8.3.2}
\end{equation*}
$$

where $a(r):=\mu(r)$ (see (2.8.1.6)). It is clear that $\rho(r)+m-2$ is also a p.o. Set as in (2.6.5.1),

$$
N(r, \mu):=A(m) \int_{0}^{r} \frac{\mu(t)}{t^{m-1}} d t
$$

where $A(m)=\max (1, m-2)$. Set also

$$
\rho_{N}[\mu]:=\rho[a], \bar{\Delta}_{N}[\mu, \rho(r)]:=\sigma[a, \rho(r)],
$$

where $a(r):=N(r, \mu)$. This is the $N$-order of $\mu$ and the $N$-type of $\mu$ with respect to p.o. $\rho(r)$.

Theorem 2.8.3.2 ( $N$-order and Convergence Exponent) The following holds:
Nce1) $\rho_{N}[\mu]$ and $\rho[\mu]$ are finite simultaneously and $\rho_{N}[\mu]=\rho[\mu]$,
Nce2) for $\rho>0$ there exists such $A_{j}:=A_{j}(\rho, m), j=1,2$, that

$$
A_{1} \bar{\Delta}[\mu, \rho(r)] \leq \bar{\Delta}_{N}[\mu, \rho(r)] \leq A_{2} \bar{\Delta}[\mu, \rho(r)]
$$

Proof. We have the inequality

$$
N(2 r, \mu) \geq A(m) \int_{r}^{2 r} \frac{\mu(t)}{t^{m-1}} d t \geq A(m) \mu(r) \int_{r}^{2 r} \frac{d t}{t^{m-1}} \geq A(m) B(m) \frac{\mu(r)}{(2 r)^{m-2}}
$$

where $B(m):=1-2^{2-m}$ for $m \geq 3$ and $B(2):=\log 2$.
From here one can obtain the inequality $\rho[\mu] \geq \rho_{N}[\mu]$ and the left side of Nce2) for $A_{1}(\rho, m):=A(m) B(m) 2^{-\rho}$. For proving the opposite inequalities we use the l'Hôspital Rule (slightly improved):

$$
\limsup _{r \rightarrow \infty} \frac{N(r, \mu)}{r^{\rho(r)}} \leq \limsup _{r \rightarrow \infty} \frac{N^{\prime}(r, \mu)}{\left(r^{\rho(r))^{\prime}}\right.}=\limsup _{r \rightarrow \infty} \frac{\mu(r) r^{2-m}}{r^{\rho(r)}\left[\rho(r)+r \log r \rho^{\prime}(r)\right]}=\frac{1}{\rho} \bar{\Delta}[\mu] .
$$

Thus $\rho_{N}[\mu] \leq \rho[\mu]$ and the right side of Nce2) holds.

We shall denote as $\delta \mathcal{M}\left(\mathbb{R}^{m}\right)$ the set of charges (signed measures) of the form $\nu:=\mu_{1}-\mu_{2}$ where $\mu_{1}, \mu_{2} \in \mathcal{M}\left(\mathbb{R}^{m}\right)$. Let us remember that $|\nu| \in \mathcal{M}\left(\mathbb{R}^{m}\right)$ is the full variation of $\nu$ (see 2.2.1).
Theorem 2.8.3.3 (Jensen) Let $u:=u_{1}-u_{2} \in \delta S H\left(\mathbb{R}^{m}\right)$ and $\nu:=\mu_{1}-\mu_{2}$ be a corresponding charge. Then
J1) $\rho[|\nu|] \leq \max \left(\rho\left[\mu_{1}\right], \rho\left[\mu_{2}\right]\right) \leq \rho[u]$,
J2) $\bar{\Delta}[|\nu|, \rho(r)] \leq \bar{\Delta}\left[\mu_{1}, \rho(r)\right]+\bar{\Delta}\left[\mu_{2}, \rho(r)\right] \leq A \sigma_{T}[u, \rho(r)]$ for some $A:=A(\rho, m)$.
Proof. We can suppose without loss of generality that $u(0)=0$ because the function $u(x)-u(0)$ has the same order and the same number type if $\rho>0$. We apply the Jensen-Privalov formula (Theorem 2.6.5.1) to the functions $u_{1}, u_{2}$ and obtain

$$
N\left(r, \mu_{j}\right) \leq \mathcal{M}\left(r, u_{j}\right) \leq T(r, u)
$$

Thus $N(r,|\nu|) \leq N\left(r, \mu_{1}\right)+N\left(r, \mu_{2}\right) \leq 2 T(r, u)$. From here one can obtain J1) and J2) for $\rho_{N}[|\nu|]$ and $\bar{\Delta}_{N}[|\nu|, \rho(r)]$. However, we can delete the subscript $N$ because of Theorem 2.8.3.2.

### 2.9 The representation theorem of subharmonic functions in $\mathbb{R}^{m}$

2.9.1 Set

$$
H(z, \cos \gamma, m):= \begin{cases}-\frac{1}{2} \log \left(z^{2}-2 z \cos \gamma+1\right), & \text { for } m=2  \tag{2.9.1.1}\\ \left(z^{2}-2 z \cos \gamma+1\right)^{-\frac{m-2}{2}}, & \text { for } m \geq 3\end{cases}
$$

The function $H(z, \cos \gamma, m)$ is holomorphic on $z$ in the disk $\{|z|<1\}$. It can be represented there in the form

$$
\begin{equation*}
H(z, \cos \gamma, m)=\sum_{k=0}^{\infty} C_{k}^{\frac{m-2}{2}}(\cos \gamma) z^{k} \tag{2.9.1.2}
\end{equation*}
$$

where every coefficient $C_{k}^{\beta}(\bullet), k=0,1, \ldots$ is a polynomial of degree $k$.
Such polynomials are called the Gegenbauer polynomials. Note that $C_{k}^{\frac{1}{2}}(\bullet)$ are the Legendre polynomials and

$$
C_{k}^{0}(\lambda)= \begin{cases}0, & \text { for } k=0 \\ \frac{1}{k} \cos (k \arccos \lambda), & \text { for } k \geq 1,\end{cases}
$$

i.e., they are proportional to the Chebyshev polynomials.

Thus for $m=2$ we have the equality

$$
-\frac{1}{2} \log \left(z^{2}-2 z \cos \gamma+1\right)=\sum_{k=1}^{\infty} \frac{\cos k \gamma}{k} z^{k}
$$

that can be checked directly.
Let $x \in \mathbb{R}^{m}$. Set $x^{0}:=x /|x|$. Then the scalar product $\left(x^{0}, y^{0}\right)$ is equal to $\cos \gamma$ where $\gamma$ is the angle between $x$ and $y$.

Let $\mathcal{E}_{m}(x)$ be defined by (2.4.1.1). For $m \geq 3$ the function $\mathcal{E}_{m}(x-y)$ is the Green function for $\mathbb{R}^{m}$. One can see that it is represented in the form

$$
G\left(x, y, \mathbb{R}^{m}\right):=\mathcal{E}_{m}(x-y)=-|y|^{2-m} H(|x| /|y|, \cos \gamma, m)
$$

where $\cos \gamma=\left(x^{0}, y^{0}\right)$.
For $m=2$ the function $-H(|x| /|y|, \cos \gamma, 2)$ plays the same role. Thus we will denote it as $G\left(x, y, \mathbb{R}^{2}\right)$.
Theorem 2.9.1.1 (Expansion of $G\left(x, y, \mathbb{R}^{m}\right)$ ) The following holds:

$$
\begin{equation*}
G\left(x, y, \mathbb{R}^{m}\right)=-\sum_{k=0}^{\infty} C_{k}^{\frac{m-2}{2}}(\cos \gamma) \frac{|x|^{k}}{|y|^{k+m-2}} \tag{2.9.1.3}
\end{equation*}
$$

for $|x|<|y|$, and the functions

$$
\begin{equation*}
D_{k}(x, y):=C_{k}^{\frac{m-2}{2}}(\cos \gamma) \frac{|x|^{k}}{|y|^{k+m-2}} \tag{2.9.1.4}
\end{equation*}
$$

are homogeneous harmonic functions in $x$ and harmonic in $y$ for $y \neq 0$.
Proof. The expansion (2.9.1.3) follows from (2.9.1.2). The function $G\left(z x, y, \mathbb{R}^{m}\right)$ is harmonic for $|x|<|y|$ and, hence, for any real $0 \leq z<1$. Hence, for any $\psi \in \mathcal{D}\left(K_{r}\right)$ while $r:=0.5|y|$ the function $g(z):=\left\langle G\left(z \bullet, y, \mathbb{R}^{m}\right), \Delta \psi\right\rangle=0$ for $z \in(0,1)$. The function $g$ is holomorphic for all complex $z \in\{|z|<1\}$ because $G\left(z x, y, \mathbb{R}^{m}\right)$ is holomorphic. Thus $g(z) \equiv 0$, i.e., all its coefficients are zero.

From the expansion (2.9.1.3) we can see that the coefficients of $G\left(z x, y, \mathbb{R}^{m}\right)$ are $D_{k}(x, y)$. Hence, $\left\langle D_{k}(\bullet, y), \Delta \psi\right\rangle=0$ for every $\psi \in \mathcal{D}\left(K_{r}\right)$. Thus $D_{k}(\bullet, y)$ is a harmonic distribution. By Theorem 2.4.1.1 it is an ordinary harmonic function for $|x|<0.5|y|$.
$C_{k}^{\frac{m-2}{2}}(\cos \gamma)$ is a polynomial of degree $k$ with respect to $\left(x^{0}, y^{0}\right)$. Thus $D_{k}(x, y)$ is a homogeneous polynomial of $x$ and is harmonic for all $x$.

Let us prove the harmonicity in $y$.
By Theorem 2.4.1.10 the function $D_{k}\left(y^{*}, x^{0}\right)|y|^{2-m}$ ( $*$ stands for inversion) is harmonic in $y$. We have

$$
D_{k}\left(y^{*}, x^{0}\right)|y|^{2-m}=|y|^{2-m} D_{k}\left(y /|y|^{2}, x^{0}\right)=D_{k}\left(x^{0}, y\right)
$$

Set

$$
\begin{equation*}
H(z, \cos \gamma, m, p)=H(z, \cos \gamma, m)-\sum_{k=0}^{p} C_{k}^{\frac{m-2}{2}}(\cos \gamma) z^{k} \tag{2.9.1.5}
\end{equation*}
$$

Theorem 2.9.1.2 The following holds:

$$
\begin{equation*}
|H(z, \cos \gamma, m, p)| \leq A_{1}(m, p)|z|^{p+1} \tag{2.9.1.6}
\end{equation*}
$$

for $|z| \leq 1 / 2$, and

$$
\begin{equation*}
|H(z, \cos \gamma, m, p)| \leq A_{2}(m, p)|z|^{p} \tag{2.9.1.7}
\end{equation*}
$$

for $|z| \geq 2,-\pi<\arg z \leq \pi$.
The factor $|z|^{p}$ should be replaced by $\log |z|$ if $m=2, p=0$.
Proof. Consider the function $\phi(z):=H(z, \cos \gamma, m, p) z^{-p-1}$. It is holomorphic in the disk $\{|z| \leq 1 / 2\}$. We apply the maximum principle and obtain (2.9.1.6) where

$$
A_{1}(m, p)=2^{p+1} \max _{|z|=1 / 2}|\phi(z)|
$$

For proving (2.9.1.7) we consider the function $\psi(z):=H(z, \cos \gamma, m, p) z^{-p}$ that is holomorphic in the domain $D:=\{z:|z| \geq 2,-\pi<\arg z \leq \pi\}$ and continuous in its closure. Applying the maximum principle we obtain (2.9.1.7) where

$$
A_{2}(m, p)=2^{p} \max _{z \in \partial D}|\psi(z)|
$$

Set

$$
G_{p}(x, y, m):=-|y|^{2-m} H(|x| /|y|, \cos \gamma, m, p)
$$

where $\cos \gamma=\left(x^{0}, y^{0}\right)$.
Note the equality

$$
G_{p}(x, y, m)=G\left(x, y, \mathbb{R}^{m}\right)+\sum_{k=0}^{p} D_{k}(x, y)
$$

Exercise 2.9.1.1 Check this using (2.9.1.3), (2.9.1.4) and (2.9.1.5).
It looks like a Green function for $\mathbb{R}^{m}$ but it tends more quickly to zero at infinity and generally speaking it is not negative.

For $m=2$ it can be represented in the form

$$
G_{p}(z, \zeta, 2)=\log |E(z / \zeta, p)|
$$

where $E(z / \zeta, p)$ is the primary Weierstrass factor:

$$
E(z / \zeta, p):=\left(1-\frac{z}{\zeta}\right) \exp \left[\left(\frac{z}{\zeta}\right)+\frac{1}{2}\left(\frac{z}{\zeta}\right)^{2}+\cdots+\frac{1}{p}\left(\frac{z}{\zeta}\right)^{p}\right]
$$

We will call it the primary kernel analogously to the primary factor.

Theorem 2.9.1.3 (Estimate of Primary Kernel) The following holds:

$$
\begin{equation*}
\left|G_{p}(x, y, m)\right| \leq A(m, p) \frac{|x|^{p+1}}{|y|^{p+m-1}} \tag{2.9.1.8}
\end{equation*}
$$

for $|x|<2|y|$,

$$
\begin{equation*}
\left|G_{p}(x, y, m)\right| \leq A(m, p) \frac{|x|^{p}}{|y|^{p+m-2}} \tag{2.9.1.9}
\end{equation*}
$$

for $|y|<2|x|$, and

$$
\begin{equation*}
G_{p}(x, y, m) \leq A(m, p) \min \left(\frac{|x|^{p+1}}{|y|^{p+m-1}}, \frac{|x|^{p}}{|y|^{p+m-2}}\right) \tag{2.9.1.10}
\end{equation*}
$$

for all $x, y \in \mathbb{R}^{m}$, where $A(m, p)$ does not depend on $x, y$.
For $m=2, p=0$ we have $G_{p}(z, \zeta, 2) \leq A(2,0) \log \left(1+\frac{|z|}{|\zeta|}\right)$.

Proof. The inequality (2.9.1.8) follows directly from (2.9.1.6) and (2.9.1.9) follows from (2.9.1.7). By the condition $2 \leq|x| /|y|$ (2.9.1.10) follows from (2.9.1.9).

Suppose $1 / 2 \leq|x| /|y| \leq 2$. Since all the summands in (2.9.1.5) are bounded from below, for $1 / 2 \leq z \leq 2$ we have

$$
G_{p}(x, y, m) \leq A_{1}(m, p)|y|^{2-m} \leq A(m, p) \min \left(\frac{|x|^{p+1}}{|y|^{p+m-1}}, \frac{|x|^{p}}{|y|^{p+m-2}}\right)
$$

also under these conditions.
The case $m=2, p=0$ is obvious.
2.9.2 Let $\mu \in \mathcal{M}\left(\mathbb{R}^{m}\right)$. We suppose below that its support does not contain the origin.

We will say that the integral $\int_{\mathbb{R}^{m}} f(x, y) \mu(d y)$ converges uniformly on $x \in D$ if

$$
\sup _{x \in D}\left|\int_{|y|>R} f(x, y) \mu(d y)\right| \rightarrow 0
$$

when $R \rightarrow \infty$.
Hence, the integral is permitted to be equal to infinity for some finite $x$.
Let $\mu$ have genus $p$ (see, 2.8.3). Set

$$
\begin{equation*}
\Pi(x, \mu, p):=\int_{\mathbb{R}^{m}} G_{p}(x, y, m) \mu(d y) \tag{2.9.2.1}
\end{equation*}
$$

It is called the canonical potential.

In particular, let $m=2$ and $\mu:=n$ be a zero distribution, i.e., it has unit masses concentrated on a discrete point set $\left\{z_{j}: j=1,2, \ldots\right\}$. Then

$$
\Pi(z, n, p)=\log \left|\prod_{j=1}^{\infty} E\left(\frac{z}{z_{j}}, p\right)\right|
$$

where

$$
\prod_{j=1}^{\infty} E\left(\frac{z}{z_{j}}, p\right)
$$

is the canonical Weierstrass product.
Theorem 2.9.2.1 (Brelot-Weierstrass) The canonical potential (2.9.2.1) converges uniformly on any bounded domain. It is a subharmonic function with $\mu$ as its Riesz measure.

Proof. Let $|x|<R_{0}$ and $|y|>R$. From the estimate of the primary kernel (Theorem 2.9.1.3) we have

$$
\begin{aligned}
\left|\int_{|y|>R} G_{p}(x, y, m) \mu(d y)\right| & \leq A(m, p)|x|^{p+1} \int_{|y|>R}|y|^{-p-m+1} \mu(d y) \\
& =A(m, p)|x|^{p+1} \int_{R}^{\infty} t^{-p-m+1} \mu(d t)
\end{aligned}
$$

Integrating by part we obtain

$$
\int_{R}^{\infty} t^{-p-m+1} \mu(d t)=\frac{\mu(R)}{R^{p+m-1}}+(p+m-1) \int_{R}^{\infty} \frac{\mu(t)}{t^{p+m}} d t
$$

The last integral converges since the genus of $\mu$ is $p$. Hence, both summands tend to zero when $R \rightarrow \infty$. Thus

$$
\sup _{|x|<R_{0}}\left|\int_{|y|>R} G_{p}(x, y, m) \mu(d y)\right| \rightarrow 0
$$

while $R_{0}$ is fixed and $R \rightarrow \infty$, i.e., the canonical potential converges uniformly on any bounded domain.

Let us represent the canonical potential for $R>R_{0}$ in the form

$$
\begin{aligned}
\Pi(x, \mu, p)=\int_{|y|<R} G\left(x, y, \mathbb{R}^{m}\right) \mu(d y) & +\int_{|y|<R} \sum_{k=0}^{p} D_{k}(x, y) \mu(d y) \\
& +\int_{|y|>R} G_{p}(x, y, m) \mu(d y)
\end{aligned}
$$

The first summand is a potential, hence a subharmonic function and its Riesz measure coincide with $\mu$. The other summands are harmonic for $|x|<R_{0}$.

The following proposition estimates the growth of the canonical potential in terms of its masses.

Theorem 2.9.2.2 (Estimation of Canonical Potential) The following inequality holds:

$$
\begin{equation*}
M(r, \Pi(\bullet, \mu, p)) \leq A\left[\int_{0}^{\infty} \frac{\mu(r \tau)}{r^{m-2}} \frac{\min \left(1, \tau^{-1}\right)}{\tau^{p+m-1}} d \tau+\frac{\mu(r)}{r^{m-1}}\right] \tag{2.9.2.2}
\end{equation*}
$$

where $A:=A(m, p)$ does not depend on $r$ and $\mu$.
Proof. From (2.9.1.10),

$$
\Pi(x, \mu, p) \leq A(m, p) \int_{\mathbb{R}^{m}} \min \left(\frac{|x|^{p+1}}{|y|^{p+m-1}}, \frac{|x|^{p}}{|y|^{p+m-2}}\right) \mu(d y)
$$

Set $r:=|x|, t:=|y|$. Then we have

$$
\begin{equation*}
M(r, \Pi(\bullet, \mu, p)) \leq A \int_{0}^{\infty} \min \left(\frac{r^{p+1}}{t^{p+m-1}}, \frac{r^{p}}{t^{p+m-2}}\right) \mu(d t) \tag{2.9.2.3}
\end{equation*}
$$

The integral on the right side of (2.9.2.3) can be represented in the form

$$
\int_{0}^{r} \frac{r^{p}}{t^{p+m-2}} \mu(d t)+\int_{r}^{\infty} \frac{r^{p+1}}{t^{p+m-1}} \mu(d t)
$$

Integrating every integral by parts we obtain

$$
\begin{aligned}
& (p+m-2) \int_{0}^{r} \frac{r^{p}}{t^{p+m-1}} \mu(t) d t+(p+m-1) \int_{r}^{\infty} \frac{r^{p+1}}{t^{p+m}} \mu(t) d t+\frac{\mu(r)}{r^{m-1}} \\
& \quad \leq(p+m-1) \int_{0}^{\infty} \min \left(1, \frac{r}{t}\right) \frac{r^{p}}{t^{p+m-1}} \mu(t) d t+\frac{\mu(r)}{r^{m-1}}
\end{aligned}
$$

After the change $t=r \tau$ we obtain (2.9.2.2) where the new $A(m, p)$ is equal to $A(m, p)(p+m-1)$.

Theorem 2.9.2.3 (Brelot-Borel) The order of the canonical potential is equal to the convergence exponent of its mass distribution, i.e.,

$$
\rho[\Pi(\bullet, \mu, p)]=\rho[\mu],
$$

if the genus of $\mu$ is equal to $p$.
Proof. First assume $\rho[\mu]<p+1$. Let us choose $\lambda$ such that $\rho[\mu]<\lambda<p+1$.
For some constant $C$ that does not depend on $t$ we have $\mu(t) \leq C t^{\lambda+m-2}$.
Actually, $\mu(t) / t^{\lambda+m-2} \rightarrow 0$, because $\lambda>\rho[\mu]$. Since $\mu(t)=0$ for small $t$, this function is bounded and we can take its lower bound as $C$.

Now we have

$$
\begin{equation*}
f(r, \tau):=\frac{\mu(r \tau)}{r^{\lambda+m-2}} \frac{\min \left(1, \tau^{-1}\right)}{\tau^{p+m-1}} \leq C \tau^{\lambda-p-1} \min (1,1 / \tau) \tag{2.9.2.4}
\end{equation*}
$$

for all $\tau \in(0, \infty)$.
We also have

$$
\begin{equation*}
\lim _{r \rightarrow \infty} f(r, \tau)=0 \tag{2.9.2.5}
\end{equation*}
$$

because of $\lambda>\rho[\mu]$.
Let us divide (2.9.2.2) by $r^{\lambda}$ and pass to the upper limit. By Fatou's lemma (Theorem 2.2.2.3)

$$
\begin{equation*}
\limsup _{r \rightarrow \infty} \frac{M(r, \Pi(\bullet, \mu, p))}{r^{\lambda}} \leq A(m, p)\left[\int_{0}^{\infty} \limsup _{r \rightarrow \infty} f(r, \tau) d \tau+\limsup _{r \rightarrow \infty} \frac{\mu(r)}{r^{\lambda+m-1}}\right]=0 . \tag{2.9.2.6}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
\lambda \geq \rho[\Pi(\bullet, \mu, p)] . \tag{2.9.2.7}
\end{equation*}
$$

Since this holds for any $\lambda>\rho[\mu]$, we have $\rho[\mu] \geq \rho[\Pi(\bullet, \mu, p)]$ under the assumption $\lambda<p[\mu]+1$.

Let $\rho[\mu]=p[\mu]+1$. By Theorem 2.8.3.1, $\bar{\Delta}[\mu]=0$. Hence, $\mu(t) t^{-p-m+1} \leq C$ and

$$
f(r, \tau):=\frac{\mu(r \tau)}{(r \tau)^{p+m-1}} \min \left(1, \tau^{-1}\right) \leq C \min (1,1 / \tau)
$$

The function $\min (1,1 / \tau)$ is not summable on $(0, \infty)$. Therefore we will act in a slightly different way. From Theorem 2.9.2.2 we have

$$
\begin{aligned}
\limsup _{r \rightarrow \infty} \frac{M(r, \Pi(\bullet, \mu, p))}{r^{p+1}} \leq & A(m, p)\left[\int_{0}^{1} \limsup _{r \rightarrow \infty} f(r, \tau) d \tau\right] \\
& +A(m, p)\left[\limsup _{r \rightarrow \infty} \int_{r}^{\infty} \frac{\mu(t)}{t^{p+m}} d t++\limsup _{r \rightarrow \infty} \frac{\mu(r)}{r^{p+m}}\right] .
\end{aligned}
$$

The first integral is equal to zero because $\bar{\Delta}[\mu]=0$. The second addend vanishes since the integral converges. Thus we have $p+1=\rho[\mu] \geq \rho[\Pi(\bullet, \mu, p)]$.

The reverse inequality holds for any subharmonic function in $\mathbb{R}^{m}$ by the Jensen theorem (Theorem 2.8.3.3).
2.9.3 Let us denote as $\delta S H(\rho)$ the class of functions $u \in \delta S H\left(\mathbb{R}^{m}\right)$ for which $\rho_{T}[u] \leq \rho$.
Theorem 2.9.3.1 (Brelot-Hadamard) Let $u=u_{1}-u_{2} \in \delta S H(\rho)$, and let $p_{1}$, $p_{2}$ be the genuses of the mass distributions $\mu_{j}:=\mu_{u_{j}}, j=1,2$. Suppose $\operatorname{supp}\left[\mu_{1}-\mu_{2}\right] \cap$ $\{0\}=\varnothing$.

Then the following equality holds:

$$
u(x)=\Pi\left(x, \mu_{1}, p_{1}\right)-\Pi\left(x, \mu_{2}, p_{2}\right)+\Phi_{q}(x)
$$

where $\Phi_{q}(x)$ is a harmonic polynomial of degree $q \leq \rho$.

Proof. The function $v(x):=u(x)-\Pi\left(x, \mu_{1}, p_{1}\right)+\Pi\left(x, \mu_{2}, p_{2}\right)$ is harmonic by the Brelot-Weierstrass theorem (Theorem 2.9.2.1). We also have the inequality

$$
\begin{equation*}
\rho_{T}[v] \leq \max \left(\rho_{T}[u], \rho_{T}\left[\Pi\left(\bullet, \mu_{1}, p_{1}\right)\right], \rho_{T}\left[\Pi\left(\bullet, \mu_{2}, p_{2}\right)\right]\right) \tag{2.9.3.1}
\end{equation*}
$$

by Theorem 2.8.2.2 ( $\rho_{T}$-properties). The property $\rho \mathrm{MT} 1$ ) (Theorem 2.8.2.3) implies

$$
\rho_{T}\left[\Pi\left(\bullet, \mu_{j}, p_{j}\right)\right]=\rho_{M}\left[\Pi\left(\bullet, \mu_{j}, p_{j}\right)\right]:=\rho\left[\Pi\left(\bullet, \mu_{j}, p_{j}\right)\right], j=1,2 .
$$

The Brelot-Borel theorem (Theorem 2.9.2.3) implies

$$
\rho\left[\Pi\left(\bullet, \mu_{j}, p_{j}\right)\right]=\rho\left[\mu_{j}\right], \quad j=1,2
$$

The Jensen theorem (Theorem 2.8.3.3) implies

$$
\max \left(\rho\left[\mu_{1}\right], \rho\left[\mu_{2}\right]\right) \leq \rho_{T}[u]
$$

From (2.9.3.1) we have

$$
\rho_{T}[v] \leq \rho_{T}[u] \leq \rho
$$

Since $v$ is subharmonic, $\rho_{T}[v]=\rho_{M}[v]:=\rho[v]$ by Theorem 2.8.2.3, and $\rho[v] \leq \rho$. Therefore

$$
\lim _{r \rightarrow \infty} \frac{M(r, v)}{r^{\rho+\epsilon}}=0
$$

for arbitrarily small $\epsilon>0$.
By the Liouville theorem (Theorem 2.4.2.3) $v(x)$ is a harmonic polynomial of degree $q \leq \rho+\epsilon$, and thus $v(x)=\Phi_{q}(x)$ for $q \leq \rho$.

For a non-integer $\rho$ the Brelot-Hadamard theorem allows us to connect the growth of functions and masses more tightly than in the Jensen theorem.

Theorem 2.9.3.2 (Sharpening of Jensen) Let $\rho>0$ and be non-integer, $u=u_{1}-$ $u_{2} \in \delta S H\left(\mathbb{R}^{m}\right)$ with $\rho_{T}[u]=\rho$, and let $\nu_{u}=\mu_{1}-\mu_{2}$ the corresponding charge. Then
pJ1) $\quad \rho\left[\nu_{u}\right]=\max \left(\rho\left[\mu_{1}\right], \rho\left[\mu_{2}\right]\right)=\rho$,
pJ2) $\quad A_{1} \sigma_{T}[u, \rho(r)] \leq \bar{\Delta}\left[\nu_{u}, \rho(r)\right] \leq \bar{\Delta}\left[\mu_{1}, \rho(r)\right]+\bar{\Delta}\left[\mu_{2}, \rho(r)\right] \leq A_{2} \sigma_{T}[u, \rho(r)]$, where $A_{j}=A_{j}(m, \rho)$ and $\rho(r)$ is an arbitrarily proximate order such that $\rho(r) \rightarrow \rho$ when $r \rightarrow \infty$.

For proving this theorem we need
Theorem 2.9.3.3 Let $\Pi(x, \mu, p)$ be a canonical potential with non-integer $\rho[\mu]:=$ [ $\rho$ ], and let $\rho(r)(\rightarrow \rho)$ be a proximate order. Then

$$
\begin{equation*}
\sigma[\Pi(\bullet, \mu, p), \rho(r)] \leq A(m, \rho, p) \bar{\Delta}[\mu, \rho(r)] \tag{2.9.3.2}
\end{equation*}
$$

Proof. We can suppose without loss of generality that $\bar{\Delta}[\mu, \rho(r)]<\infty$. By this condition and since $\mu(t)=0,0<t<c$ for some $c>0$, we have the inequality

$$
\mu(t) t^{-\rho(t)-m+2} \leq C
$$

for all $t \in(0, \infty)$ and some $C>0$ that does not depend on $t$. Set

$$
I(r):=\int_{c / r}^{\infty} \frac{\mu(r t)}{r^{\rho(r)+m-2}} \frac{\min (1,1 / t)}{t^{p+m-1}} d t
$$

By Theorem 2.9.2.2 we have

$$
\begin{equation*}
\sigma[\Pi(\bullet, \mu, p), \rho(r)]=\limsup _{r \rightarrow \infty} \frac{M(r, \Pi(\bullet, \mu, p)}{r^{\rho(r)}} \leq A(m, p) \limsup _{r \rightarrow \infty} I(r) . \tag{2.9.3.3}
\end{equation*}
$$

Let us choose $r_{\epsilon}$ such that

$$
\sup _{r>r_{\epsilon}} \frac{\mu(r \epsilon)}{(r \epsilon)^{\rho(r \epsilon)+m-2}} \leq \bar{\Delta}[\mu, \rho(r)]+\epsilon .
$$

For such $r$ we have

$$
\begin{aligned}
I(r)= & \int_{c / r}^{\infty} \frac{\mu(r t)}{(r t)^{\rho(r t)+m-2}} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} \frac{\min (1,1 / t)}{t^{p+1}} d t \\
\leq & \sup _{c / r \leq t \leq \epsilon} \frac{\mu(r t)}{(r t)^{\rho(r t)+m-2}} \int_{c / r}^{\epsilon} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} \frac{\min (1,1 / t)}{t^{p+1}} d t \\
& +\sup _{\epsilon \leq t \leq 1 / \epsilon} \cdots \int_{\epsilon}^{1 / \epsilon} \cdots d t+\sup _{1 / \epsilon \leq t \leq \infty} \cdots \int_{1 / \epsilon}^{\infty} \cdots d t \\
\leq & C \int_{c / r}^{\epsilon} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}} \frac{\min (1,1 / t)}{t^{p+1}} d t+(\bar{\Delta}[\mu, \rho(r)]+\epsilon) \int_{\epsilon}^{1 / \epsilon} \cdots d t+C \int_{1 / \epsilon}^{\infty} \cdots d t .
\end{aligned}
$$

The function

$$
f(t):=\frac{\min (1,1 / t)}{t^{p+1}}
$$

satisfies the conditions of Gol'dberg's theorem (Theorem 2.8.1.5) with $p+1-\rho<$ $\delta<1$ and $0<\gamma<p+1-\rho$. Passing to the limit we have

$$
\begin{aligned}
\limsup _{r \rightarrow \infty} I(r) \leq C & \int_{0}^{\epsilon} t^{\rho-p} d t+(\bar{\Delta}[\mu, \rho(r)]+\epsilon) \\
& \times \int_{\epsilon}^{1 / \epsilon} t^{\rho-p-1} \min (1,1 / t) d t+C \int_{1 / \epsilon}^{\infty} t^{\rho-p-2} d t
\end{aligned}
$$

Passing to the limit as $\epsilon \rightarrow 0$ we obtain with the help of (2.9.3.3)

$$
\sigma[\Pi(\bullet, \mu, p), \rho(r)] \leq A(m, p) \bar{\Delta}[\mu, \rho(r)] \int_{0}^{\infty} t^{\rho-p-1} \min (1,1 / t) d t
$$

Proof of Theorem 2.9.3.2. The inequality $\rho\left[\nu_{u}\right] \leq \rho$ and the last inequality in pJ2) follow from the Jensen theorem (Theorem 2.8.3.3). Let us prove the reverse inequality and the left side.

Since $\rho$ is non-integer, $q<\rho$ in the Brelot-Hadamard theorem (Theorem 2.9.3.1). Hence $M\left(r, \Phi_{q}\right)=o\left(r^{\rho}\right)$ and

$$
T(r, u) \leq T\left(r, \Pi\left(\bullet, \mu_{1}, p\right)\right)+T\left(r, \Pi\left(\bullet, \mu_{2}, p\right)\right)+o\left(r^{\rho}\right)
$$

Thus

$$
\begin{aligned}
\rho_{T}[u] & \leq \max \left(\rho\left[\Pi\left(\bullet, \mu_{1}, p\right)\right], \rho\left[\Pi\left(\bullet, \mu_{2}, p\right]\right)\right. \\
\sigma_{T}[u, \rho(r)] & \leq \max \left(\sigma_{T}\left[\Pi\left(\bullet, \mu_{1}, p\right), \rho(r)\right], \sigma_{T}\left[\Pi\left(\bullet, \mu_{2}, p\right], \rho(r)\right]\right) .
\end{aligned}
$$

From Theorem 2.9.3.3 we obtain

$$
\begin{aligned}
\rho_{T}[u] & \leq \max \left(\rho\left[\mu_{1}\right], \rho\left[\mu_{2}\right]\right) \\
\sigma_{T}[u, \rho(r)] & \leq A(m, \rho, p) \max \left(\bar{\Delta}\left[\mu_{1}, \rho(r)\right], \bar{\Delta}\left[\mu_{2}, \rho(r)\right]\right. \\
& =A(m, \rho, p) \bar{\Delta}[|\nu|, \rho(r)]
\end{aligned}
$$

We can set $A_{1}:=A^{-1}(m, \rho, p)$ and obtain the left side of pJ 2$)$.
2.9.4 Let $u \in \delta S H\left(\mathbb{R}^{m}\right)$ and $\rho:=\rho_{T}[u]$ be an integer number. We can always represent the function $u$ in the form

$$
\begin{equation*}
u(x)=\Pi(x, \nu, \rho)+\Phi_{\rho}(x) \tag{2.9.4.1}
\end{equation*}
$$

where $\Phi_{\rho}(x)$ is a harmonic polynomial of degree at most $\rho$. Actually, such a representation can be obtained from Theorem 2.9.3.1 by addition and subtraction of terms of the form

$$
\Phi_{k_{j}}(x):=\int_{\mathbb{R}^{m}} D_{k_{j}}(x, y) \mu_{j}(d y), j=1,2
$$

where $p_{j}<k_{j} \leq \rho$. All $\Phi_{k_{j}}(x)$ of such a kind are harmonic polynomials of degree at most $\rho$. Set

$$
\begin{align*}
\Pi_{<}^{R}(x, \nu, \rho-1) & :=\int_{|y|<R} G_{\rho-1}(x, y, m) \nu(d y)  \tag{2.9.4.2}\\
\Pi_{>}^{R}(x, \nu, \rho) & :=\int_{|y| \geq R} G_{\rho}(x, y, m) \nu(d y)  \tag{2.9.4.3}\\
\delta_{R}(x, \nu, \rho) & :=\int_{|y|<R} D_{\rho}(x, y) \nu(d y) \tag{2.9.4.4}
\end{align*}
$$

In particular, for $m=2$,

$$
\begin{equation*}
\delta_{R}(z, \nu, \rho):=\frac{1}{\rho} \int_{|\zeta|<R} \Re\left(\frac{z}{\zeta}\right)^{\rho} \nu(d \zeta) . \tag{2.9.4.4a}
\end{equation*}
$$

Let $Y_{\rho}(x)$ be the homogeneous polynomial of degree $\rho$ from the polynomial $\Phi_{\rho}$ in (2.9.4.1). Set also

$$
\begin{align*}
\delta_{R}(x, u, \rho) & :=\delta_{R}(x, \nu, \rho)+Y_{\rho}(x), \\
M(r, \delta) & :=\max _{|y|=1}\left|\delta_{r}(r y, u, \rho)\right|,  \tag{2.9.4.5}\\
\bar{\Delta}_{\delta}[u, \rho] & :=\limsup _{r \rightarrow \infty} M(r, \delta) r^{-\rho(r)} .
\end{align*}
$$

The functions $\delta_{R}(x, \nu, \rho)$ are homogeneous polynomials that are determined completely by their values on the unit sphere. Thus, by the Harnack theorem (Theorem 2.4.1.7) we have

Theorem 2.9.4.1 $\bar{\Delta}_{\delta}[u, \rho(r)]<\infty$ if and only if the family $\delta_{R}(x, u, \rho) R^{\rho-\rho(R)}, R>$ 0 is precompact in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$.

Let $\rho$ be an integer number and $\rho(r) \rightarrow \rho$ be a p.o. Set

$$
\Omega[u, \rho(r)]:=\max \left(\bar{\Delta}_{\delta}[u, \rho(r)], \bar{\Delta}\left[\left|\nu_{u}\right|, \rho(r)\right] .\right.
$$

Theorem 2.9.4.2 (Brelot-Lindelöf) The following holds:

$$
A_{1} \Omega[u, \rho(r)] \leq \sigma_{T}[u, \rho(r)] \leq A_{2} \Omega[u, \rho(r)]
$$

where $A_{j}:=A_{j}(m, \rho)$.
For proving this theorem we will first study the function $\Pi_{<}^{R}$ and $\Pi_{>}^{R}$. Set

$$
\begin{aligned}
T(r, \lambda,>) & :=T\left(r, \Pi_{>}^{\lambda r}(\bullet, \nu, \rho)\right) \\
T(r, \lambda,<) & :=T\left(r, \Pi_{<}^{\lambda r}(\bullet, \nu, \rho-1)\right) .
\end{aligned}
$$

Theorem 2.9.4.3 (Estimate of $T(\bullet,>, T(\bullet,<))$ The following holds:

$$
\begin{align*}
T(r, \lambda,>) \leq & A\left(\int_{\lambda}^{\infty} \frac{|\nu|(r t)}{r^{m-2}} \frac{\min \left(1, t^{-1}\right)}{t^{\rho+m-1}} d t+\frac{|\nu|(r)}{r^{m-1}}\right)  \tag{2.9.4.6}\\
T(r, \lambda,<) \leq & A\left(\int_{0}^{\lambda} \frac{|\nu|(r t)}{r^{m-2}} \frac{\min \left(1, t^{-1}\right)}{t^{\rho+m-2}} d t\right) \\
& +A\left(\frac{|\nu|(r \lambda)}{r^{m-2}} \int_{\lambda}^{\infty} \frac{\min \left(1, t^{-1}\right)}{t^{\rho+m-2}} d t+\frac{|\nu|(r)}{r^{m-1}}\right) \tag{2.9.4.7}
\end{align*}
$$

where $A:=A(m, \rho)$.
Proof. Let $\nu=\mu_{1}-\mu_{2}$. Then $|\nu|=\mu_{1}+\mu_{2}$. We have

$$
\begin{equation*}
\Pi_{<}^{R}(x, \nu, \rho-1)=\Pi_{<}^{R}\left(x, \mu_{1}, \rho-1\right)-\Pi_{<}^{R}\left(x, \mu_{2}, \rho-1\right) . \tag{2.9.4.8}
\end{equation*}
$$

Since $\Pi_{<}^{R}\left(0, \mu_{2}, \rho-1\right)=0$, we have (see t3),t4), Theorem 2.8.2.1)

$$
\begin{equation*}
T\left(r, \Pi_{<}^{R}(\bullet, \nu, \rho-1)\right) \leq T\left(r, \Pi_{<}^{R}\left(\bullet, \mu_{1}, \rho-1\right)\right)+T\left(r, \Pi_{<}^{R}\left(\bullet, \mu_{2}, \rho-1\right)\right) \tag{2.9.4.9}
\end{equation*}
$$

Set

$$
\Pi_{1}:=\Pi_{<}^{R}\left(\bullet, \mu_{1}, \rho-1\right), \Pi_{2}:=\Pi_{<}^{R}\left(\bullet, \mu_{2}, \rho-1\right) .
$$

Let us estimate, for example, $T\left(r, \Pi_{1}\right)$. The masses of the canonical potential $\Pi_{1}$ are concentrated in $K_{R}$. Applying Theorem 2.9.2.2 (Estimation of Canonical Potential) for $p=\rho-1$ we obtain

$$
\begin{aligned}
T\left(r, \Pi_{1}\right) & \leq M\left(r, \Pi_{1}\right) \\
& \leq A \int_{0}^{\frac{R}{r}} \frac{\mu_{1}(r t)}{r^{m-2}} \frac{\min \left(1, t^{-1}\right)}{t^{\rho+m-2}} d t+A \frac{\mu_{1}(R)}{r^{m-2}} \int_{\frac{R}{r}}^{\infty} \frac{\min \left(1, t^{-1}\right)}{t^{\rho+m-2}} d t+\frac{\mu_{1}(r)}{r^{m-1}}
\end{aligned}
$$

Set $R:=r \lambda$. Then we obtain the inequality (2.9.4.7) for $\nu:=\mu_{1}$. Analogously one can do the same for $\nu:=\mu_{2}$. The inequality (2.9.4.9) allows us to pass to the limit in (2.9.4.7) in the general case.

Set $\Pi_{1}:=\Pi_{>}^{R}\left(\bullet, \mu_{1}, \rho\right)$. Applying (2.9.2.2) for $p=\rho$ we obtain

$$
T\left(r, \Pi_{1}\right) \leq M\left(r, \Pi_{1}\right) \leq \int_{\frac{R}{r}}^{\infty} \frac{\mu_{1}(r t)}{r^{m-2}} \frac{\min \left(1, t^{-1}\right)}{t^{\rho+m-1}} d t+\frac{\mu_{1}(r)}{r^{m-1}}
$$

In the same way we obtain (2.9.4.6).
Set

$$
\begin{aligned}
& \sigma\left[\Pi_{>}, \rho(r)\right]:=\limsup _{r \rightarrow \infty} \frac{T\left(r, \Pi_{>}^{r}(\bullet, \nu, \rho)\right)}{r^{\rho(r)}}, \\
& \sigma\left[\Pi_{<,}, \rho(r)\right]:=\limsup _{r \rightarrow \infty} \frac{T\left(r, \Pi_{<}^{r}(\bullet, \nu, \rho)\right)}{r^{\rho(r)}} .
\end{aligned}
$$

Theorem 2.9.4.4 Let $\nu:=\mu_{1}-\mu_{2} \in \delta \mathcal{M}(\rho)$ and $\rho$ an integer number. Then for any p.o. $\rho(r) \rightarrow \rho$,

$$
\max \left(\sigma\left[\Pi_{>}, \rho(r)\right], \sigma\left[\Pi_{<}, \rho(r)\right]\right) \leq A \bar{\Delta}[|\nu|, \rho(r)]
$$

where $A:=A(m, \rho)$.

Proof. From (2.9.4.6) we have

$$
T\left(r, \Pi_{>}^{r}(\bullet, \nu, \rho)\right)=T(r, 1,>) \leq A \int_{1}^{\infty} \frac{|\nu|(r t)}{r^{m-2}} \frac{1}{t^{\rho+m}} d t+\frac{|\nu|(r)}{r^{m-1}}
$$

Now we repeat the reasoning of Theorem 2.9.3.3 for $\mu:=|\nu|$ and $p:=\rho$. We will obtain

$$
\sigma\left[\Pi_{>}, \rho(r)\right] \leq A \bar{\Delta}[|\nu|, \rho(r)] \int_{1}^{\infty} t^{-2} d t
$$

For the other case we have from (2.9.4.7),

$$
\begin{aligned}
& T\left(r, \Pi_{<}^{r}(\bullet, \nu, \rho-1)\right)=T(r, 1,<) \\
& \quad \leq A \int_{0}^{1} \frac{|\nu|(r t)}{r^{m-2}} \frac{1}{t^{\rho+m-1}} d t+A \frac{|\nu|(r)}{r^{m-2}}\left(\int_{1}^{\infty} t^{-\rho-m+1} d t+r^{-1}\right) .
\end{aligned}
$$

We divide this inequality by $r^{\rho(r)}$ and pass to the upper limit while $r \rightarrow \infty$.
The first summand of the right side gives

$$
A \bar{\Delta}[|\nu|, \rho(r)] \int_{0}^{1} d t
$$

by the reasoning of Theorem 2.9.3.3.
The second one can be computed directly, yielding

$$
A \bar{\Delta}[|\nu|, \rho(r)] \int_{1}^{\infty} t^{-\rho-m+1} d t
$$

Combining all these inequalities we obtain the assertion of the theorem.
Proof of Theorem 2.9.4.2. Let us represent $u(x)$ in the form

$$
\begin{equation*}
u(r y)=\Pi_{<}^{r}\left(r y, \nu_{u}, \rho-1\right)+\Pi_{>}^{r}\left(r y, \nu_{u}, \rho\right)+\delta_{r}(r y, u, \rho)+o\left(r^{\rho-1}\right) \tag{2.9.4.10}
\end{equation*}
$$

where $|y|=1$.
Then we have

$$
T(r, u) \leq T\left(r, \Pi_{<}^{r}\left(\bullet, \nu_{u}, \rho-1\right)\right)+T\left(r, \Pi_{>}^{r}\left(\bullet, \nu_{u}, \rho\right)\right)+M(r, \delta)+o\left(r^{\rho-1}\right) .
$$

Let us divide this by $r^{\rho(r)}$ and pass to the upper limit. By Theorem 2.9.4.4 we obtain

$$
\sigma_{T}[u, \rho(r)] \leq A \max \left(\bar{\Delta}[|\nu|, \rho(r)], \bar{\Delta}_{\delta}[u, \rho(r)]\right)=A_{2} \Omega[u, \rho(r)]
$$

where $A_{2}=A(m, \rho)$. Let us write (2.9.4.11) in the form

$$
\delta_{r}(r y, u, \rho)=u(r y)-\Pi_{<}^{r}\left(r y, \nu_{u}, \rho-1\right)-\Pi_{>}^{r}\left(r y, \nu_{u}, \rho\right)+o\left(r^{\rho-1}\right)
$$

We obtain

$$
T\left(r, \delta_{r}(\bullet, u, \rho)\right) \leq T(r, u)+T\left(r, \Pi_{<}^{r}\left(\bullet, \nu_{u}, \rho-1\right)\right)+T\left(r, \Pi_{>}^{r}\left(\bullet, \nu_{u}, \rho\right)+o\left(r^{\rho-1}\right) .\right.
$$

Since $\delta_{R}(\bullet, u, \rho)$ is harmonic and homogeneous, we have by (2.8.2.5)

$$
M\left(r, \delta_{R}\right) \leq 2^{m-1} T\left(2 r, \delta_{R}\right)=2^{m-1+\rho} T\left(r, \delta_{R}\right)
$$

Therefore we obtain the inequality

$$
\bar{\Delta}_{\delta}[u, \rho(r)] \leq \sigma_{T}[u, \rho(r)]+2 A \bar{\Delta}[|\nu|, \rho(r)] .
$$

By the Jensen theorem (Theorem 2.8.3.3) we have

$$
\Omega[u, \rho(r)] \leq A_{1}^{-1} \sigma_{T}[u, \rho(r)]
$$

for some $A_{1}=A_{1}(m, \rho)$.

## Chapter 3

## Asymptotic Behavior of Subharmonic Functions of Finite Order

### 3.1 Limit sets

3.1.1 Let $\left\{V_{t}: t \in(0, \infty)\right\}$ be a family of rotations of $\mathbb{R}^{m}$ that form a oneparametric group, i.e.,

$$
\begin{equation*}
V_{t_{1}} V_{t_{2}}=V_{t_{1} t_{2}}, \quad V_{1}=I \tag{3.1.1.0}
\end{equation*}
$$

where $I$ is the identity map.
The family of linear transformations

$$
\begin{equation*}
P_{t}:=t V_{t} \tag{3.1.1.1}
\end{equation*}
$$

is also a one-parametric group.
In particular, for $m=2$ the general form of the rotations is

$$
V_{t} z=z \exp (i \gamma \log t)
$$

where $\gamma$ is real.
The orbit $\left\{P_{t} z: t \in(0, \infty)\right\}$ of every point $z \neq 0$ is a logarithmic spiral if $\gamma \neq 0$ and a ray if $\gamma=0$.

For $m \geq 3$ and $V_{t} \equiv I, t \in(0, \infty)$ the orbit of every point $x \neq 0$ is a ray from the origin. For other $V_{\bullet}$ it is a spiral connecting the origin to infinity.

It is clear that only one orbit $\left\{P_{t} x: t \in(0, \infty)\right\}$ passes through every $x \neq 0$. The behavior of every point $y(t):=P_{t} x$ is completely determined by a system of
differential equations with constant coefficients:

$$
\frac{d}{d t} y=\left(I+V^{\prime}\right) y, \quad V^{\prime}:=\left.\frac{d}{d t} V_{t}\right|_{t=1}
$$

with the initial condition of $y(1)=x$.
3.1.2 Let $u \in S H(\rho)$ and $\sigma_{M}[u, \rho(r)]<\infty$ for some p.o. $\rho(r) \rightarrow \rho$. We will write $u \in S H\left(\mathbb{R}^{m}, \rho, \rho(r)\right)$ or shorter, $u \in S H(\rho(r))$.

For $u \in S H(\rho(r))$ set

$$
\begin{equation*}
u_{t}(x):=u\left(P_{t} x\right) t^{-\rho(t)} \tag{3.1.2.1}
\end{equation*}
$$

We will denote this transformation as $(\bullet)_{t}$.
Theorem 3.1.2.1 (Existence of Limit Set) The following holds:
els1) $\quad u_{t} \in S H(\rho(r))$ for any $t \in(0, \infty)$,
els2) the family $\left\{u_{t}\right\}$ is precompact at infinity.
I.e., for any sequence $t_{k} \rightarrow \infty$ there exists a subsequence $t_{k_{j}} \rightarrow \infty$ and a function $v \in S H\left(\mathbb{R}^{m}\right)$ such that $u_{t_{k_{j}}} \rightarrow v$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$ (see Section 2.7.1).

Proof. The functions $u_{t}$ are subharmonic by sh1) and sh5), Theorem 2.6.1.1. (Elementary Properties), and

$$
M\left(r, u_{t}\right)=M(r t, u) t^{-\rho(t)}
$$

Now we have

$$
\sigma_{M}\left[u_{t}, \rho(r)\right]=t^{-\rho(t)} \limsup _{r \rightarrow \infty} \frac{M(r t, u)}{(r t)^{\rho(r t)}} \cdot \lim _{r \rightarrow \infty} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}}=\sigma_{M}[u, \rho(r)] t^{\rho-\rho(t)}
$$

because

$$
\begin{equation*}
\lim _{r \rightarrow \infty} \frac{(r t)^{\rho(r t)}}{r^{\rho(r)}}=t^{\rho} \lim _{r \rightarrow \infty} \frac{L(r t)}{L(r)}=t^{\rho} \tag{3.1.2.2}
\end{equation*}
$$

(see, ppo3), Theorem 2.8.1.3 (Properties of P.O)). Therefore els1) is proved.
Let us check the conditions of Theorem 2.7.1.1 (Compactness in $\mathcal{D}^{\prime}$ ). We have

$$
\begin{equation*}
\limsup _{t \rightarrow \infty} M\left(r, u_{t}\right)=\limsup _{t \rightarrow \infty} \frac{M(r t, u)}{(r t)^{\rho(r t)}} \cdot \lim _{t \rightarrow \infty} \frac{(r t)^{\rho(r t)}}{t^{\rho(t)}}=\sigma_{M}[u, \rho(r)] r^{\rho} \tag{3.1.2.3}
\end{equation*}
$$

Thus, the family is bounded from above on every compact set and

$$
\lim _{t \rightarrow \infty} u_{t}(0)=\lim _{t \rightarrow \infty} u(0) t^{-\rho(t)}=0
$$

Therefore $u_{t}(0)$ are bounded from below for large $t$.

We will call the set of all functions $v$ from Theorem 3.1.2.1 the limit set of the function $u(x)$ with respect to $V_{\bullet}$ and denote it by $\operatorname{Fr}\left[u, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right]$ or shortly $\mathrm{Fr}[u]$.

The limit set does not depend on values of the subharmonic function on a bounded set, hence, it is a characteristic of asymptotic behavior.

Set

$$
\begin{gather*}
U[\rho, \sigma]:=\left\{v \in S H\left(\mathbb{R}^{m}\right): M(r, v) \leq \sigma r^{\rho}, r \in[0, \infty) ; v(0)=0\right\}, \\
U[\rho]:=\bigcup_{\sigma>0} U[\rho, \sigma] \tag{3.1.2.4}
\end{gather*}
$$

and

$$
\begin{equation*}
v_{[t]}(x):=t^{-\rho} v\left(P_{t} x\right), t \in(0, \infty) \tag{3.1.2.4a}
\end{equation*}
$$

Let us emphasize that the transformation $(\bullet)_{[t]}$ coincides with $(\bullet)_{t}$ from (3.1.2.1) for $\rho(r) \equiv \rho$ and satisfies the condition

$$
\begin{equation*}
(\bullet)_{[t \tau]}=\left((\bullet)_{[t]}\right)_{[\tau]} \tag{3.1.2.4b}
\end{equation*}
$$

Theorem 3.1.2.2 (Properties of Fr) The following holds:
fr1) $\operatorname{Fr}[u]$ is a connected compact set;
fr2) $\operatorname{Fr}[u] \subset U[\rho, \sigma]$, for $\sigma \geq \sigma_{M}[u]$;
fr3) $(\operatorname{Fr}[u])_{[t]}=\mathbf{F r}[u], t \in(0, \infty)$. I.e., $v \in \operatorname{Fr}[u]$ implies $v_{[t]} \in \mathbf{F r}[u]$;
$\mathrm{fr} 4)$ if $\rho_{1}(r)$ and $\rho(r)$ are equivalent (see (2.8.1.5)), then

$$
\mathbf{F r}\left[u, \rho_{1}(r), \bullet\right]=\mathbf{F r}[u, \rho(r), \bullet] .
$$

We need the following assertion.
Theorem 3.1.2.3 (Continuity $u_{t}$ ) The functions

$$
u_{t}, v_{[t]}:(0, \infty) \times \mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right) \mapsto \mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)
$$

are continuous in the natural topology.
Proof. For any $\psi \in \mathcal{D}\left(\mathbb{R}^{m}\right)$ consider

$$
\left\langle u_{t}, \psi\right\rangle:=\int u_{t}(x) \psi(x) d x=\int u(y) \psi(y / t) t^{m-\rho(t)} d y:=\langle u, \psi(\bullet, t)\rangle,
$$

where $\psi(y, t):=\psi(y / t) t^{m-\rho(t)}$.
The function $\psi(\bullet, t)$ is continuous in $t$ in $\mathcal{D}\left(\mathbb{R}^{m}\right)$. By Theorem 2.3.4.6 (Continuity $\langle\bullet, \bullet\rangle)\langle u, \psi(\bullet, t)\rangle$ is continuous in $(u, t)$.

Proof of Theorem 3.1.2.2. Let us denote as $\operatorname{clos}\{\bullet\}$ the closure in $\mathcal{D}^{\prime}$-topology.
The set $F_{N}:=\operatorname{clos}\left\{u_{t}: t \geq N\right\} \supset \mathbf{F r}[u]$ is compact in $\mathcal{D}^{\prime}$-topology. Indeed, let $t_{j} \rightarrow t$ and $t<\infty$; then $u_{t_{j}} \rightarrow u_{t}$ because of Theorem 3.1.2.3. If $t_{j} \rightarrow \infty$ and $u_{t_{j}} \rightarrow v$, then $v \in \operatorname{Fr}[u]$ by its definition, hence, $v \in F_{N}$. Since $\operatorname{Fr}[u]=\cap_{N=1}^{\infty} F_{N}$, it is compact.

Let us prove the connectedness. Suppose $\mathbf{F r}[u]$ is not connected. Then it can be written as a union of two disjoint nonempty closed sets $F^{1}$ and $F^{2}$. Let $V^{1}, V^{2}$ be disjoint open neighborhoods of $F^{1}, F^{2}$ respectively in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$. Since $F^{1}, F^{2}$ are nonempty there exist sequences $\left\{s_{j}\right\},\left\{t_{j}\right\}$ such that $s_{j}<t_{j}, s_{j} \rightarrow \infty, u_{s_{j}} \in$ $V^{1}, u_{t_{j}} \in V^{2}$. Since the mapping $u_{t}:(0, \infty) \mapsto \mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$ is continuous, by Theorem 3.1.2.3 its image is connected. Thus there exists a sequence $\left\{p_{j}\right\}$ with $s_{j}<p_{j}<t_{j}$ such that $u_{p_{j}} \notin V^{1} \cup V^{2}$. This sequence has a subsequence that converges to a function $v \in \operatorname{Fr}[u]$ and $v \notin F^{1} \cup F^{2}$. This is a contradiction. Hence, $\operatorname{Fr}[u]$ is connected and fr1) is proved.

Set

$$
\psi(r):=\limsup _{r \rightarrow \infty} M\left(r, u_{t}\right)
$$

This function is convex with respect to $-r^{2-m}$ for $m \geq 3$ and with respect to $\log r$ for $m=2$ and hence continuous.

Indeed, $M\left(|x|, u_{t}\right)$ are subharmonic (see Theorem 2.6.5.2 (Convexity $M(\bullet, u)$ and $\mathcal{M}(r, u)$ ). By Theorem 2.7.3.3 (H.Cartan + ) the function $\psi^{*}(|x|)$ is subharmonic and $\psi(|x|)=\psi^{*}(|x|)$ quasi-everywhere. However, if $\psi(|x|)<\psi^{*}(|x|)$ at some point, the same inequality holds on a sphere which has a positive capacity (see Example 2.5.2.2). Hence, $\psi(|x|)=\psi^{*}(|x|)$ everywhere, and $\psi(|x|)$ is subharmonic. Thus $\psi(r)$ is convex with respect to $-r^{2-m}$ for $m \geq 3$ and with respect to $\log r$ for $m=2$ by Theorem 2.6.3.2 (Subharmonicity and Convexity).

One can also see that for $u \in S H\left(\mathbb{R}^{m}\right)$,

$$
M\left(r, u_{\epsilon}\right) \leq M(r+\epsilon, u)
$$

where $(\bullet)_{\epsilon}$ is defined by (2.6.2.3).
Let $v \in \operatorname{Fr}[u]$ and $u_{t_{j}} \rightarrow v$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$. By property reg3), Theorem 2.3.4.5 $\left(u_{t_{j}}\right)_{\epsilon} \rightarrow v_{\epsilon}$ uniformly on any compact set. Thus

$$
\begin{align*}
v_{\epsilon}(x)=\lim _{j \rightarrow \infty}\left(u_{t_{j}}\right)_{\epsilon} & \leq \limsup _{t \rightarrow \infty)} M\left(|x|,\left(u_{t}\right)_{\epsilon}\right) \\
& \leq \limsup _{t \rightarrow \infty)} M\left(|x|+\epsilon, u_{t}\right)=\psi(|x|+\epsilon) . \tag{3.1.2.5}
\end{align*}
$$

If $\epsilon \downarrow 0$, then $v_{\epsilon} \downarrow v$ by Theorem 2.6.2.3 and $\psi(r+\epsilon) \rightarrow \psi(r)$ because of continuity. Passing to the limit in (3.1.2.5) and using (3.1.2.3) we obtain

$$
\begin{equation*}
v(x) \leq \sigma_{M}[u, \rho(r)]|x|^{\rho} . \tag{3.1.2.6}
\end{equation*}
$$

Since $u(0) \leq u_{\epsilon}(0)$ we have $u(0) t^{-\rho(t)} \leq\left(u_{t}\right)_{\epsilon}(0)$. Let us pass to the limit as $t:=t_{j} \rightarrow \infty$. We obtain $v_{\epsilon}(0) \geq 0$. Passing to the limit as $\epsilon \downarrow 0$ we have

$$
\begin{equation*}
v(0) \geq 0 . \tag{3.1.2.7}
\end{equation*}
$$

The inequalities (3.1.2.6) and (3.1.2.7) imply fr2).
One can check the equality

$$
\begin{equation*}
\left(u_{t}\right)_{[\tau]}=u_{t \tau} \cdot \frac{(t \tau)^{\rho(t \tau)}}{t^{\rho(t)} \tau^{\rho}} \tag{3.1.2.8}
\end{equation*}
$$

By using properties of p.o. we have

$$
\lim _{t \rightarrow \infty} \frac{(t \tau)^{\rho(t \tau)}}{t^{\rho(t)} \tau^{\rho}}=1
$$

(compare (3.1.2.2)).
Let $v \in \operatorname{Fr}[u]$ and $u_{t_{j}} \rightarrow v$. Set $t:=t_{j}, \tau:=t$ in (3.1.2.8) and pass to the limit. Then

$$
v_{[t]}=\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} u_{t_{j} t}
$$

Thus $v_{[t]} \in \operatorname{Fr}[u]$. The property f 3 ) is proved.
Let us prove f4). We have

$$
\frac{u\left(P_{t} x\right)}{t^{\rho_{1}(t)}}=\frac{u\left(P_{t} x\right)}{t^{\rho(t)}} \times e^{\left(\rho_{1}(t)-\rho(t)\right) \log t}=\frac{u\left(P_{t} x\right)}{t^{\rho(t)}} \times(1+o(1))
$$

as $t \rightarrow \infty$ because of (2.8.1.5).
This implies f4).
Exercise 3.1.2.1 Check this in detail.
We can consider the limit sets as a mapping $u \mapsto \operatorname{Fr}[u]$. The following theorem describes some properties of this mapping.

Set

$$
\begin{equation*}
U[\rho]:=\bigcup_{\sigma>0} U[\rho, \sigma] \tag{3.1.2.9}
\end{equation*}
$$

where $U[\rho, \sigma]$ is defined by (3.1.2.4).
Let $X, Y$ be subsets of a cone (i.e., a subset of a linear space that is closed with respect to sum and multiplication by a positive number). The set $U[\rho]$ is such a cone. Set

$$
\begin{equation*}
X+Y:=\{z=x+y: x \in X, y \in Y\} ; \lambda X:=\{z=\lambda x: x \in X\} \tag{3.1.2.10}
\end{equation*}
$$

Theorem 3.1.2.4 (Properties of $u \mapsto \mathbf{F r}[u]$ ) The following holds:
fru1) $\operatorname{Fr}\left[u_{1}+u_{2}\right] \subset \mathbf{F r}\left[u_{1}\right]+\operatorname{Fr}\left[u_{2}\right]$,
fru2) $\quad \operatorname{Fr}[\lambda u]=\lambda \operatorname{Fr}[u]$.

Proof. Let $v \in \operatorname{Fr}\left[u_{1}+u_{2}\right]$. Then there exists $t_{j} \rightarrow \infty$ such that $\left(u_{1}+u_{2}\right)_{t_{j}} \rightarrow v$ in $\mathcal{D}^{\prime}$. We can find a subsequence $t_{j_{k}}$ such that $\left(u_{1}\right)_{t_{j_{k}}} \rightarrow v_{1}$ and $\left(u_{2}\right)_{t_{j_{k}}} \rightarrow v_{2}$. Then $v=v_{1}+v_{2}$. The property fru1) has been proved.

The property fru2) is proved analogically.
3.1.3 We will write $\mu \in \mathcal{M}\left(\mathbb{R}^{m}, \rho(r)\right)$ or shortly, $\mu \in \mathcal{M}(\rho(r))$ if $\mu \in \mathcal{M}\left(\mathbb{R}^{m}\right)$ (see 2.8.3) and $\bar{\Delta}[\mu, \rho(r)]<\infty$ (see 2.8.3.2).

Let us define a distribution $\mu_{t}$ for $\mu \in \mathcal{M}(\rho(r))$ by

$$
\begin{equation*}
\left\langle\mu_{t}, \phi\right\rangle:=t^{-\rho(t)-m+2} \int \phi\left(P_{t}^{-1} x\right) \mu(d x) \tag{3.1.3.1}
\end{equation*}
$$

for $\phi \in \mathcal{D}\left(\mathbb{R}^{m}\right)$.
It is positive. Hence, it defines uniquely a measure $\mu_{t}$.
Theorem 3.1.3.1 (Explicit form of $\left.\mu_{t}\right)$ For any $E \in \sigma\left(\mathbb{R}^{m}\right)$ the following holds:

$$
\begin{equation*}
\mu_{t}(E)=t^{-\rho(t)-m+2} \mu\left(P_{t} E\right) \tag{3.1.3.2}
\end{equation*}
$$

Proof. It is enough to prove the assertion for some dense ring (see Theorem 2.2.3.5), for example, for all compact sets.

Let $\chi_{K}$ be a characteristic function of a compact set $K$ and let $\phi_{\epsilon} \downarrow \chi_{K}$ be a monotonically converging sequence of functions that belong to $\mathcal{D}\left(\mathbb{R}^{m}\right)$ (see Theorems 2.1.2.1, 2.1.2.9 and 2.3.4.4). Then

$$
\int \phi_{\epsilon}(x) \mu_{t}(d x)=t^{-\rho(t)-m+2} \int \phi_{\epsilon}\left(P_{t}^{-1} x\right) \mu(d x) .
$$

Since $\phi_{\epsilon}\left(P_{t}^{-1} x\right) \downarrow \chi_{P_{t} K}(x)$,

$$
\mu_{t}(K)=\int \chi_{K}(x) \mu_{t}(d x)=t^{-\rho(t)-m+2} \int \chi_{P_{t} K}(x) \mu(d x)=t^{-\rho(t)-m+2} \mu\left(P_{t} K\right)
$$

Theorem 3.1.3.2 (Existence of $\mu$-Limit Set) The following holds:
mels1) $\mu_{t} \in \mathcal{M}(\rho(r))$ for any $t \in(0, \infty)$;
mels2) the family $\left\{\mu_{t}\right\}$ is precompact in infinity.
I.e., for any sequence $t_{k} \rightarrow \infty$ there exists a subsequence $t_{k_{j}} \rightarrow \infty$ and a measure $\nu \in \mathcal{M}\left(\mathbb{R}^{m}\right)$ such that $\mu_{t_{k_{j}}} \rightarrow \nu$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$ (see Section 2.7.1).
Proof. We have

$$
\mu_{t}(r)=\mu(r t) t^{-\rho(t)-m+2}
$$

Thus

$$
\begin{align*}
\limsup _{r \rightarrow \infty} \frac{\mu_{t}(r)}{r^{\rho(r)+m-2}} & =\limsup _{r \rightarrow \infty} \frac{\mu(r t)}{(r t)^{\rho(r t)+m-2}} \frac{(r t)^{\rho(r t)}}{t^{\rho(t)+m-2} r^{\rho(r)}} \\
& =t^{\rho-\rho(t)-(m-2)} \bar{\Delta}[\mu, \rho(r)] \tag{3.1.3.3}
\end{align*}
$$

Therefore mels1) holds.

We also have

$$
\limsup _{t \rightarrow \infty} \mu_{t}(r)=\bar{\Delta}[\mu, \rho(r)] r^{\rho+m-2}
$$

Thus $\mu_{t}$ satisfies the assumption of the Helly theorem (Theorem 2.2.3.2). Using also Theorem 2.3.4.4 we obtain mels2).

We will call the set of all measures $\nu$ from Theorem 3.1.2.1 the limit set of the mass distribution $\mu$ with respect to $V_{\bullet}$ and denote it $\operatorname{Fr}\left[\mu, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right]$ or shortly, $\operatorname{Fr}[\mu]$.

Set

$$
\begin{gather*}
\mathcal{M}[\rho, \Delta]:=\left\{\nu: \nu(r) \leq \Delta r^{\rho+m-2}, \forall r>0\right\} .  \tag{3.1.3.4}\\
\mathcal{M}[\rho]:=\bigcup_{\Delta>0} \mathcal{M}[\rho, \Delta],
\end{gather*}
$$

and

$$
\begin{equation*}
\nu_{[t]}(E):=t^{-\rho-m+2} \nu\left(P_{t} E\right) \tag{3.1.3.5}
\end{equation*}
$$

for $E \in \sigma\left(\mathbb{R}^{m}\right)$.
Theorem 3.1.3.3 (Properties of $\operatorname{Fr}[\mu]$ ) The following holds:
frm1) $\operatorname{Fr}[\mu]$ is connected and compact;
frm2) $\operatorname{Fr}[\mu] \subset \mathcal{M}[\Delta, \rho]$, for $\Delta \geq \bar{\Delta}[\mu, \rho(r)]$;
frm3) $(\mathbf{F r}[\mu])_{[t]}=\mathbf{F r}[\mu], t \in(0, \infty)$.
Proof. We will only prove frm2) because frm1) and frm3) are proved word by word as in Theorem 3.1.2.2.

Suppose $t_{n} \rightarrow \infty$ and $\mu_{t_{n}} \rightarrow \nu \in \operatorname{Fr}[\mu]$. Let us choose $r^{\prime}>r$ such that the open ball $K_{r^{\prime}}$ is squarable with respect to $\nu$. It is possible because of Theorem 2.2.3.3, sqr2). By Theorems 2.2.3.7 and 2.3.4.4, $\mu_{t_{n}}\left(r^{\prime}\right) \rightarrow \nu\left(r^{\prime}\right)$. Thus (compare with (3.1.2.3))

$$
\nu\left(r^{\prime}\right)=\lim _{t_{n} \rightarrow \infty} \mu_{t_{n}}\left(r^{\prime}\right) \leq \limsup _{t \rightarrow \infty} \mu_{t}\left(r^{\prime}\right)=\bar{\Delta}[\mu, \rho(r)]\left(r^{\prime}\right)^{\rho+m-2} .
$$

Choosing $r^{\prime} \downarrow r$ we obtain

$$
\lim _{r^{\prime} \rightarrow r} \nu\left(r^{\prime}\right)=\nu(r)
$$

because (2.2.3.3). Thus frm2) holds.
The following assertion is a "copy" of Theorem 3.1.2.4.
Theorem 3.1.3.4 (Properties of $\mu \mapsto \operatorname{Fr}[\mu]$ ) The following holds:
frmu1) $\operatorname{Fr}\left[\mu_{1}+\mu_{2}\right] \subset \mathbf{F r}\left[\mu_{1}\right]+\mathbf{F r}\left[\mu_{2}\right]$,
frmu2) $\operatorname{Fr}[\lambda \mu]=\lambda \mathbf{F r}[\mu]$.
The proof is also a "copy" and we omit it.
Exercise 3.1.3.1 Prove Theorem 3.1.3.4
3.1.4 We are going to study the class $U[\rho]$ and obtain for it "non-asymptotic" analogies of Theorem 2.8.3.3 (Jensen), 2.9.2.3 (Brelot-Borel), 2.9.3.1 (Brelot-Hadamard)

Theorem 3.1.4.1 (*Jensen) Let $v \in U[\rho]$. Then its Riesz measure $\nu_{v} \in \mathcal{M}[\rho]$.
Proof. As in Theorem 2.8.3.2 we have an inequality

$$
\begin{equation*}
\frac{\nu_{v}(r)}{r^{m-2}} \leq A(m) N\left(2 r, \nu_{v}\right) \tag{3.1.4.1}
\end{equation*}
$$

Since $v(0)=0$ we have (Theorem 2.6.5.1. (Jensen-Privalov))

$$
\begin{equation*}
N\left(2 r, \nu_{v}\right)=\mathcal{M}(2 r, v) \leq M(2 r, v) \leq 2^{\rho} \sigma r^{\rho} . \tag{3.1.4.2}
\end{equation*}
$$

Substituting (3.1.4.2) in (3.1.4.1) we obtain $\nu_{v} \in \mathcal{M}[\rho, \Delta]$ for some $\Delta$. Thus $\nu_{v} \in$ $\mathcal{M}[\rho]$.

Let $\rho$ be non-integer and $\nu \in \mathcal{M}[\rho]$. Consider the canonical potential $\Pi(x, \nu, p)$ where $p:=[\rho]$ (see (2.9.2.1)). Let us emphasize that the support of $\nu$ may contain the origin but $\nu(0)=0$, i.e., there is no concentrated mass in the origin. Thus we must also check its convergence in the origin.

Theorem 3.1.4.2 (*Brelot-Borel) Let $\rho$ be non-integer and let $\nu \in \mathcal{M}[\rho]$. Then $\Pi(x, \nu, p)$ converges and belongs to $U[\rho]$.

Proof. Using (2.9.1.9) we have

$$
\begin{equation*}
\left|\int_{|y|<2|x|} G_{p}(x, y, m) \nu(d y)\right| \leq A(m, p)|x|^{p} \int_{0}^{2|x|} \frac{\nu(d t)}{t^{p+m-2}} . \tag{3.1.4.3}
\end{equation*}
$$

Let us estimate the integral in (3.1.4.3). Integrating by parts we obtain

$$
I_{<}(x):=\int_{0}^{2|x|} \frac{\nu(d t)}{t^{p+m-2}}=\left.\frac{\nu(t)}{t^{p+m-2}}\right|_{0} ^{2|x|}+(p+m-2) \int_{0}^{2|x|} \frac{\nu(t) d t}{t^{p+m-1}} .
$$

Since $\nu \in \mathcal{M}[\rho, \Delta]$ for some $\Delta$,

$$
I_{<}(x) \leq A(m, \rho, p) \Delta|x|^{\rho-p}
$$

Substituting this in (3.1.4.3) we obtain

$$
\begin{equation*}
\left|\int_{|y|<2|x|} G_{p}(x, y, m) \nu(d y)\right| \leq A(m, \rho, p) \Delta|x|^{\rho} . \tag{3.1.4.4}
\end{equation*}
$$

Analogously, using (2.9.1.8) we obtain

$$
\begin{equation*}
\left|\int_{|x|<2|y|} G_{p}(x, y, m) d \nu(d y)\right| \leq A(m, \rho, p) \Delta|x|^{\rho} . \tag{3.1.4.5}
\end{equation*}
$$

In particular, these estimates show that $\Pi(x, \nu, p)$ exists. Now using (2.9.1.10) we have also

$$
\int_{\frac{|x|}{2} \leq|y| \leq 2|x|} G_{p}(x, y, m) \nu(d y) \leq A(m, p) \int_{\frac{|x|}{2}}^{2|x|} \nu(d t) \min \left(\frac{|x|^{p+1}}{t^{p+m-1}}, \frac{|x|^{p}}{t^{p+m-2}}\right) .
$$

The latter integral can also be easily estimated by $\Delta A(m, p, \rho)|x|^{\rho}$. Thus we have

$$
\int_{\frac{|x|}{2} \leq|y| \leq 2|x|} G_{p}(x, y, m) \nu(d y) \leq A(m, p, \rho)|x|^{\rho}
$$

Therefore by (3.1.4.5) and (3.1.4.4) we obtain $M(r, \Pi) \leq \sigma r^{\rho}$ for some $\sigma$.
Since $G_{p}(0, y, m)=0$ for all $y \neq 0$ and the integral converges, $\Pi(0, \nu, p)=0$.
We will need an assertion that looks like the Liouville theorem (Theorem 2.4.2.3).

Theorem 3.1.4.3 (*Liouville) Let $H$ be a harmonic function in $\mathbb{R}^{m}$ and $H \in U[\rho]$. Then $H \equiv 0$ if $\rho$ is non-integer and $H$ is a homogeneous polynomial of degree $p$ if $\rho=p$ is integer.

In particular, for $m=2$ we have $H\left(r e^{i \phi}\right)=r^{p} \Re\left(c e^{i p \phi}\right)$.
Proof. Like in the proof of the Liouville theorem we obtain the inequality (2.4.2.9) and

$$
\left|c_{k}\right| \leq A R^{-k} \max _{|x|=R} H(x) \leq A \sigma R^{\rho-k}
$$

for some $\sigma>0$.
If $k>\rho$, we will pass to the limit when $R \rightarrow \infty$ and obtain $c_{k}=0$. If $k<\rho$, we will do that when $R \rightarrow 0$ and obtain $c_{k}=0$.

The following theorem can be considered as an analogy of the Brelot-Hadamard theorem (Theorem 2.9.3.1):

Theorem 3.1.4.4 (*Hadamard) Let $\rho$ be non-integer and $v \in U[\rho]$. Then

$$
\begin{equation*}
v(x)=\Pi\left(x, \nu_{v}, p\right) \tag{3.1.4.6}
\end{equation*}
$$

for $p=[\rho]$.
Proof. Consider the function $H(x):=v(x)-\Pi\left(x, \nu_{v}, p\right)$. It is harmonic. We also have by(2.8.2.5)

$$
M(r, H) \leq A(m) T(r, H) \leq A(m)[T(r, v)+T(r, \Pi)] \leq \sigma r^{\rho}
$$

for some $\sigma$.
Hence, $H(x) \equiv 0$ by Theorem 3.1.4.3.

Let us consider the case of integer $\rho$. Let $\nu \in \mathcal{M}[\rho]$ for an integer $\rho=p$. Set

$$
\begin{align*}
& \Pi_{<}(x, \nu, \rho):=\int_{|y|<1} G_{p-1}(x, y, m) \nu(d y)  \tag{3.1.4.7}\\
& \Pi_{>}(x, \nu, \rho):=\int_{|y| \geq 1} G_{p}(x, y, m) \nu(d y) . \tag{3.1.4.8}
\end{align*}
$$

Both potentials converge and belong to $U[\rho]$.
Theorem 3.1.4.5 (**Hadamard) Let $\rho$ be integer and let $v \in U[\rho]$. Then

$$
\begin{equation*}
v=H_{\rho}(x)+\Pi_{<}(x, \nu, \rho)+\Pi_{>}(x, \nu, \rho), \tag{3.1.4.9}
\end{equation*}
$$

where $H_{\rho}$ is a homogeneous harmonic polynomial of degree $\rho$.
The proof is exactly the same as in the *Hadamard theorem, but we use the second case of Theorem 3.1.4.3. We also note that the polynomial may be equal to zero identically.

Exercise 3.1.4.1 Check this in detail.
Let as check that $\nu$ from (3.1.4.9) has the following property that is analogous to Theorem 2.9.4.2.

Theorem 3.1.4.6 (*Lindelöf) Let $\rho$ be integer and let $v \in U[\rho]$. Then

$$
\begin{equation*}
\lim _{\epsilon \rightarrow 0} \int_{\epsilon \leq|y|<1} D_{\rho}(x, y) \mu(d y)=H_{\rho}(x) \tag{3.1.4.10}
\end{equation*}
$$

Proof. Consider the function

$$
v_{\epsilon}^{*}(x):=v(x)+\left\{\begin{array}{l}
\int_{|y|<\epsilon} \frac{\nu(d y)}{x-\left.y\right|^{m-2}} \text { for } m>2  \tag{3.1.4.11}\\
-\int_{|y|<\epsilon} \log |x-y| \nu(d y), \text { for } m=2
\end{array}\right.
$$

It is subharmonic with supp $\nu \cap\{0\}=\varnothing$. We represent this function as in (2.9.4.10) in the form

$$
v_{\epsilon}^{*}(x)=\Pi_{<}^{1}\left(x, \nu_{\epsilon}^{*}, \rho\right)+\Pi_{>}^{1}\left(x, \nu_{\epsilon}^{*}, \rho\right)+P_{\rho-1}^{*}\left(x, v_{\epsilon}^{*}\right)+\delta_{1}\left(x, v_{\epsilon}^{*}, \rho\right) .
$$

In this representation we can pass to the limit as $\epsilon \rightarrow 0$ in the left side and in all the summands except perhaps the last two from the right side.

Exercise 3.1.4.2 Check this, using that all the integrals converge for $\nu \in \mathcal{M}(\rho, \Delta)$ and showing that the integral in (3.1.4.11) tends to zero.

The last two summands form a harmonic polynomial, the limit of which is also a harmonic polynomial. Comparing the limit with the representation (3.1.4.9), we obtain that $P_{\rho-1}^{*}\left(\bullet, v_{\epsilon}^{*}\right)$ tends to zero and $\delta_{1}\left(x, v_{\epsilon}^{*}, \rho\right)$ tends to $H_{\rho}(x)$.

Theorem 3.1.4.7 (**Liouville) If $v \in U[\rho]$ satisfies inequality $v(x) \leq 0$ for $z \in \mathbb{R}^{m}$, then $v(x) \equiv 0$.

Otherwise it contradicts subharmonicity in 0 .
3.1.5 Let us study the connection between $\operatorname{Fr}[u]$ and $\operatorname{Fr}\left[\mu_{u}\right]$.

Note the following properties of the transformations $(\bullet)_{t}$ and $(\bullet)_{[t]}$.
Theorem 3.1.5.0 (Connection between $u_{t}$ and $\mu_{t}$ ) One has

$$
\begin{equation*}
\left(\mu_{u}\right)_{t}=\mu_{u_{t}} ; \quad\left(\mu_{v}\right)_{[t]}=\mu_{v_{[t]}} . \tag{3.1.5.0}
\end{equation*}
$$

Proof. By the F. Riesz theorem (Theorem 2.6.4.3) and Theorem 2.5.1.1, GPo3) we have for any $\psi \in \mathcal{D}\left(\mathbb{R}^{m}\right)$,

$$
\left\langle\mu_{u}, \psi\right\rangle=\theta_{m}\langle\Delta u, \psi\rangle=\theta_{m}\langle u, \Delta \psi\rangle .
$$

Using the definition (3.1.3.1), we obtain

$$
\left\langle\left(\mu_{u}\right)_{t}, \psi\right\rangle=\left\langle\left(\mu_{u}\right)_{t}, \psi\left(\left(P_{t}\right)^{-1} \bullet\right)\right\rangle t^{-\rho(t)-m+2}
$$

Thus

$$
\left\langle\left(\mu_{u}\right)_{t}, \psi\right\rangle=\theta_{m}\left\langle u, \Delta\left[\psi\left(\left(P_{t}\right)^{-1} \bullet\right)\right]\right\rangle t^{-\rho(t)-m+2}
$$

Since the Laplace operator is invariant with respect to $V_{t}$ for any $t$ we have

$$
\Delta\left[\psi\left(\left(P_{t}\right)^{-1} \bullet\right)\right]=t^{-2}[\Delta \psi]\left(\left(P_{t}\right)^{-1} \bullet\right) .
$$

Thus we obtain

$$
\begin{aligned}
\left\langle\left(\mu_{u}\right)_{t}, \psi\right\rangle & =\theta_{m} t^{-\rho(t)-m}\left\langle u,[\Delta \psi]\left(\left(P_{t}\right)^{-1} \bullet\right)\right\rangle \\
& =\theta_{m}\left\langle u\left(P_{t} \bullet\right) t^{-\rho(t)}, \Delta \psi\right\rangle=\theta_{m}\left\langle u_{t}, \Delta \psi\right\rangle=\left\langle\mu_{u_{t}}, \psi\right\rangle .
\end{aligned}
$$

Exercise 3.1.5.1 Do this for $(\bullet)_{[t]}$.
We begin from the case of a non-integer $\rho$.
Theorem 3.1.5.1 (Connection between Fr's for non-integer $\rho$ ) Let $u \in U(\rho(r))$ and $\mu_{u}$ be its Riesz measure. Then

$$
\begin{align*}
\operatorname{Fr}\left[\mu_{u}\right] & =\left\{\nu_{v}: v \in \mathbf{F r}[u]\right\},  \tag{3.1.5.1}\\
\operatorname{Fr}[u] & =\left\{\Pi(\bullet, \nu, p): \nu \in \mathbf{F r}\left[\mu_{u}\right]\right\} . \tag{3.1.5.2}
\end{align*}
$$

Proof. Let $\nu \in \operatorname{Fr}\left[\mu_{u}\right]$. There exists $t_{n} \rightarrow \infty$ such that $\left(\mu_{u}\right)_{t_{n}} \rightarrow \nu$ in $\mathcal{D}^{\prime}$. We can find a subsequence $t_{n}^{\prime}$ such that $u_{t_{n}^{\prime}} \rightarrow v \in \operatorname{Fr}[u]$. Thus $\left(\mu_{u}\right)_{t_{n}^{\prime}} \rightarrow \nu_{v}$ and therefore
$\nu=\nu_{v}$. Hence, $\operatorname{Fr}\left[\mu_{u}\right] \subset\left\{\nu_{v}: v \in \operatorname{Fr}[u]\right\}$. Analogously we can prove that every $\nu_{v} \in \operatorname{Fr}\left[\mu_{u}\right]$ and hence (3.1.5.1) holds.

Let $\nu \in \operatorname{Fr}\left[\mu_{u}\right]$. We find a sequence $t_{n} \rightarrow \infty$ such that $\left(\mu_{u}\right)_{t_{n}} \rightarrow \nu$ in $\mathcal{D}^{\prime}$. We find a subsequence $t_{n}^{\prime}$ such that $u_{t_{n}^{\prime}} \rightarrow v \in \mathbf{F r}[u]$ and $\nu_{v}=\nu$. By the *Hadamard theorem (Theorem 3.1.4.4) $v=\Pi(\bullet, \nu, p)$. Hence, $\left\{\Pi(\bullet, \nu, p): \nu \in \operatorname{Fr}\left[\mu_{u}\right]\right\} \subset \operatorname{Fr}[u]$. And vice versa, since $\operatorname{Fr}[u] \subset U[\rho]$ (Theorem 3.1.2.2, fr2)), every $v \in \operatorname{Fr}[u]$ is represented as $\Pi\left(\bullet, \nu_{v}, p\right)$ and $\nu_{v} \in \operatorname{Fr}[\mu]$ by (3.1.5.1).

Let $\rho$ be integer and $u \in U(\rho(r))$. Let us consider the precompact family of homogeneous polynomials $\delta_{t}(x, u, \rho) t^{\rho-\rho(t)}$ from Theorem 2.9.4.1. For every $t_{n} \rightarrow$ $\infty$ we can find a subsequence $t_{n}^{\prime}$ such that the pair $\left(\delta_{t_{n}^{\prime}}(\bullet, u, \rho) t_{n}^{\prime}{ }^{\rho-\rho\left(t_{n}^{\prime}\right)},\left(\mu_{u}\right)_{t_{n}^{\prime}}\right)$ tends to a pair $\left(H_{\nu}, \nu\right)$ where $H_{\nu}$ is a homogeneous harmonic polynomial of degree $p$. We denote the set of all such pairs as $(\mathcal{H}, \mathbf{F r})[u]$. Every $v \in U[\rho]$ can be represented in the form (3.1.4.7). Thus for every $v$ the polynomial $H^{v}:=H_{p}$ is determined.

Theorem 3.1.5.2 (Connection between Fr's for integer $\rho$ ) Let $u \in U(\rho(r))$. Then

$$
\begin{gather*}
(\mathcal{H}, \mathbf{F r})[u]=\left\{\left(H^{v}, \nu_{v}\right): v \in \mathbf{F r}[u]\right\}  \tag{3.1.5.3}\\
\operatorname{Fr}[u]=\left\{v:=H_{\nu}+\Pi_{<}(\bullet, \nu, \rho)+\Pi_{>}(\bullet, \nu, \rho):\left(H_{\nu}, \nu\right) \in(\mathcal{H}, \mathbf{F r})[u]\right\} . \tag{3.1.5.4}
\end{gather*}
$$

The proof is clear.
3.1.6 Up to now we supposed that the family of rotations $V_{\bullet}$ was fixed. Now we take in consideration that it can vary and use the notation $\mathbf{F r}\left[u, V_{\mathbf{0}}\right]$.

Theorem 3.1.6.1 (Dependence of $\mathbf{F r}$ on $V_{\bullet}$ ) Let $\operatorname{Fr}\left[u, V_{\bullet}\right]$ and $\operatorname{Fr}\left[u, W_{\bullet}\right]$ be limit sets of $u$ with respect to rotation families $V_{\bullet}$ and $W_{\bullet}$ accordingly. Then for any $v \in \mathbf{F r}\left[u, V_{\bullet}\right]$ there exist a rotation $V^{v}$ and $w^{v} \in \mathbf{F r}\left[u, W_{\bullet}\right]$ such that

$$
v(x)=w^{v}\left(V^{v} x\right)
$$

for all $x \in \mathbb{R}^{m}$.
Proof. Let $v \in \mathbf{F r}\left[u, V_{\bullet}\right]$ and let $t_{n} \rightarrow \infty$ be a sequence such that

$$
t_{n}^{-\rho\left(t_{n}\right)} u\left(t_{n} V_{t_{n}} \bullet\right) \rightarrow v
$$

Since the family $V_{t}$ is obviously precompact there exists a subsequence (for which we keep the same notation), and a rotation $V^{v}$ such that $W_{t_{n}}^{-1} V_{t_{n}} \rightarrow V^{v}$ and $w \in \operatorname{Fr}\left[u, W_{\bullet}\right]$ such that $t_{n}^{-\rho\left(t_{n}\right)} u\left(t_{n} W_{t_{n}} \bullet\right) \rightarrow w$.

Now we have

$$
\begin{aligned}
v(\bullet) & =\mathcal{D}^{\prime}-\lim t_{n}^{-\rho\left(t_{n}\right)} u\left(t_{n} V_{t_{n}} \bullet\right) \\
& =\mathcal{D}^{\prime}-\lim t_{n}^{-\rho\left(t_{n}\right)} u\left(t_{n} W_{t_{n}} W_{t_{n}}^{-1} V_{t_{n}} \bullet\right) \\
& =w\left(V^{v} \bullet\right) .
\end{aligned}
$$

### 3.2 Indicators

3.2.1 Let $u \in S H(\rho(r))$ and let $\mathbf{F r}[u]$ be the limit set. Set

$$
\begin{align*}
h(x, u) & :=\sup \{v(x): v \in \mathbf{F r}[u]\},  \tag{3.2.1.1}\\
\underline{h}(x, u) & :=\inf \{v(x): v \in \operatorname{Fr}[u]\} . \tag{3.2.1.2}
\end{align*}
$$

These functionals reflect the asymptotic behavior of $u$ along curves of the form

$$
\begin{equation*}
l_{\boldsymbol{x}^{0}}:=\left\{x=t V_{t} \boldsymbol{x}^{0}: t \in(0, \infty)\right\} \tag{3.2.1.3}
\end{equation*}
$$

and are called indicator of growth of $u$ and lower indicator respectively.
Of course, the indicators depend on $\rho(r)$ and $V_{t}$, but we will only note that if necessary.

Theorem 3.2.1.1 (Properties of Indicators) The following holds:
h1) $\underline{h}$ is upper semicontinuous, $h$ is subharmonic;
h2) they are semiadditive and positively homogeneous, i.e.,

$$
\begin{align*}
h\left(x, u_{1}+u_{2}\right) & \leq h\left(x, u_{1}\right)+h\left(x, u_{2}\right)  \tag{3.2.1.4}\\
\underline{h}\left(x, u_{1}+u_{2}\right) & \geq \underline{h}\left(x, u_{1}\right)+\underline{h}\left(x, u_{2}\right)  \tag{3.2.1.5}\\
h, \underline{h}(x, C u) & =C h, \underline{h}(x, u), C \geq 0 \tag{3.2.1.6}
\end{align*}
$$

h3) invariance:

$$
\begin{equation*}
h, \underline{h}_{[t]}(x, u)=h, \underline{h}(x, u) . \tag{3.2.1.7}
\end{equation*}
$$

Proof. Semicontinuity of $\underline{h}$ follows from Theorem 2.1.2.8 (Commutativity of inf and $\mathrm{M}($.$) . Semicontinuity and subharmonicity of h$ follow from Theorem 2.7.3.4 (Sigurdsson's Lemma). The properties h2) follow from properties of infimum and supremum. The invariance follows from invariance of $\mathbf{F r}[u]$ (Theorem 3.1.2.2, fr3)).

Set

$$
\begin{equation*}
x^{0}(x):=P_{|x|}^{-1}(x) \tag{3.2.1.8}
\end{equation*}
$$

where $P_{t}$ is defined by (3.1.1.1).
This is an intersection of the orbit of $P_{t}$ that passes through a point $x$ with the unit sphere.

$$
\begin{align*}
& \text { If } V_{t} \equiv I, \\
& \qquad x^{0}(x)=x /|x|:=x^{0} \tag{3.2.1.9}
\end{align*}
$$

Theorem 3.2.1.2 (Homogeneity $h, \underline{h})$ One has

$$
\begin{equation*}
h, \underline{h}(x, \bullet)=|x|^{\rho} h, \underline{h}\left(x^{0}(x), \bullet\right) \tag{3.2.1.10}
\end{equation*}
$$

Thus the indicators are determined uniquely by their values on the unit sphere, i.e., they are "functions of direction". In particular, they are homogeneous for $V_{t} \equiv I$ :

$$
\begin{equation*}
h, \underline{h}(x, \bullet)=|x|^{\rho} h, \underline{h}\left(x^{0}, \bullet\right) . \tag{3.2.1.11}
\end{equation*}
$$

The proof of (3.2.1.10) follows from h3), Theorem 3.2.1.1 if we set $t:=|x| ; x:=$ $P_{t}^{-1} x$.
3.2.2 In this item we will suppose that $V_{t} \equiv I$ and study the indicator.

Let $\Delta_{\boldsymbol{x}^{0}}$ be as defined in Section 2.4.1. Its coefficients depend on a choice of the spherical coordinate system. However, one has
Theorem 3.2.2.1 Let $\psi(y)$ have continuous second derivatives on the unit sphere $S_{1}$. Then the differential form $\Delta_{\boldsymbol{x}^{0}} \psi(y) d y$ is invariant with respect to the choice of spherical coordinate system.

Proof. Let $\phi(x)$ be a smooth function in $\mathbb{R}^{m}$. Then $\Delta \phi(x) d x$ is invariant with respect to the choice of an orthogonal system because $\Delta$ (the Laplace operator) and an element of volume are invariant. Set $\phi(x)=\psi(y)$, where $y:=x^{0}=x /|x|$. Then

$$
\Delta \phi d x=\Delta_{\boldsymbol{x}^{0}} \psi(y) d y r^{m-3} d r .
$$

Since $r$ is invariant with respect to rotations of the coordinate system, $\Delta_{\boldsymbol{x}^{0}} \psi(y) d y$ is invariant with respect to the choice of a spherical coordinate system.

Note that for $m=2$ this theorem is obvious because

$$
\begin{equation*}
\Delta_{\boldsymbol{x}^{0}}=\frac{d^{2}}{d \theta^{2}} \tag{3.2.2.1}
\end{equation*}
$$

and it does not depend on translations with respect to $\theta$.
We define the operator $\Delta_{\boldsymbol{x}^{0}}$ on $f \in \mathcal{D}^{\prime}\left(S_{1}\right)$ by

$$
\left\langle\Delta_{\boldsymbol{x}^{0}} f, \psi\right\rangle:=\left\langle f, \Delta_{\boldsymbol{x}^{0}} \psi\right\rangle, \psi \in \mathcal{D}\left(S_{1}\right)
$$

in a fixed spherical coordinate system.
The definition is correct. Indeed, suppose in a fixed system

$$
\begin{equation*}
\operatorname{supp} \psi \subset S_{1} \backslash\left\{\theta_{j}=0 ; \pi: j=1,2, \ldots, m-2\right\} \tag{3.2.2.2}
\end{equation*}
$$

Then all the coefficients of $\Delta_{\boldsymbol{x}^{0}}$ are infinitely differentiable and $\Delta_{\boldsymbol{x}^{0}} \psi(y) \in \mathcal{D}\left(S_{1}\right)$. By Theorem 3.2.2.1 we obtain that the condition of Theorem 2.3.5.2 ( $\mathcal{D}^{\prime}$ on Sphere) are fulfilled.

Note that for $m=2$ the operator $\Delta_{\boldsymbol{x}^{0}}$ is realized by the formula (3.2.2.1) on functions of the form $f=f\left(e^{i \theta}\right)$, i.e., on $2 \pi$-periodic functions.

Theorem 3.2.2.2 (Subsphericality of Indicator) One has

$$
\begin{equation*}
\left[\Delta_{\boldsymbol{x}^{0}}+\rho(\rho+m-2)\right] h(y, u):=s>0 \tag{3.2.2.3}
\end{equation*}
$$

in $\mathcal{D}^{\prime}\left(S_{1}\right)$.
I.e., $s$ is a measure on $S_{1}$.

Proof. It is sufficient to prove this locally, in any spherical system. Let $R(r)$ be finite, infinitely differentiable and nonnegative in $(0 ; \infty)$ and let $\psi \in \mathcal{D}\left(S_{1}\right)$ be nonnegative and satisfy (3.2.2.2). Set $\phi(x):=R(|x|) \psi\left(x^{0}\right)$. Using the subharmonicity of $h(x, u)$ (h1), Theorem 3.2.1.1 and (3.2.2.2), we have

$$
\begin{aligned}
0 & \leq \int h(x, u) \Delta \phi(x) d x \\
& =\int_{(y, r) \in S_{1} \times(0 ; \infty)} r^{\rho} h(y, u)\left[\frac{1}{r^{m-1}} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r}+\frac{1}{r^{2}} \Delta_{x^{0}}\right] \psi(y) r^{m-1} d y d r
\end{aligned}
$$

Transforming the last integral we obtain

$$
\begin{align*}
\int h(x, u) \Delta \phi(x) d x= & \int_{0}^{\infty} r^{\rho}\left[\frac{1}{r^{m-1}} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r} R(r)\right] r^{m-1} d r \int_{S_{1}} h(y, u) \psi(y) d y \\
& +\int_{0}^{\infty} r^{\rho-2} r^{m-1} R(r) d r \int_{S_{1}} h(y, u) \Delta_{\boldsymbol{x}^{0}} \psi(y) d y \tag{3.2.2.4}
\end{align*}
$$

Integrating by parts in the first summand we obtain

$$
\begin{equation*}
\int_{0}^{\infty} r^{\rho}\left[\frac{1}{r^{m-1}} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r} R(r)\right] r^{m-1} d r=\int_{0}^{\infty} R(r) \rho(\rho+m-2) r^{\rho+m-3} d r \tag{3.2.2.5}
\end{equation*}
$$

Substituting (3.2.2.5) into (3.2.2.4), we have

$$
0 \leq \int_{0}^{\infty} r^{\rho+m-3} R(r) d r \int_{S_{1}} h(y, u)\left[\Delta_{\boldsymbol{x}^{0}}+\rho(\rho+m-2)\right] \psi(y) d y
$$

Since $R(r)$ is an arbitrarily nonnegative function,

$$
\int_{S_{1}} h(y, u)\left[\Delta_{\boldsymbol{x}^{0}}+\rho(\rho+m-2)\right] \psi(y) d y \geq 0
$$

for arbitrary $\psi$.
We will call an upper semicontinuous function which satisfies (3.2.2.3) a $\rho$ subspherical one. Now we are going to study properties of these functions.
3.2.3 We consider the case $m=2$. A $\rho$-subspherical function for $m=2$ is called $\rho$-trigonometrically convex ( $\rho$-t.c.). We will obtain for such a function a representation like in Theorems 3.1.4.4, 3.1.4.5.(*,** Hadamard). Set

$$
T_{\rho}:=\frac{d^{2}}{d \phi^{2}}+\rho^{2} .
$$

Let us find a fundamental solution of this operator. Let $\rho$ be non-integer. Let us denote as $\widetilde{\cos \rho}(\phi)$ the periodic continuation of $\cos \rho \phi$ from the interval $(-\pi, \pi)$.

Theorem 3.2.3.1 (Fundamental Solution of $T_{\rho}$ ) One has

$$
\frac{1}{2 \rho \sin \pi \rho} T_{\rho} \widetilde{\cos \rho}(\phi-\pi)=\delta(\phi) \text { in } \mathcal{D}^{\prime}\left(S_{1}\right) .
$$

Proof. Let $f \in \mathcal{D}\left(S_{1}\right)$. We have

$$
\begin{equation*}
\int_{0}^{2 \pi} \widetilde{\cos \rho}(\phi-\pi)\left[f^{\prime \prime}+\rho^{2} f\right] d \phi=\lim _{\epsilon \rightarrow 0} \int_{\epsilon}^{2 \pi-\epsilon} \cos \rho(\phi-\pi)\left[f^{\prime \prime}+\rho^{2} f\right] d \phi \tag{3.2.3.1}
\end{equation*}
$$

Integrating by parts we obtain

$$
\begin{aligned}
& \int_{\epsilon}^{2 \pi-\epsilon} \cos \rho(\phi-\pi)\left[f^{\prime \prime}+\rho^{2} f\right] d \phi \\
& =\left.\cos \rho(\phi-\pi) f^{\prime}(\phi)\right|_{\epsilon} ^{2 \pi-\epsilon}+\left.\rho \sin \rho(\phi-\pi) f(\phi)\right|_{\epsilon} ^{2 \pi-\epsilon}+\int_{\epsilon}^{2 \pi-\epsilon} f(\phi) T_{\rho} \cos \rho(\phi-\pi) d \phi
\end{aligned}
$$

However, $T_{\rho} \cos \rho(\phi-\pi)=0$ for $\phi \in(\epsilon, 2 \pi-\epsilon)$. Thus the limit in (3.2.3.1) is equal to $f(0) 2 \rho \sin \pi \rho$.

Let $s$ be a measure on the circle $S_{1}$. Set

$$
\Pi(\phi, s):=\int_{0}^{2 \pi} \widetilde{\cos \rho}(\phi-\psi-\pi) s(d \psi) .
$$

Theorem 3.2.3.2 One has

$$
T_{\rho} \Pi(\phi, s)=(2 \rho \sin \pi \rho) s(\bullet) \text { in } \mathcal{D}^{\prime}\left(S_{1}\right) .
$$

The proof is the same as GPo3) in Theorem 2.5.1.1.
Theorem 3.2.3.3 (Representation of $\rho$-t.c.f for a non-integer $\rho$ ) Let h be $\rho$-t.c. on $S_{1}$ for non-integer $\rho$ and let $s:=T_{\rho} h$. Then

$$
h(\phi)=\frac{1}{2 \rho \sin \pi \rho} \Pi(\phi, s) .
$$

The proof is like in Theorem 3.1.4.4 (*Hadamard).
3.2.4 We will suppose in this item that $V_{t}=I, m=2, \rho$ is integer.

Theorem 3.2.4.1 (Condition on $s$ ) Let $\rho$ be integer, $h$ be $\rho$-t.c. and $T_{\rho} h=s$. Then

$$
\begin{equation*}
\int_{0}^{2 \pi} e^{i \rho \phi} s(d \phi)=0 \tag{3.2.4.1}
\end{equation*}
$$

Proof. We have for $f \in \mathcal{D}\left(S_{1}\right)$ :

$$
\langle s, f\rangle=\left\langle T_{\rho} h, f\right\rangle=\left\langle h, T_{\rho} f\right\rangle
$$

Since $e^{i \rho \phi} \in \mathcal{D}\left(S_{1}\right)$ for integer $\rho$ and $T_{\rho} e^{i \rho \phi}=0$, we have for $f:=e^{i \rho \phi}$,

$$
\left\langle s, e^{i \rho \bullet}\right\rangle=\left\langle h, T_{\rho} e^{i \rho \bullet}\right\rangle=0
$$

Let us denote the periodic continuation of the function $f(\phi):=\phi$ from the interval $[0,2 \pi)$ to $(-\infty, \infty)$ as $\tilde{\phi}$.

Theorem 3.2.4.2 (Generalized Fundamental Solution for $T_{\rho}$ ) One has

$$
T_{\rho}\left[-\frac{1}{2 \pi \rho} \tilde{\phi} \sin \rho \phi\right]=\delta(\phi)-\frac{1}{\pi} \cos \rho \phi
$$

in $\mathcal{D}^{\prime}\left(S_{1}\right)$.
Proof. Let $\phi \in(\epsilon, 2 \pi-\epsilon)$. Then

$$
T_{\rho} \tilde{\phi} \sin \rho \phi=2 \rho \cos \rho \phi
$$

because $\tilde{\phi}=\phi$ when $\phi \in(\epsilon, 2 \pi-\epsilon)$. We have also

$$
(\phi \sin \rho \phi)^{\prime}=\sin \rho \phi+\phi \rho \cos \rho \phi
$$

Thus

$$
\left\langle T_{\rho} \tilde{\bullet} \sin \rho \bullet, f\right\rangle=\int_{0}^{2 \pi} \tilde{\phi} \sin \rho \phi T_{\rho} f d \phi=\lim _{\epsilon \rightarrow 0} \int_{\epsilon}^{2 \pi-\epsilon} \phi \sin \rho \phi T_{\rho} f d \phi
$$

Integrating by parts we obtain

$$
\begin{aligned}
\int_{\epsilon}^{2 \pi-\epsilon} \phi \sin \rho \phi T_{\rho} f d \phi=\left.\phi \sin \rho \phi f^{\prime}(\phi)\right|_{\epsilon} ^{2 \pi-\epsilon} & -f(\phi)\left[\sin \rho \phi+\left.\phi \rho \cos \rho \phi\right|_{\epsilon} ^{2 \pi-\epsilon}\right. \\
& +\int_{\epsilon}^{2 \pi-\epsilon} T_{\rho}[\phi \sin \rho \phi] f(\phi) d \phi
\end{aligned}
$$

Passing to the limit as $\epsilon \rightarrow 0$ and taking in account that $f$ is periodic and continuous we obtain

$$
\begin{aligned}
\left\langle T_{\rho}[\tilde{\bullet} \sin \rho \bullet], f\right\rangle & =-2 \pi \rho f(0)+2 \rho \int_{0}^{2 \pi} \cos \rho \phi f(\phi) d \phi \\
& =-2 \pi \rho f(0)+2 \rho\langle\cos \rho \bullet, f\rangle
\end{aligned}
$$

Set

$$
\hat{\Pi}(\phi, s):=\int_{0}^{2 \pi} \widetilde{(\phi-\psi)} \sin \rho(\phi-\psi) s(d \psi)
$$

Theorem 3.2.4.3 One has

$$
T_{\rho} \hat{\Pi}(\bullet, s)=-2 \pi \rho s \text { in } \mathcal{D}^{\prime}\left(S_{1}\right)
$$

for $s$ that satisfies (3.2.4.1).
Proof. Using Theorem 3.2.4.2 we obtain

$$
\left\langle T_{\rho} \hat{\Pi}(\bullet, s), f\right\rangle=\langle s, f\rangle-\frac{1}{\pi}\left\langle\int_{0}^{2 \pi} \cos \rho(\bullet-\psi) s(d \psi), f\right\rangle
$$

The last integral is zero because of Theorem 3.2.4.1.
Theorem 3.2.4.4 (Representation of $\rho$-t.c.f. for an integer $\rho$ ) Let $h$ be a $\rho$-t.c.f.for an integer $\rho$ and $T_{\rho} h:=s$. Then

$$
h(\phi)=\Re c e^{i \phi}+\hat{\Pi}(\phi, s)
$$

for some complex constant c.
Proof. The function $H(\phi):=h(\phi)-\hat{\Pi}(\phi, s)$ satisfies the equation $T_{\rho} H=0$ in $\mathcal{D}^{\prime}\left(S_{1}\right)$ because of Theorem 3.2.4.3 and it is real. Thus $H(\phi)=\Re c e^{i \phi}$.
3.2.5 The class $T C_{\rho}$ of $\rho$-t.c.functions has a number of properties of subharmonic functions.

The function $\widetilde{\cos \rho} \phi$ is continuous and $\widetilde{\phi} \sin \rho \phi$ is continuous for integer $\rho$. Therefore any $\rho$-t.c.f is continuous as follows from Theorem 3.2.3.3 and 3.2.4.4.

Set

$$
\mathcal{E}(\phi):=\frac{1}{2 \rho} \sin \rho|\phi| .
$$

For any interval $I:=(\alpha, \beta) \Subset(-\pi, \pi)$ this function satisfies the equality

$$
T_{\rho} \mathcal{E}=\delta
$$

in $\mathcal{D}^{\prime}(\alpha, \beta)$, where $\delta$ is the Dirac function in zero.
Let $G_{I}(\psi, \phi)$ be the Green function of $T_{\rho}$ for the interval $I$. By definition it must be symmetric with respect to $\phi, \psi$ and have the form

$$
\begin{equation*}
G_{I}(\phi, \psi):=\frac{1}{2 \rho} \sin \rho|\phi-\psi|+A_{I} \cos \rho \phi \cos \rho \psi+B_{I} \sin \rho \phi \sin \rho \psi \tag{3.2.5.1}
\end{equation*}
$$

where $A_{I}, B_{I}$ are chosen such that $G_{I}(\phi, \psi)$ is equal to zero on $\partial\{I \times I\}$. An explicit form of $G_{I}$ is given by

$$
G_{I}(\phi, \psi)= \begin{cases}\frac{\sin \rho(\beta-\phi) \sin \rho(\psi-\alpha)}{\rho \sin \rho(\beta-\alpha)}, & \text { for } \psi<\phi \\ \frac{\sin \rho(\beta-\psi) \sin \rho(\phi-\alpha)}{\rho \sin \rho(\beta-\alpha)}, & \text { for } \phi<\psi\end{cases}
$$

The following assertion is analogous to the Riesz theorem (Theorem 2.6.4.3):
Theorem 3.2.5.1 (Representation on $I$ ) Let $h \in T C_{\rho}$ and let $I$ be an interval of length mes $I<\pi / \rho$. Then

$$
h(\phi)=Y_{\rho}(\phi, h)-\int_{\alpha}^{\beta} G_{I}(\phi, \psi) s(d \psi),
$$

where $Y_{\rho}(\phi, h)$ is the only solution of the boundary problem:

$$
\begin{equation*}
T_{\rho} Y=0, Y(\alpha)=h(\alpha), Y(\beta)=h(\beta) \tag{3.2.5.2}
\end{equation*}
$$

and $s:=T_{\rho} h$.
Proof. Set

$$
\Pi_{I}(\phi, s):=\int_{\alpha}^{\beta} G_{I}(\phi, \psi) s(d \psi)
$$

One can check as in Theorem 3.2.3.2 that $T_{\rho} \Pi_{I}=-s$ in $\mathcal{D}^{\prime}(I)$. Then the function

$$
Y_{\rho}(\phi):=h(\phi)+\Pi_{I}(\phi, s)
$$

satisfies the conditions (3.2.5.2).
The explicit form of $Y_{\rho}(\phi)$ is

$$
\begin{equation*}
Y_{\rho}(\phi)=\frac{h(\alpha) \sin \rho(\beta-\phi)+h(\beta) \sin \rho(\phi-\alpha)}{\sin \rho(\beta-\phi)} . \tag{3.2.5.3}
\end{equation*}
$$

Since $\Pi_{I}(\phi) \geq 0$ we have
Theorem 3.2.5.2 ( $\rho$-Trigonometric Majorant) Suppose $h \in T C_{\rho}$ and $Y_{\rho}(\phi)$ is the solution of (3.2.5.2). Then

$$
h(\phi) \leq Y_{\rho}(\phi), \phi \in I
$$

if $\beta-\alpha<\pi / \rho$.
This inequality can be written in the symmetric form

$$
\begin{equation*}
h(\alpha) \sin \rho(\beta-\phi)+h(\phi) \sin \rho(\alpha-\beta)+h(\beta) \sin \rho(\phi-\alpha) \geq 0 \tag{3.2.5.4}
\end{equation*}
$$

for $\max (\alpha, \phi, \beta)-\min (\alpha, \phi, \beta)<\pi / \rho$. It is called the fundamental relation of indicator.

Theorem 3.2.5.3 (Subharmonicity and $\rho$-t.c.) A function $h(\phi) \in T C_{\rho}$ iff the function $u\left(r e^{i \phi}\right):=h(\phi) r^{\rho}$ is subharmonic in $\mathbb{R}^{2}$.

Proof. Sufficiency follows from Theorem 3.2.2.2. Let us prove necessity. The function $u_{1}(z):=r^{\rho} \sin \rho|\phi|$ is subharmonic. Actually, it is harmonic for $\phi \neq 0, r \neq 0$ and can be represented in the form

$$
u_{1}(z)=\max \left(r^{\rho} \sin \phi,-r^{\rho} \sin \phi\right)
$$

in a neighborhood of the line $\phi=0$. Hence, it is subharmonic because of sh2), Theorem 2.6.1.1 (Elementary Properties).

The function $u_{2}(z):=r^{\rho} \Pi_{I}(\phi, s)$ is subharmonic because of sh5) and sh4), Theorem 2.6.1.1. The function $r^{\rho} Y_{\rho}(\phi)$ is harmonic for $r>0$. This can be checked directly. Hence, $u(z)$ is subharmonic for $r>0$ because of Theorem 3.2.5.1. By Theorem 2.6.2.2 $u(z)$ is also subharmonic for $r=0$, because it is, obviously, continuous at $z=0$.

Theorem 3.2.5.4 (Elementary Properties of $\rho$-t.c.Functions) One has
tc1) If $h \in T C_{\rho}$, then $A h \in T C_{\rho}$ for $A>0$;
$\mathrm{tc} 2)$ If $h_{1}, h_{2} \in T C_{\rho}$, then $h_{1}+h_{2}, \max \left(h_{1}, h_{2}\right) \in T C_{\rho}$.
These properties follow from Theorem 3.2.5.3 and properties of subharmonic functions.

Exercise 3.2.5.1 Prove Theorem 3.2.5.4.
Similarly to (usual) convexity, $\rho$-t.convexity of functions implies several analytic properties.

Theorem 3.2.5.5 Let $h \in T C_{\rho}$; then there exist right ( $h_{+}^{\prime}$ ) and left ( $h_{-}^{\prime}$ ) derivatives and they coincide everywhere except, maybe, for a countable set of points.

Proof. It is enough to prove these properties for the potential

$$
\Pi(\phi):=\int_{\alpha}^{\beta} \sin \rho|\phi-\psi| s(d \psi)
$$

because of (3.2.5.1) and Theorem 3.2.5.1.
We will prove that

$$
\begin{align*}
& \Pi_{+}^{\prime}(\phi)=\rho \int_{\alpha}^{\phi-0} \cos \rho(\phi-\psi) s(d \psi)+\rho \mu(\phi)-\rho \int_{\phi+0}^{\beta} \cos \rho(\phi-\psi) s(d \psi)  \tag{3.2.5.5}\\
& \Pi_{-}^{\prime}(\phi)=\rho \int_{\alpha}^{\phi-0} \cos \rho(\phi-\psi) s(d \psi)-\rho \mu(\phi)-\rho \int_{\phi+0}^{\beta} \cos \rho(\phi-\psi) s(d \psi) \tag{3.2.5.6}
\end{align*}
$$

where $\mu(\phi)$ is the measure, concentrated in the point $\phi$.
We have for $\Delta>0$ :

$$
\begin{aligned}
\frac{\Pi(\phi+\Delta)-\Pi(\phi)}{\Delta}= & \int_{\alpha}^{\phi-0} \\
& \frac{\sin \rho|\phi+\Delta-\psi|-\sin \rho|\phi-\psi|}{\Delta} s(d \psi) \\
& +\frac{\sin \rho \Delta}{\Delta} \mu(\phi)+\int_{\phi+0}^{\phi+\Delta} \cdots+\int_{\phi+\Delta}^{\beta} \cdots
\end{aligned}
$$

Let us estimate the second integral.We have
$\int_{\phi+0}^{\phi+\Delta}\left|\frac{\sin \rho|\phi+\Delta-\psi|-\sin \rho|\phi-\psi|}{\Delta}\right| s(d \psi) \leq \frac{2 \sin \rho \Delta}{\Delta}[s(\phi+\Delta)-s(\phi+0)]=o(1)$ when $\Delta \rightarrow+0$.

Passing to the limit, we obtain (3.2.5.5). The equality (3.2.5.6) is obtained in the same way when $\Delta<0$.

Since $\mu(\phi) \neq 0$ at most in a countable set, for all the other points $\Pi_{+}^{\prime}(\phi)=$ $\Pi_{-}^{\prime}(\phi)$.
3.2.6 Now we consider the case $m \geq 3$. We will obtain for the $\rho$-subspherical function a representation like for the $\rho$-trigonometrically convex functions.

Theorem 3.2.6.1 (Subharmonicity and Subsphericality) Let $h$ be subspherical in a neighborhood of $y \in S_{1}$. Then the function $u(x):=h(y) r^{\rho}, x=r y$ is subharmonic in the corresponding neighborhood of the ray $x=r y: 0<r<\infty$.

Proof. Let $f \in \mathcal{D}^{\prime}\left(\mathbb{R}^{m} \backslash 0\right)$. We can represent it in the form $f:=f(r x), x \in S_{1}$ where $f(\bullet x) \in \mathcal{D}^{\prime}(0, \infty)$ for any $x$.

Then

$$
\begin{aligned}
\langle u, \Delta f\rangle= & \int_{0}^{\infty} \int_{S_{1}} u(r x) \Delta f(r x) r^{m-1} d r d x \\
= & \int_{0}^{\infty} \int_{S_{1}} u(r x) \frac{1}{r^{m-1}} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r} f(r x) r^{m-1} d r d x \\
& +\int_{0}^{\infty} \int_{S_{1}} u(r x) \frac{1}{r^{2}} \Delta_{x^{0}} f(r x) r^{m-1} d r d x
\end{aligned}
$$

Integrating by parts in the first integral, we obtain

$$
\int_{0}^{\infty} \rho(\rho+m-2) r^{\rho+m-3} \int_{S_{1}} h(x) f(r x) d r d x
$$

Set

$$
\begin{equation*}
\mathcal{S}_{\rho}:=\Delta_{\boldsymbol{x}^{0}}+\rho(\rho+m-2) . \tag{3.2.6.1}
\end{equation*}
$$

Together with the second summand we obtain

$$
\langle u, f\rangle=\int_{0}^{\infty}\left[\int_{S_{1}} h(x) \mathcal{S}_{\rho} f(r x) d x\right] r^{\rho+m-3} d r>0
$$

if $f(r x) \geq 0$.

Note that the Riesz measure for such $u$ has the form

$$
\mu\left(r^{m-1} d r d x\right)=r^{\rho+m-3} d r \nu_{h}(d x)
$$

where $\nu_{h}$ is a positive measure on $S_{1}$, that is equal to $\mathcal{S}_{\rho} h$ in $\mathcal{D}^{\prime}\left(S_{1}\right)$.
For a non-integer $\rho$ set

$$
\mathcal{E}_{\rho}(x, y):=\int_{0}^{\infty} G_{p}(x, r y, m) r^{\rho+m-3} d r
$$

where $G_{p}$ is the primary kernel and $x, y \in S_{1}$.
Theorem 3.2.6.2 For non-integer $\rho$ and any $\rho$-subspherical function $h$ one has

$$
h(x)=\int_{S_{1}} \mathcal{E}_{\rho}(x, y) \nu_{h}(d y)
$$

Proof. Set in (3.1.4.6) $v:=r^{\rho} h(x)$. It is clear that $v \in U[\rho]$. We have

$$
r^{\rho} h(x)=\int_{S_{1}} \int_{0}^{\infty} G_{p}(r x, t y, m) t^{\rho+m-3} d t \nu_{h}(d x)
$$

Now we make the change $t^{\prime}:=t / r$ and use the homogeneity of $G_{p}(r x, t y, m)$.
Exercise 3.2.6.1 Show that $\mathcal{E}_{\rho}(x, y)$ is a fundamental solution of the operator $\mathcal{S}_{\rho}$.
For an integer $\rho=p$ set

$$
\mathcal{E}_{\rho}^{\prime}(x, y):=\int_{0}^{1} G_{p-1}(x, r y) r^{\rho+m-3} d r+\int_{1}^{\infty} G_{p}(x, r y) r^{\rho+m-3} d r
$$

Exercise 3.2.6.2 Prove the next
Theorem 3.2.6.3 For any integer $\rho=p$ and any $\rho$-subspherical function hone has

$$
h(x)=Y_{p}(x)+\int_{S_{1}} \mathcal{E}_{\rho}^{\prime}(x, y) \nu_{h}(d y)
$$

where $Y_{p}$ is some $p$-spherical function.
For any p-spherical function $Y$,

$$
\int_{S_{1}} Y(x) \nu_{h}(d y)=0
$$

3.2.7 We return to the general case when $x \in \mathbb{R}^{m}, V_{t}$ is a one-parametric group, $\rho(r)$ is a proximate order and $u \in S H(\rho(r))$. The following theorem represents indicators in a form of limits in the usual topology.

Theorem 3.2.7.1 (Classic Indicators) One has

$$
\begin{equation*}
h(x, u)=\sup _{T}\left[\limsup _{t_{j} \rightarrow \infty} u_{t_{j}}\right]^{*}(x)=\left[\limsup _{t \rightarrow \infty} u_{t}(x)\right]^{*} \tag{3.2.7.1}
\end{equation*}
$$

where $*$ can be deleted outside a set of zero capacity, and

$$
\begin{equation*}
\underline{h}(x, u)=\inf _{T}\left[\limsup _{t_{j} \rightarrow \infty} u_{t_{j}}\right]^{*}(x), \tag{3.2.7.2}
\end{equation*}
$$

where $T$ is the set of all the sequences that tend to infinity.
Proof. Let us prove (3.2.7.1). Set

$$
\begin{equation*}
h\left(x, u,\left\{t_{j}\right\}\right):=\limsup _{t_{j} \rightarrow \infty} u_{t_{j}}(x) . \tag{3.2.7.3}
\end{equation*}
$$

Let $v \in \operatorname{Fr}[u]$ and $u_{t_{j}} \rightarrow v$ in $\mathcal{D}^{\prime}$. Then

$$
\begin{equation*}
h^{*}\left(x, u,\left\{t_{j}\right\}\right)=v(x) \tag{3.2.7.4}
\end{equation*}
$$

by Theorem 2.7.3.3. (H. Cartan+). Thus

$$
\begin{equation*}
\sup _{T} h^{*}\left(x, u,\left\{t_{j}\right\}\right) \geq h(x, u) . \tag{3.2.7.5}
\end{equation*}
$$

Let $\epsilon>0$ be arbitrarily small, and $t_{j}:=t_{j}(x)$ be a sequence such that

$$
h^{*}\left(x, u,\left\{t_{j}\right\}\right) \geq \sup _{T} h^{*}\left(x, u,\left\{t_{j}\right\}\right)-\epsilon .
$$

We can find a subsequence $\left\{t_{j}\right\}$ (we keep the same notation for it) and $v \in \mathbf{F r}[u]$ such that $u_{t_{j}} \rightarrow v$ in $\mathcal{D}^{\prime}$. From (3.2.7.4) we obtain

$$
h(x, u) \geq v(x) \geq \sup _{T} h^{*}\left(x, u,\left\{t_{j}\right\}\right)-\epsilon .
$$

Thus the reverse inequality to (3.2.7.5) holds. Therefore

$$
h(x, u)=\sup _{T} h^{*}\left(x, u,\left\{t_{j}\right\}\right) .
$$

Let us prove the second equality in (3.2.7.1). Since

$$
\sup _{T} h\left(x, u,\left\{t_{j}\right\}\right)=\limsup _{t \rightarrow \infty} u_{t}(x)
$$

we have

$$
\begin{equation*}
h(x, u) \geq\left[\limsup _{t \rightarrow \infty} u_{t}\right]^{*}(x) \tag{3.2.7.6}
\end{equation*}
$$

Let us prove the opposite inequality. Let $v \in \operatorname{Fr}[u]$. There exists a sequence $t_{j} \rightarrow \infty$ such that $u_{t_{j}} \rightarrow v$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$. By (3.2.7.4)

$$
\left[\limsup _{t \rightarrow \infty} u_{t}\right]^{*} \geq h^{*}\left(x, u,\left\{t_{j}\right\}\right)=v(x)
$$

Since it holds for every $v \in \operatorname{Fr}[u]$ we have the reverse inequality to (3.2.7.6). Hence, (3.2.7.1) is proved completely.

Let us prove (3.2.7.2). From (3.2.7.4) we have

$$
\inf _{T} h^{*}\left(x, u,\left\{t_{j}\right\}\right) \leq v(x)
$$

for all $v \in \operatorname{Fr}[u]$. Therefore

$$
\begin{equation*}
\inf _{T} h^{*}\left(x, u,\left\{t_{j}\right\}\right) \leq \underline{h}(x, u) . \tag{3.2.7.7}
\end{equation*}
$$

Let us prove the opposite inequality. Let $\left\{t_{j}\right\}$ be any sequence that tends to $\infty$. Let us find a subsequence $\left\{t_{j^{\prime}}\right\}$ such that $u_{t_{j^{\prime}}} \rightarrow v$ in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$. Then

$$
h\left(x, u,\left\{t_{j}\right\}\right) \geq \limsup _{j^{\prime} \rightarrow \infty} u_{t_{j^{\prime}}}(x)
$$

Taking $*$ from the two sides of this inequality and using Theorem 2.7.3.3, we obtain

$$
h^{*}\left(x, u,\left\{t_{j}\right\}\right) \geq\left[\limsup _{j^{\prime} \rightarrow \infty} u_{t_{j^{\prime}}}\right]^{*}(x)=v(x) \geq \underline{h}(x, u) .
$$

This implies the reverse inequality to (3.2.7.7). Hence (3.2.7.2) holds.
Corollary 3.2.7.2 If all the functions (3.2.7.3) are upper semicontinuous, then

$$
h(x, u)=\limsup _{t \rightarrow \infty} u_{t}(x), \quad \underline{h}(x, u)=\liminf _{t \rightarrow \infty} u_{t}(x) .
$$

Proof. We have $h^{*}\left(x, u,\left\{t_{j}\right\}\right)=h\left(x, u,\left\{t_{j}\right\}\right)$ and thus

$$
\begin{aligned}
& h(x, u)=\sup _{T}\left[\limsup _{t_{j} \rightarrow \infty} u_{t_{j}}\right](x)=\liminf _{t \rightarrow \infty} u_{t}(x), \\
& \underline{h}(x, u)=\inf _{T}\left[\limsup _{t_{j} \rightarrow \infty} u_{t_{j}}\right](x)=\liminf _{t \rightarrow \infty} u_{t}(x) .
\end{aligned}
$$

Theorem 3.2.7.3 (Indicators of Harmonic Function) Let $u \in S H(\rho(r))$ be harmonic for all the large $|y|$ in a "cone" of the form

$$
C o_{\Omega}:=\left\{y=P_{t} x: x \in \Omega, t \in(0 ; \infty)\right\}
$$

where $\Omega \subset S_{1}$. Then

$$
\begin{equation*}
h(x, u)=\limsup _{t \rightarrow \infty} u_{t}(x) \tag{3.2.7.8}
\end{equation*}
$$

and

$$
\begin{equation*}
\underline{h}(x, u)=\liminf _{t \rightarrow \infty} u_{t}(x) \tag{3.2.7.9}
\end{equation*}
$$

for $x \in C o_{\Omega}$.

Proof. The harmonicity of $u$ in $C o_{\Omega}$ implies $\left[u_{t}\right]_{\epsilon}(x)=u_{t}(x)$ for large $t$ and sufficiently small $\epsilon$ when $x \in C o_{\Omega}$.

The family $\left[u_{t}\right]_{\epsilon}$ is uniformly continuous by reg3), Theorem 2.3.4.5 (Properties of Regularizations). Thus the function (3.2.7.5) is continuous. Therefore we can use Corollary 3.2.7.2.

Theorem 3.2.7.4 (Indicator for $m=2$ ) Let $u \in S H\left(\mathbb{R}^{2}\right)$. Then

$$
\begin{equation*}
h(z, u)=\limsup _{t \rightarrow \infty} u_{t}(x) \tag{3.2.7.10}
\end{equation*}
$$

I.e., the star in (3.2.7.1) can be deleted.

Proof. Let as denote as $h_{1}(z, u)$ the right part of (3.2.7.10). The "homogeneity" of the indicator (3.2.1.10) and also of $h_{1}(z, u)$ implies the following property: if the inequality $h_{1}(z, u)<h(z . u)$ holds for some $z_{0}$, it holds on the whole orbit

$$
z=\left\{P_{t} z_{0}: 0<t<\infty\right\}
$$

that has a positive capacity in $\mathbb{R}^{2}$. This contradicts Theorem 3.2.7.1.

### 3.3 Densities

3.3.1 In the sequel $G$ is an open set, $K$ is a compact set and $E$ a bounded Borel set. Let $\mu \in \mathcal{M}(\rho(r))$ and $\operatorname{Fr}[\mu]:=\operatorname{Fr}\left[\mu, \rho(r), V_{t}, \mathbb{R}^{m}\right]$ be the limit set of $\mu$. Set

$$
\begin{aligned}
\bar{\Delta}(G, \mu) & :=\sup \{\nu(G): \nu \in \mathbf{F r}[\mu]\} ; \\
\bar{\Delta}(E, \mu) & :=\inf \{\bar{\Delta}(G, \mu): G \supset E\} ; \\
\underline{\Delta}(K, \mu) & :=\inf \{\nu(K): \nu \in \mathbf{F r}[\mu]\} ; \\
\underline{\Delta}(E, \mu) & :=\sup \{\underline{\Delta}(K): K \subset E\}
\end{aligned}
$$

The quality $\bar{\Delta}(E, \mu),(\underline{\Delta}(E, \mu))$ is called the upper (lower) density of $\mu$ relative to the proximate order $\rho(r)$ and the family $V_{t}$.

Theorem 3.3.1.1 (Properties of Densities) The following properties hold:
dens1) if $E=\varnothing$, then $\bar{\Delta}(E, \bullet)=\underline{\Delta}(E, \bullet)=0$;
dens2) $\forall E, \underline{\Delta}(E, \bullet) \leq \bar{\Delta}(E, \bullet)$;
dens3) monotonicity: $\underline{\Delta}, \bar{\Delta}\left(E_{1}, \bullet\right) \leq \underline{\Delta}, \bar{\Delta}\left(E_{2}, \bullet\right)$ for $E_{1} \subset E_{2}$;
dens4) generalized semi-additivity ${ }^{1}$ with respect to a set:

$$
\begin{align*}
& \bar{\Delta}\left(E_{1} \cup E_{2}, \bullet\right)+\underline{\Delta}\left(E_{1} \cap E_{2}, \bullet\right) \leq \bar{\Delta}\left(E_{1}, \bullet\right)+\bar{\Delta}\left(E_{2}, \bullet\right)  \tag{3.3.1.1}\\
& \underline{\Delta}\left(E_{1} \cup E_{2}, \bullet\right)+\bar{\Delta}\left(E_{1} \cap E_{2}, \bullet\right) \geq \underline{\Delta}\left(E_{1}, \bullet\right)+\underline{\Delta}\left(E_{2}, \bullet\right) \tag{3.3.1.2}
\end{align*}
$$

[^2]dens5) continuity from the right and from the left.
\[

$$
\begin{align*}
& E_{n} \uparrow E \Longrightarrow \bar{\Delta}\left(E_{n}, \bullet\right) \uparrow \bar{\Delta}(E, \bullet) ; K_{n} \downarrow K \Longrightarrow \bar{\Delta}\left(K_{n}, \bullet\right) \downarrow \bar{\Delta}(K, \bullet),  \tag{3.3.1.3}\\
& E_{n} \downarrow E \Longrightarrow \underline{\Delta}\left(E_{n}, \bullet\right) \downarrow \underline{\Delta}(E, \bullet) ; G_{n} \uparrow G \Longrightarrow \underline{\Delta}\left(G_{n}, \bullet\right) \uparrow \underline{\Delta}(G, \bullet) \tag{3.3.1.4}
\end{align*}
$$
\]

dens6) semi-additivity and positive homogeneity with respect to $\mu$, i.e.,

$$
\begin{align*}
\bar{\Delta}\left(E, \mu_{1}+\mu_{2}\right) & \leq \bar{\Delta}\left(E, \mu_{1}\right)+\bar{\Delta}\left(E, \mu_{2}\right) ;  \tag{3.3.1.5}\\
\underline{\Delta}\left(E, \mu_{1}+\mu_{2}\right) & \geq \underline{\Delta}\left(E, \mu_{1}\right)+\underline{\Delta}\left(E, \mu_{2}\right) ;  \tag{3.3.1.6}\\
\bar{\Delta}, \underline{\Delta}(E, \lambda \mu) & =\lambda \bar{\Delta}, \lambda \underline{\Delta}(E, \mu) \tag{3.3.1.7}
\end{align*}
$$

for $\lambda \geq 0$;
dens7) invariance with respect to the map $(\bullet)_{[t]}$ (see, 3.1.2.4a), i.e.,

$$
t^{-\rho-m+2} \bar{\Delta}, \underline{\Delta}\left(P_{t} E, \bullet\right)=\bar{\Delta}, \underline{\Delta}(E, \bullet)
$$

Proof of Theorem 3.3.1.1. The property dens1) holds because the empty set is open by definition. The properties dens2) and dens3) hold because of the monotonicity of $\nu$.

Let us prove dens4). Since $\nu$ is a measure we have

$$
\nu\left(G_{1} \cup G_{2}, \mu\right)+\nu\left(G_{1} \cap G_{2}, \mu\right)=\nu\left(G_{1}, \mu\right)+\nu\left(G_{2}, \mu\right)
$$

for any $G_{1} \supset E_{1}$ and $G_{2} \supset E_{2}$.
From this we obtain

$$
\begin{equation*}
\nu\left(G_{1} \cup G_{2}, \mu\right)+\nu\left(K_{1} \cap K_{2}, \mu\right) \leq \nu\left(G_{1}, \mu\right)+\nu\left(G_{2}, \mu\right) \tag{3.3.1.8}
\end{equation*}
$$

for $K_{1} \subset E_{1}$ and $K_{2} \subset E_{2}$.
The right side of (3.3.1.8) is no larger than $\bar{\Delta}\left(G_{1}, \bullet\right)+\bar{\Delta}\left(G_{1}, \bullet\right)$. Now we can take supremum over $\nu \in \mathrm{Fr} \mu$ in the first summand of the left side and infimum in the second summand. Thus we obtain

$$
\begin{equation*}
\bar{\Delta}\left(G_{1} \cup G_{2}, \bullet\right)+\underline{\Delta}\left(K_{1} \cap K_{2}, \bullet\right) \leq \bar{\Delta}\left(G_{1}, \bullet\right)+\bar{\Delta}\left(G_{1}, \bullet\right) \tag{3.3.1.9}
\end{equation*}
$$

Since $\bar{\Delta}(E, \bullet)$ and $\underline{\Delta}(E, \bullet)$ are monotonic with respect to $E$,

$$
\inf \left\{\bar{\Delta}\left(G_{1} \cup G_{2}, \bullet\right): G_{1} \supset E_{1}, G_{2} \supset E_{2}\right\}=\bar{\Delta}\left(E_{1} \cup E_{2}, \bullet\right)
$$

and

$$
\sup \left\{\underline{\Delta}\left(K_{1} \cap K_{2}, \bullet\right): K_{1} \subset E_{1}, K_{2} \subset E_{2}\right\}=\underline{\Delta}\left(E_{1} \cup E_{2}\right)
$$

Thus we obtain the first inequality in dens4) from (3.3.1.9). The second one can be obtained analogously ${ }^{2}$.

[^3]Let us prove dens5). For arbitrary $G \supset K$ there exists $n_{0}$ such that $K_{n} \subset G$ for $n>n_{0}$. According to dens3),

$$
\bar{\Delta}(K, \bullet) \leq \bar{\Delta}\left(K_{n}, \bullet\right) \leq \bar{\Delta}(G, \bullet)
$$

Hence,

$$
\bar{\Delta}(K, \bullet) \leq \lim _{n \rightarrow \infty} \bar{\Delta}\left(K_{n}, \bullet\right) \leq \bar{\Delta}(G, \bullet)
$$

Taking infimum over all $G \supset K$, we obtain the second assertion in (3.3.1.3).
For $G_{n} \uparrow G$ we have the equality

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \bar{\Delta}\left(G_{n}, \bullet\right)=\sup _{n} \bar{\Delta}\left(G_{n}, \bullet\right)=\bar{\Delta}(G, \bullet) \tag{3.3.1.10}
\end{equation*}
$$

because one can change the order of taking the supremum on $n$ and on $\nu \in \operatorname{Fr}[\mu]$.
Let $E_{n} \uparrow E$ and let $\epsilon$ be arbitrarily small. One can find $G_{n} \supset E_{n}$ such that

$$
\bar{\Delta}\left(G_{n}, \bullet\right) \leq \bar{\Delta}\left(E_{n}, \bullet\right)+\epsilon
$$

Since $G:=\bigcup_{1}^{\infty} G_{n} \supset E$ we have

$$
\bar{\Delta}\left(G_{n}, \bullet\right)-\epsilon \leq \bar{\Delta}\left(E_{n}, \bullet\right) \leq \bar{\Delta}(E, \bullet) \leq \bar{\Delta}(G, \bullet)
$$

Using (3.3.1.10), we obtain

$$
\bar{\Delta}(E, \bullet)-\lim _{n \rightarrow \infty} \bar{\Delta}\left(E_{n}, \bullet\right) \leq \epsilon
$$

Since $\epsilon$ is arbitrarily small,

$$
\bar{\Delta}(E, \bullet) \leq \lim _{n \rightarrow \infty} \bar{\Delta}\left(E_{n}, \bullet\right)
$$

and hence the first assertion in (3.3.1.3) holds.
The assertion (3.3.1.4) can be proved analogously. ${ }^{3}$
Let us prove dens6). One has

$$
\bar{\Delta}\left(G, \mu_{1}+\mu_{2}\right)=\sup \left\{\nu(G): \nu \in \mathbf{F r}\left[\mu_{1}+\mu_{2}\right]\right\}
$$

Since

$$
\operatorname{Fr}\left[\mu_{1}+\mu_{2}\right] \subset \mathbf{F r}\left[\mu_{1}\right]+\mathbf{F r}\left[\mu_{2}\right]
$$

(see frmu1), Theorem 3.1.3.4 (Properties of $\mu \mapsto \operatorname{Fr}[\mu]$ )) one can continue the previous equality as

$$
\begin{aligned}
& \leq \sup \left\{\nu(G): \nu \in \mathbf{F r}\left[\mu_{1}\right]+\mathbf{F r}\left[\mu_{2}\right]\right\} \\
& =\sup \left\{\nu(G): \nu \in \mathbf{F r}\left[\mu_{1}\right]\right\}+\sup \left\{\nu(G): \nu \in \mathbf{F r}\left[\mu_{2}\right]\right\}=\bar{\Delta}\left(G, \mu_{1}\right)+\bar{\Delta}\left(G, \mu_{2}\right)
\end{aligned}
$$

Passing to the infimum over $G \supset E$, we obtain (3.3.1.5). The assertions (3.3.1.6) and (3.3.1.7) can be proved analogously. ${ }^{4}$

The properties dens7) follow from the invariance of $\mathbf{F r}[\mu]$ (see frm3), Theorem 3.1.3.3. (Properties of $\operatorname{Fr}[\mu])$ ).

[^4]Exercise 3.3.1.1 Prove the subadditivity of $\bar{\Delta}(E, \bullet)$ :

$$
\bar{\Delta}\left(E_{1} \cup E_{2}, \bullet\right) \leq \bar{\Delta}\left(E_{1}, \bullet\right)+\bar{\Delta}\left(E_{2}, \bullet\right)
$$

and the superadditivity of $\underline{\Delta}(E, \bullet)$ :

$$
\underline{\Delta}\left(E_{1} \cup E_{2}, \bullet\right) \geq \underline{\Delta}\left(E_{1}, \bullet\right)+\underline{\Delta}\left(E_{2}, \bullet\right)
$$

from Theorem 3.3.1.1.
Exercise 3.3.1.2 Prove (3.3.1.2).
Exercise 3.3.1.3 Prove (3.3.1.4).
Exercise 3.3.1.4 Prove (3.3.1.6) and (3.3.1.7).
Set for $I \subset(0, \infty)$ and $\Omega \subset S_{1}$,

$$
\operatorname{Co}_{\Omega}(I):=\left\{x=P_{t} y: y \in \Omega, t \in I\right\} .
$$

Also set $I_{t}:=(0, t)$.
Theorem 3.3.1.2 (Cone's Densities) One has

$$
\bar{\Delta}, \underline{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{t}\right)\right)=t^{\rho+m-2} \bar{\Delta}, \underline{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{1}\right)\right)
$$

We obtain this from dens7), Theorem 3.3.1.1, taking $E:=\operatorname{Co}_{\Omega}\left(I_{1}\right)$.
Exercise 3.3.1.5 Show that for $m=2, S_{1}=\{|z|=1\}, \Omega=\left\{z=e^{i \phi}: \phi \in(\alpha, \beta)\right\}$ $\mathrm{Co}_{\Omega}\left(I_{t}\right)$ is a sector of radius $t$ corresponding to the $\operatorname{arc}(\alpha, \beta)$ on the unit circle.
3.3.2 Let $\delta(E)$ be a monotonic function of $E \in \mathbb{R}^{m}$. A set $E$ is called $\delta$-squarable if

$$
\begin{equation*}
\sup _{K \subset E} \delta(K)=\inf _{G \supset E} \delta(G) . \tag{3.3.2.1}
\end{equation*}
$$

Example 3.3.2.1 Let $\delta(E)$ be a measure. Then (3.3.2.1) implies $\delta(\partial E)=0$, i.e., $E$ is $\delta$-squarable in the sense of Section 2.2.3.

Exercise 3.3.2.2 Prove the next
Theorem 3.3.2.1 If $\bar{\Delta}(\partial E)=0$, then $E$ is $\bar{\Delta}$-squarable. If $E$ is $\underline{\Delta}$-squarable, then $\Delta(\partial E)=0$.

Set

$$
E_{t}:=\{x: \exists y \in E:|x-y|<t\} .
$$

This is a $t$-extension of $E$.

A family of sets $\mathcal{A}_{1}$ is said to be dense in a family $\mathcal{A}_{2}$ if for each set $E_{2} \in \mathcal{A}_{2}$ and an arbitrarily small $\epsilon>0$ there exists a set $E_{1} \in \mathcal{A}_{1}$ such that

$$
\begin{equation*}
\overline{E_{1} \Delta E_{2}}:=\overline{\left(E_{1} \backslash E_{2}\right) \cup\left(E_{2} \backslash E_{1}\right)} \subset\left(\partial E_{2}\right)_{\epsilon} . \tag{3.3.2.2}
\end{equation*}
$$

## Exercise 3.3.2.3 Prove

Theorem 3.3.2.2 The relation "to be dense in" is reflexive and transitive.
I.e., $\mathcal{A}_{1}$ is dense in $\mathcal{A}_{1}$, and
$\left\{\mathcal{A}_{1}\right.$ is dense in $\left.\mathcal{A}_{2}\right\} \wedge\left\{\mathcal{A}_{2}\right.$ is dense in $\left.\mathcal{A}_{3}\right\} \Longrightarrow\left\{\mathcal{A}_{1}\right.$ is dense in $\left.\mathcal{A}_{3}\right\}$.
There are lots of squarable sets.
Theorem 3.3.2.3 For any monotonic $\delta(E)$ the class of $\delta$-squarable sets is dense in the class of all the subsets of $\mathbb{R}^{m}$.

Proof. For any $E \subset \mathbb{R}^{m}$ set

$$
\begin{equation*}
E(t):=E \cup(\partial E)_{t} . \tag{3.3.2.4}
\end{equation*}
$$

One can check that

$$
\begin{equation*}
\overline{E \Delta E\left(t_{1}\right)} \subset(\partial E)_{t_{2}} \tag{3.3.2.5}
\end{equation*}
$$

and

$$
\begin{equation*}
\overline{E\left(t_{1}\right)} \subset \stackrel{\circ}{E}\left(t_{2}\right) \tag{3.3.2.6}
\end{equation*}
$$

for $t_{1}<t_{2}$.
The function $f(t):=\delta(\stackrel{\circ}{E}(t))$ is monotonic. Hence, its set of continuity points has a concentration point at $t=0$.

Suppose $\epsilon>0$ is arbitrarily small, and $t_{0}<\epsilon$ is a continuity point for $f(t)$. From (3.3.2.6) we have

$$
\lim _{t \rightarrow t_{0}-\epsilon} \delta(\bar{E}(t)) \leq \sup _{K \subset E_{t_{0}}} \delta(K) \leq \inf _{G \supset E_{t_{0}}} \delta(G) \leq \lim _{t \rightarrow t_{0}+\epsilon} \delta(\stackrel{\circ}{E}(t)) .
$$

Hence, $E_{t_{0}}$ is $\delta$-squarable. From (3.3.2.5) we have

$$
\overline{E \Delta E\left(t_{0}\right)} \subset(\partial E)_{\epsilon} .
$$

Set

$$
\bar{\Delta}^{c l}(E)=\limsup _{t \rightarrow \infty} \mu_{t}(E) ; \underline{\Delta}^{c l}(E)=\liminf _{t \rightarrow \infty} \mu_{t}(E)
$$

These are classic densities determined without $\mathcal{D}^{\prime}$-topology. They are monotonic.

The following assertion connects these densities to $\bar{\Delta}$ and $\underline{\Delta}$.
Theorem 3.3.2.4 (Classic Densities) For any $\bar{\Delta}^{c l}$-squarable set $E$,

$$
\begin{equation*}
\bar{\Delta}^{\mathrm{cl}}(E)=\sup \{\nu(E): \nu \in \operatorname{Fr}[\mu]\}=\bar{\Delta}(E, \mu) . \tag{3.3.2.7}
\end{equation*}
$$

For any $\underline{\Delta}^{\mathrm{cl}}$-squarable set $E$,

$$
\begin{equation*}
\underline{\Delta}^{c l}(E)=\inf \{\nu(E): \nu \in \operatorname{Fr}[\mu]\}=\underline{\Delta}(E, \mu) . \tag{3.3.2.7}
\end{equation*}
$$

The theorem follows obviously from the following assertion.
Theorem 3.3.2.5 One has

$$
\begin{align*}
& \sup _{K \subset E} \bar{\Delta}^{\mathrm{cl}}(K) \leq \sup _{\nu \in \mathbf{F r}} \nu(E) \leq \bar{\Delta}(E) \leq \inf _{G \supset E} \bar{\Delta}^{\mathrm{cl}}(G) ;  \tag{3.3.2.8}\\
& \sup _{K \subset E} \underline{\Delta}^{\mathrm{cl}}(K) \leq \inf _{\nu \in \mathbf{F r}} \nu(E) \leq \underline{\Delta}(E) \leq \inf _{G \supset E} \underline{\Delta}^{\mathrm{cl}}(G) . \tag{3.3.2.9}
\end{align*}
$$

Proof. Let us prove, for example, (3.3.2.9). Let us choose any $G$ and $K$ such that $K \subset E \subset G$. We can find a sequence $t_{j} \rightarrow \infty$ such that

$$
\lim _{j \rightarrow \infty} \mu_{t_{j}}(G)=\underline{\Delta}^{\mathrm{cl}}(G) .
$$

Choose a subsequence $t_{j_{n}}$ such that $\mu_{t_{j_{n}}} \rightarrow \nu$ in $\mathcal{D}^{\prime}$ for some $\nu \in \mathrm{Fr}$.
Using Theorems 2.3.4.4. ( $\mathrm{D}^{\prime}$ and $\mathrm{C}^{*}$ ) and 2.2.3.1. ( $\mathrm{C}^{*}$-limits), we obtain

$$
\begin{equation*}
\nu(G) \leq \liminf _{n \rightarrow \infty} \mu_{t_{j_{n}}}(G)=\underline{\Delta}^{\mathrm{cl}}(G) . \tag{3.3.2.10}
\end{equation*}
$$

By the same theorems

$$
\begin{equation*}
\underline{\Delta}^{\mathrm{cl}}(K) \leq \limsup _{n \rightarrow \infty} \mu_{t_{j_{n}}}(K) \leq \nu(K) . \tag{3.3.2.11}
\end{equation*}
$$

From (3.3.2.10) and (3.3.2.11) we obtain

$$
\begin{equation*}
\underline{\Delta}^{\mathrm{cl}}(K) \leq \nu(E) \leq \nu(G) \leq \underline{\Delta}^{\mathrm{cl}}(G) \tag{3.3.2.12}
\end{equation*}
$$

because of monotonicity of $\nu(E)$. Taking supremum over all $K \subset E$ and infimum over all $G \supset E$, we obtain (3.3.2.9).

Exercise 3.3.2.4 Prove (3.3.2.8).
Corollary 3.3.2.6 The following holds:

$$
\begin{align*}
& \bar{\Delta}^{\mathrm{cl}}\left(K_{t}\right)=\bar{\Delta}\left(K_{t}, \mu\right)=t^{\rho+m-2} \bar{\Delta}\left(K_{1}, \mu\right), t \geq 0,  \tag{3.3.2.13}\\
& \underline{\Delta}^{\mathrm{cl}}\left(K_{t}\right)=\underline{\Delta}\left(K_{t}, \mu\right)=t^{\rho+m-2} \underline{\Delta}\left(K_{1}, \mu\right), t \geq 0 . \tag{3.3.2.13'}
\end{align*}
$$

where $K_{t}=\{x:|x|<t\}$ is the ball.
Proof. The right equalities follow from Theorem 3.3.1.2 with $\Omega:=S_{1}$. The left equalities hold at least for one $t$ because of Theorem 3.3.2.4 and hence for all $t$.
3.3.3 Let us note generally speaking that values of $\bar{\Delta}$ and $\underline{\Delta}$ on the sets $\mathrm{Co}_{\Omega}\left(I_{t}\right)$ do not determine their values even on the sets $\operatorname{Con}_{\Omega}(I)$ for $I=\left(t_{1}, t_{2}\right)$. However the following assertion holds.

Theorem 3.3.3.1 (Existence of Density) Let $\Phi$ be a dense ring (see, 2.2.3) on $S_{1}$. Then the conditions

$$
\begin{equation*}
\bar{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{t}\right)\right)=\underline{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{t}\right)\right) \tag{3.3.3.1}
\end{equation*}
$$

for $\Omega \in \Phi$ and some $t$ determine uniquely a measure $\Delta(\Omega)$ on $S_{1} . \operatorname{Fr}[\mu]$ consists of one single measure $\nu$ and

$$
\begin{equation*}
\nu\left(\operatorname{Co}_{\Omega}\left(I_{t}\right)\right)=t^{\rho+m-2} \Delta(\Omega) \tag{3.3.3.2}
\end{equation*}
$$

for all the $t \in(0, \infty)$.
To prove this we need an assertion that is valuable by itself. Set

$$
\begin{equation*}
\bar{\Delta}(\Omega):=\bar{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{1}\right)\right) ; \underline{\Delta}(\Omega):=\underline{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{1}\right)\right) \text { for } \Omega \in S_{1} . \tag{3.3.3.3}
\end{equation*}
$$

We will call them angular densities because for $m=2$ and $V_{t} \equiv I, \Omega$ determines an angle in the plane.

Let $\Omega^{G}$ denote an open set in $S_{1}$ and $\Omega^{K}$ a closed one.
Theorem 3.3.3.2 (Angular Densities) One has

$$
\begin{equation*}
\bar{\Delta}(\Omega)=\inf _{\Omega^{G} \supset \Omega} \bar{\Delta}\left(\Omega^{G}\right) ; \underline{\Delta}(\Omega)=\sup _{\Omega^{K} \subset \Omega} \underline{\Delta}\left(\Omega^{K}\right) \tag{3.3.3.4}
\end{equation*}
$$

Proof. We need to prove two assertions:

$$
\begin{align*}
& \forall \epsilon>0 \exists \Omega^{G}: \bar{\Delta}\left(\Omega^{G}\right)<\bar{\Delta}(\Omega)+\epsilon ;  \tag{3.3.3.5}\\
& \forall \epsilon>0 \exists \Omega^{K}: \underline{\Delta}\left(\Omega^{K}\right)>\underline{\Delta}(\Omega)-\epsilon . \tag{3.3.3.6}
\end{align*}
$$

Let us prove (3.3.3.5). Set

$$
\Omega^{G}(\epsilon):=\operatorname{Co}_{\Omega}\left(I_{1+\epsilon}\right) \cup\{|x|<\epsilon\} .
$$

This is an open set that contains $\mathrm{Co}_{\Omega}\left(I_{1}\right)$. One can show the following:
Exercise 3.3.3.1 For every open set $G \supset \operatorname{Co}_{\Omega}\left(I_{1}\right)$ there exists $\epsilon>0$ and $\Omega^{G} \subset S_{1}$ such that $\Omega^{G}(\epsilon) \subset G$.

We will show

$$
\begin{equation*}
\bar{\Delta}\left(\Omega^{G}(\epsilon)\right)<\bar{\Delta}\left(\Omega^{G}\right)+o(1) \tag{3.3.3.7}
\end{equation*}
$$

uniformly with respect to $\Omega^{G} \subset S_{1}$ while $\epsilon \rightarrow 0$.
We have from Exercise 3.3.1.1,

$$
\begin{equation*}
\bar{\Delta}\left(\Omega^{G}(\epsilon)\right) \leq \bar{\Delta}\left(\operatorname{Co}_{\Omega}\left(I_{1+\epsilon}\right)\right)+\bar{\Delta}(\{|x|<\epsilon\}) . \tag{3.3.3.8}
\end{equation*}
$$

The property dens7), Theorem 3.3.1.1, gives

$$
\bar{\Delta}\left(\operatorname{Co}_{\Omega^{G}}\left(I_{1+\epsilon}\right)\right)=\bar{\Delta}\left(\operatorname{Co}_{\Omega^{G}}\left(I_{1}\right)\right)(1+\epsilon)^{\rho+m-2}
$$

Since $\bar{\Delta}\left(\operatorname{Co}_{\Omega^{G}}\left(I_{1}\right)\right) \leq \bar{\Delta}(\{|x|<1\})$ we have

$$
\begin{equation*}
\bar{\Delta}\left(\mathrm{Co}_{\Omega^{G}}\left(I_{1+\epsilon}\right)\right)=\bar{\Delta}\left(\mathrm{Co}_{\Omega^{G}}\left(I_{1}\right)\right)+o(1) \tag{3.3.3.9}
\end{equation*}
$$

uniformly with respect to $\Omega^{G} \subset S_{1}$ as $\epsilon \rightarrow 0$.
By dens7) we also have

$$
\begin{equation*}
\bar{\Delta}(\{|x|<\epsilon\})=\bar{\Delta}(\{|x|<1\}) \epsilon^{\rho+m-2}=o(1) \tag{3.3.3.10}
\end{equation*}
$$

From (3.3.3.10), (3.3.3.9) and (3.3.3.8) we obtain (3.3.3.7). Hence (3.3.3.5) is proved.

Let us prove (3.3.3.6). Set

$$
\Omega^{K}(\epsilon):=\operatorname{Co}_{\Omega^{K}}\left(\bar{I}_{1-\epsilon}\right) \backslash\{|x|<\epsilon\}
$$

where $\bar{I}$ is the closure of $I$.
One can show the following:
Exercise 3.3.3.2 For any compact $K \subset \operatorname{Co}_{\Omega}\left(I_{1}\right)$ there exist $\Omega^{K} \subset \Omega$ and $\epsilon>0$ such that $K \subset \Omega^{K}(\epsilon) \subset \operatorname{Co}_{\Omega}\left(I_{1}\right)$.

From the definition of $\underline{\Delta}(\Omega)$ and the monotonicity we obtain (3.3.3.6).
Proof of Theorem 3.3.3.1. Suppose (3.3.3.1) holds. The property dens7), Theorem 3.3.1.1, implies (3.3.3.1) for all the $t \in(0, \infty)$. Set $\Delta(\Omega):=\bar{\Delta}(\Omega)=\underline{\Delta}(\Omega)$ for $\Omega \in \Phi$. Let us prove that $\Delta$ satisfies the conditions $\Delta 1)-\Delta 3$ ) from Section 2.2.3. The conditions $\Delta 1$ ) and $\Delta 2$ ) follow from dens3) and dens4), Theorem 3.3.1.1, Exercise 3.3.1.1.

Let us prove $\Delta 3$ ). By Theorem 3.3.3.2 for arbitrary $\Omega \in \Phi$ and $\epsilon>0$ we can choose $\Omega^{G} \supset \Omega$ such that $\bar{\Delta}(\Omega)>\bar{\Delta}\left(\Omega^{G}\right)-\epsilon$ and $\Omega^{K} \subset \Omega$ such that $\underline{\Delta}(\Omega)<\underline{\Delta}\left(\Omega^{K}\right)+\epsilon$.

Suppose $\Omega^{\prime} \in \Phi$ satisfies the condition $\Omega^{K} \subset \Omega^{\prime} \subset \Omega^{G}$. Then

$$
\Delta\left(\Omega^{\prime}\right)=\bar{\Delta}\left(\Omega^{\prime}\right) \leq \bar{\Delta}\left(\Omega^{G}\right) \leq \bar{\Delta}(\Omega)+\epsilon=\Delta(\Omega)+\epsilon
$$

and

$$
\Delta(\Omega)-\epsilon=\underline{\Delta}(\Omega)-\epsilon \leq \underline{\Delta}\left(\Omega^{K}\right) \leq \underline{\Delta}\left(\Omega^{\prime}\right)=\Delta\left(\Omega^{\prime}\right)
$$

implying $\Delta 3$ ).

## Chapter 4

## Structure of Limit Sets

### 4.1 Dynamical systems

4.1.1 The most complete and effective description of an arbitrary limit set can be done in terms of dynamical systems (see, [An]).

A family of the form

$$
T^{t}: M \mapsto M, t \in \mathbb{R}
$$

on a compact metric space $(M, d)$ with a metric $d(\bullet, \bullet)$ is a dynamical system $\left(T^{\bullet}, M\right)$ if it satisfies the condition

$$
T^{t+\tau}=T^{t} \circ T^{\tau}, t, \tau \in \mathbb{R}
$$

and the map $(t, m) \mapsto T^{t} m$ is continuous with respect to $(t, m)$, for all $t \in \mathbb{R}$, $m \in M$.

Let $m, m^{\prime} \in M$, and $\epsilon, s>0$. An $(\epsilon, s)$-chain from $m$ to $m^{\prime}$ is a finite sequence $m_{0}=m, m_{1}, \ldots, m_{n}=m^{\prime}$, satisfying the conditions $d\left(T^{t_{j}} m_{j}, m_{j+1}\right)<\epsilon, j=$ $0,1, \ldots, n-1$, for some $t_{j} \geq s$.

A dynamical system $\left(T^{\bullet}, M\right)$ is called chain recurrent (see, [HS]), if for an arbitrarily small $\epsilon>0$ and an arbitrarily large $s>0$ there exists an $(\epsilon, s)$-chain in $M$ from $m$ to $m$.

Theorem 4.1.1.1 (Properties of Chain Recurrence) Let $\left(T^{\bullet}, M\right)$ be a dynamical system on a compact set. Then the following conditions are equivalent:
cr1) $\quad M$ is connected and $\left(T^{\bullet}, M\right)$ is chain recurrent;
cr2) for every open proper $U \subset M$ satisfying

$$
\begin{equation*}
T^{t} U \subset U,-\infty<t<0 \tag{4.1.1.1}
\end{equation*}
$$

the boundary $\partial U$ contains a nonempty $T^{\bullet}$-invariant subset of $M$;
cr3) for every closed proper $K \subset M$ satisfying

$$
\begin{equation*}
T^{t} K \subset K, t \geq 0 \tag{4.1.1.2}
\end{equation*}
$$

the boundary $\partial K$ contains a nonempty $T^{\bullet}$-invariant subset of $M$;
cr4) there does not exist any open proper $V \subset M$ satisfying $T^{\tau} \operatorname{clos} V \subset V$ for some $\tau>0$;
cr5) for any small $\epsilon>0$, large $s>0$, and every pair of points $m, m^{\prime}$ there exists an $(\epsilon, s)$-chain from $m$ to $\mathrm{m}^{\prime}$.

Proof. The conditions cr2) and cr3) are equivalent. Let us prove, for example, $\mathrm{cr} 2) \Longrightarrow \mathrm{cr} 3)$. Set $U:=M \backslash K$. It is open. Applying to (4.1.1.2) $T^{-t}$ and, using the invariance of $M$, we obtain (4.1.1.1) for $U$. Hence $\partial U$ contains a nonempty invariant subset of $M$. Since $\partial K=\partial U$ we obtain cr2).

Let us prove the implication cr1 $\Longrightarrow \mathrm{cr} 3$ ). Let $K \subset M$ be closed, proper and satisfy (4.1.1.2). Since $M$ is proper $\partial K$ is nonempty.

Let $W$ denote the interior of $K$ in $M$. The continuity of $T$ and (4.1.1.2) imply

$$
\begin{equation*}
T^{t} W \subset W \tag{4.1.1.3}
\end{equation*}
$$

for $t \geq 0$. Indeed, $T^{t} W \subset K$. It must be open. Thus it cannot contain any point of $\partial K$, since else it would contain some neighborhood of this point, contradicting the definition of $\partial K$.

Suppose that $\partial K$ does not contain any nonempty $T$-invariant set. Let us show that there exists $s>0$ such that

$$
\begin{equation*}
T^{s} K \subset W \tag{4.1.1.4}
\end{equation*}
$$

For any $m \in \partial K$ there exists $t=t(m)$ such that $T^{t} m \in W$. There exists a neighborhood $V_{m}$ of $m$ in $\partial K$ that passes to $W$ under $T^{t(m)}$-action because of continuity of $T^{t} m$ on $m$.

We also have $T^{t} V_{m} \subset W$ for $t>t(m)$ because of (4.1.1.3). Since $\partial K$ is compact we can cover it by a finite number of neighborhoods and obtain $s$ such that

$$
\begin{equation*}
T^{s} \partial K \subset W \tag{4.1.1.5}
\end{equation*}
$$

(4.1.1.5) and (4.1.1.3) give (4.1.1.4).

Set $\epsilon:=0.5 d\left(\partial K, T^{s} K\right)$. From (4.1.1.2) we see that $T^{t} K \subset T^{s} K$ for $t>s$. Therefore there does not exist any ( $\epsilon, s$ )-chain from a small neighborhood of a point $m \in \partial K$ to itself. This contradicts the chain recurrence of $M$.

Let us prove cr3) $\Longrightarrow \mathrm{cr} 4$ ). Assume that there exists an open proper $V \subset M$ satisfying $T^{\tau} \operatorname{clos} V \subset V$ for some $\tau>0$.

We will construct $K$ that does not satisfy cr3). Set $W:=\bigcup_{0 \leq t \leq \tau} T^{t} V$ and $K:=\operatorname{clos} W$.

Then

$$
\begin{equation*}
T^{s} W \subset W, \forall s \geq 0 \tag{4.1.1.6}
\end{equation*}
$$

Indeed, let $s=k \tau+s^{\prime}, s^{\prime} \in[0, \tau), k \in \mathbb{Z}$. Then

$$
\begin{equation*}
T^{s} W=\bigcup_{t \in[0, \tau]} T^{t+s} V \tag{4.1.1.7}
\end{equation*}
$$

Since $T^{\tau} V \subset V$ we have $T^{t+k \tau} V \subset T^{t} V$ for $t>0$. From (4.1.1.7) we obtain

$$
\begin{aligned}
T^{s} W & =\bigcup_{t \in[0, \tau]} T^{t+s^{\prime}+k \tau} V \subset \bigcup_{t \in[0, \tau]} T^{t+s^{\prime}} V=\bigcup_{t^{\prime} \in\left[s^{\prime}, \tau+s^{\prime}\right]} T^{t^{\prime}} V \\
& =\bigcup_{t \in\left[s^{\prime}, \tau\right]} T^{t} V \cup \bigcup_{t \in\left[\tau, \tau+s^{\prime}\right]} T^{t} V:=W_{1} \cup W_{2}
\end{aligned}
$$

Further we have $W_{1} \subset W$ by definition. $W_{2}$ can be represented in the form

$$
W_{2}=\bigcup_{t \in\left[0, s^{\prime}\right]} T^{t+\tau} V
$$

Since

$$
T^{t+\tau} V=T^{t} T^{\tau} V \quad \text { and } \quad T^{\tau} V \subset V
$$

by the assumption we get:

$$
W_{2} \subset W_{1} \subset W
$$

This implies (4.1.1.6). The same holds for $K$ because of continuity of $T^{t}$, i.e., $K$ satisfies (4.1.1.2).

Let us prove the equality

$$
\begin{equation*}
K=\bigcup_{0 \leq t \leq \tau} T^{t} \cos V \tag{4.1.1.8}
\end{equation*}
$$

Denote as $K^{\prime}$ the right side of (4.1.1.8).
The set $K^{\prime}$ is closed because of compactness of $[0, \tau]$. Indeed, let the sequence $\left\{T^{t_{j}} v_{j}: j=1,2, \ldots\right\} \in T^{t_{j}}(\operatorname{clos} V)$ converge to $w$. Choose a subsequence $t_{j_{k}} \rightarrow$ $s \in[0, \tau]$. Then

$$
v:=\lim _{k \rightarrow \infty} v_{j_{k}}=\lim _{k \rightarrow \infty} T^{-t_{j_{k}}} w=T^{-s} w .
$$

Since clos $V$ is closed, $v \in \operatorname{clos} V$. Thus $w=T^{s} v$ for some $s \in[0, \tau]$ and some $v \in \operatorname{clos} V$, i.e., $w \in K^{\prime}$.

Now, $W \subset K^{\prime}$ because

$$
T^{t} V \subset \cos T^{t} V=T^{t} \cos V
$$

Hence,

$$
K:=\operatorname{clos} W \subset \operatorname{clos} K^{\prime}=K^{\prime} .
$$

We also have

$$
\left(T^{t} V \subset W \forall t \in[0, \tau]\right) \Longrightarrow\left(\cos T^{t} V=T^{t} \cos V \subset \cos W=K, \forall t \in[0, \tau]\right)
$$

Hence, $K^{\prime} \subset K$. Therefore $K=K^{\prime}$, i.e., (4.1.1.8) holds.
From (4.1.1.8) and $T^{\tau} \operatorname{clos} V \subset V$ we obtain $T^{\tau} \cos W \subset W$. Hence $T^{\tau} \partial K \subset$ $W$. This and $\partial K \cap W=\varnothing$ imply

$$
\begin{equation*}
T^{\tau} \partial K \cap \partial K=\varnothing \tag{4.1.1.9}
\end{equation*}
$$

To obtain a contradiction and complete the proof of cr 3$) \Longrightarrow \mathrm{cr} 4$ ) we have to show that $K$ is a proper subset, because both cases: $\partial K=\varnothing$ and $\partial K \neq \varnothing$ will contradict cr3).

Since $V$ is proper $T^{t} V$ is proper for any $t \in(-\infty, \infty)$. Otherwise $T^{t} V=M$ implies $V=T^{-t} M=M$, which is a contradiction.

Since $V$ is a neighborhood of the compact set $T^{\tau} \operatorname{clos} V$ we can find $\alpha>0$ such that $T^{t} \circ T^{\tau} \cos V \subset V$ for $t \in[0, \alpha]$. Then $T^{t} \cos V \subset T^{-\tau} V$ for $t \in[0, \alpha]$.

By iteration of this inclusion we obtain $T^{j t} \cos V \subset T^{-j \tau} V$ for any integer $j$. When $j \alpha>\tau$ it follows that $K \subset T^{-j \tau} V$. The last set is proper because we mentioned already that $T^{t} V$ is proper for any $t \in(-\infty, \infty)$. Hence $K$ is proper.

So $K$ satisfies the conditions of cr3) but $\partial K$ does not contain a nonempty $T^{\bullet}$-invariant set. This contradiction proves the implication cr 3$) \Longrightarrow \mathrm{cr} 4$ ).

Let us prove cr4) $\Longrightarrow$ cr5). Let $\epsilon>0$ be small and $s>0$ be large. Let $V$ denote the set of all $m^{\prime} \in M$ such that there exists an $(\epsilon, s)$-chain from $m$ to $m^{\prime}$. This set is open and closed. Indeed, let $m^{\prime} \in V$. There exists an $(\epsilon, s)$-chain $m=m_{0}, \ldots, m_{n-1}, m_{n}=m^{\prime}$ from $m$ to $m^{\prime}$. Choose $\epsilon_{1}<\epsilon-d\left(m_{n}, m_{n-1}\right)$ and consider the closed neighborhood $W:=\left\{m^{\prime \prime}: d\left(m^{\prime}, m^{\prime \prime}\right) \leq \epsilon_{1}\right\}$. Then for any $m^{\prime \prime} \in W$ the chain $m=m_{0}, \ldots, m_{n-1}, m_{n}=m^{\prime \prime}$ is an $(\epsilon, s)$-chain from $m$ to $m^{\prime \prime}$. Hence, with every point, $V$ contains its closed neighborhood. Therefore it is open and closed. Therefore it is a connected component of $M$.

We also have $T^{s} m \in V$ because for that case $n=1, m_{0}=m, m_{1}=T^{s} m$. Hence $T^{s} \operatorname{clos} V \subset V$. If $V$ does not coincide with the whole $M$ the latter contradicts cr4). Hence $V=M$.

Finally, let us prove cr5) $\Longrightarrow \mathrm{cr} 1$ ). If $M$ is a union of two nonempty disjoint sets $A$ and $B$, then both of them are open and closed. Since $M$ is compact, the distance $\epsilon$ between $A$ and $B$ is positive. Hence every $(\epsilon / 2, s)$-chain starting at a point of $A$ remains in $A$, contradicting cr5).

Since for every point $m \in M$ the set $V$ from the proof of cr 4$) \Longrightarrow \operatorname{cr} 5$ ) coincides with $M$, cr1) holds.

Theorem 4.1.1.2 Let $T^{\bullet}$ be chain recurrent on $M_{\alpha}, \alpha \in A$. Then $T^{\bullet}$ is chain recurrent on $M=\bigcup_{\alpha \in A} M_{\alpha}$.

This is because every $(\epsilon, s)$-chain from $m$ to $m^{\prime}$ in $M_{\alpha}$ is also $(\epsilon, s)$-chain in $M$.
4.1.2 Here we prove two auxiliary assertions that will be used further.

Theorem 4.1.2.1 Let $T^{\bullet}$ be chain recurrent on a connected compact $M$ and let $\left\{q_{j}\right\}$ be a sequence in $M$. Then there exist sequences $\left\{\alpha_{\nu}\right\}$ and $\left\{\omega_{\nu}\right\}$ of real numbers and a sequence $\left\{p_{\nu}\right\}$ in $M$ having $\left\{q_{j}\right\}$ as a subsequence, such that

$$
\begin{equation*}
\alpha_{\nu} \rightarrow-\infty ; \omega_{\nu} \rightarrow \infty \tag{4.1.2.1}
\end{equation*}
$$

and

$$
\begin{equation*}
d\left(T^{\omega_{\nu}} p_{\nu}, T^{\alpha_{\nu+1}} p_{\nu+1}\right) \rightarrow 0 \tag{4.1.2.2}
\end{equation*}
$$

as $\nu \rightarrow \infty$.
Proof. In addition to $\left\{\alpha_{\nu}\right\},\left\{\omega_{\nu}\right\}$ and $\left\{p_{\nu}\right\}$ we define, by induction, a sequence $\left\{\epsilon_{\nu}\right\}$ of positive real numbers, tending to zero, and an increasing sequence $\left\{\nu_{j}\right\}$ of positive integers, such that $\left\{p_{\nu_{j}}\right\}=\left\{q_{j}\right\}$ and

$$
\begin{equation*}
d\left(T^{\omega_{\nu}} p_{\nu}, T^{\alpha_{\nu+1}} p_{\nu+1}\right)<\epsilon_{\nu}, \nu=1,2, \ldots \tag{4.1.2.3}
\end{equation*}
$$

We start by setting $\alpha_{1}=-1, \epsilon_{1}=1, \nu_{1}=1, \omega_{1}=5$ and $p_{1}=q_{1}$. Assume now that $\alpha_{\nu}, \epsilon_{\nu}, \omega_{\nu}$ and $p_{\nu}$ have been chosen for $\nu=1,2, \ldots, \nu_{j}$. Set

$$
\begin{equation*}
\alpha=\alpha_{\nu_{j}}-1, \epsilon=\epsilon_{\nu_{j}} / 2, \omega=\omega_{\nu_{j}} . \tag{4.1.2.4}
\end{equation*}
$$

By Theorem 4.1.1.1, cr5) there exists a sequence $r_{0}:=T^{\omega} q_{j}, r_{1}, \ldots, r_{m}:=T^{\alpha} q_{j+1}$ such that $d\left(T^{t_{k}} r_{k}, r_{k+1}\right)<\epsilon$ for $k=0,1, \ldots, m-1$, where $t_{k} \geq \omega$. Now we set $\nu_{j+1}=\nu_{j}+m+1$. For $\nu=\nu_{j}+k+1, k=0,1, \ldots, m-1$, we set $\alpha_{\nu}=-t_{k} / 2, \omega_{\nu}=$ $t_{k} / 2, p_{\nu}=T^{t_{k} / 2} r_{k}$, and finally, for $\nu=\nu_{j+1}$ we set $\alpha_{\nu}=\alpha, \epsilon_{\nu}=\epsilon, \omega_{\nu}=\omega+1, p_{\nu}=$ $q_{j+1}$.

Let us check that with this setting the properties (4.1.2.1) hold . Since $\omega_{\nu_{j+1}}=$ $\omega_{\nu_{j}}+1$ we have $\omega_{\nu_{j}} \rightarrow \infty$ as $j \rightarrow \infty$. From $t_{k} \geq \omega=\omega_{\nu_{j}}$ we obtain $\alpha_{\nu} \rightarrow-\infty$ and $\omega_{\nu} \rightarrow \infty$. Hence (4.1.2.1) holds.

One can see from (4.1.2.4) that $\epsilon_{\nu}=\epsilon_{\nu_{j}} / 2 \rightarrow 0$. To prove (4.1.2.2) it is enough to check (4.1.2.3). For $k=0$ we have

$$
p_{\nu}=T^{t_{0} / 2} r_{0}=T^{\left(t_{0} / 2\right)+\omega} q_{j}=T^{\left(t_{0} / 2\right)+\omega} p_{\nu_{j}} .
$$

Hence,

$$
T^{\alpha_{\nu}} p_{\nu}=T^{\omega} p_{\nu_{j}}=T^{\omega_{\nu_{j}}} p_{\nu_{j}}
$$

Thus

$$
\begin{equation*}
d\left(T^{\omega_{\nu_{j}}} p_{\nu_{j}}, T^{\alpha_{\nu}} p_{\nu}\right)=0 \tag{4.1.2.5}
\end{equation*}
$$

for this case.
For $k=1, \ldots, m-2$ and the corresponding $\nu$ we have

$$
T^{\omega_{\nu}} p_{\nu}=T^{t_{k} / 2} \circ T^{t_{k} / 2} r_{k}=T^{t_{k}} r_{k}
$$

and

$$
T^{\alpha_{\nu+1}} p_{\nu+1}=T^{-t_{k+1} / 2} \circ T^{t_{k+1} / 2} r_{k+1}=r_{k+1}
$$

Hence,

$$
\begin{equation*}
d\left(T^{\omega_{\nu}} p_{\nu}, T^{\alpha_{\nu+1}} p_{\nu+1}\right)=d\left(T^{t_{k}} r_{k}, r_{k+1}\right)<\epsilon=\epsilon_{\nu} \tag{4.1.2.6}
\end{equation*}
$$

Finally, for the last link of the chain we obtain

$$
\begin{gathered}
k=m-1, \nu=\nu_{j}+m, \nu+1=\nu_{j+1}, \alpha_{\nu_{j+1}}=\alpha, \\
T^{\alpha_{\nu+1}} p_{\nu+1}=T^{\alpha_{\nu_{j+1}}} p_{\nu_{j+1}}=T^{\alpha} q_{j+1}=r_{m} .
\end{gathered}
$$

Thus (4.1.2.6) holds for $k=m-1$. Hence, (4.1.2.3) also holds. Therefore (4.1.2.2) holds.

Lemma 4.1.2.2 Let $p_{k}, q_{k} \in M$ and $d\left(p_{k}, q_{k}\right) \rightarrow 0$ as $k \rightarrow \infty$. Then there exists a sequence $\left\{\gamma_{k} \uparrow \infty\right\}$ such that

$$
\begin{equation*}
d\left(T^{\tau} p_{k}, T^{\tau} q_{k}\right) \rightarrow 0 \tag{4.1.2.7}
\end{equation*}
$$

uniformly with respect to $\tau \in\left[-\gamma_{k+1}, \gamma_{k}\right]$.

Proof. Let $[-\gamma, \gamma]$ be a fixed segment. Then $d\left(T^{\tau} p_{k}, T^{\tau} q_{k}\right) \rightarrow 0$ uniformly in this segment.

Indeed, suppose there exist sequences $\tau_{j}, k_{j}$ such that $d\left(T^{\tau_{j}} p_{k_{j}}, T^{\tau_{j}} q_{k_{j}}\right) \geq \epsilon>$ 0 . Choosing a subsequence we can assume that $\tau_{j} \rightarrow \tau \in[-\gamma, \gamma], p_{k_{j}} \rightarrow p \in M$ and $q_{k_{j}} \rightarrow q=p$. Using continuity of $T^{\tau} m$ on $(\tau, m)$ and continuity of $d(\bullet \bullet)$ in both arguments we obtain $0=d(p, p) \geq \epsilon>0$. This is impossible.

Denote

$$
\epsilon(\gamma, k):=\max _{\tau \in[-\gamma, \gamma]} d\left(T^{\tau} p_{k}, T^{\tau} q_{k}\right)
$$

This function increases monotonically in $\gamma$ and tends to zero for any $\gamma$ as $k \rightarrow \infty$.
Choose $l_{n}$ such that $\epsilon(n, k) \leq 1 / n$ for $k \geq l_{n}$. Set $\gamma_{k+1}:=n$ for $l_{n}<k \leq l_{n+1}$. One can see that $\epsilon\left(\gamma_{k+1}, k\right) \rightarrow 0$ as $k \rightarrow \infty$. Since

$$
\max _{\tau \in\left[-\gamma_{k+1}, \gamma_{k}\right]} d\left(T^{\tau} p_{k}, T^{\tau} q_{k}\right) \leq \epsilon\left(\gamma_{k+1}, k\right)
$$

$\left\{\gamma_{k}\right\}$ satisfies (4.1.2.7).
4.1.3 We connect the property of being chain recurrent with other well-known characteristics of dynamical systems ([AGL]).

A point $m_{0} \in M$ is called non-wandering (see [An]) if for any neighborhood $\mathcal{O}$ of $m_{0}$ and arbitrarily large number $s \in \mathbb{R}$ there exists $m \in \mathcal{O}$ and $t \geq s$ such that $T^{t} m \in \mathcal{O}$.

This means that the "returns" take place to an arbitrarily small neighborhood of the point $m_{0}$. We shall denote as $\Omega\left(T^{\bullet}\right)$ the set of non-wandering points. It is a closed invariant subset of $M$.

The set $A \subset M$ is called an attractor if it satisfies the following conditions:
attr1) for any neighborhood $\mathcal{O} \supset A$ there exists a neighborhood $\mathcal{O}^{\prime}, A \subset \mathcal{O}^{\prime} \subset \mathcal{O}$ such that $T^{t} \mathcal{O}^{\prime} \subset \mathcal{O} t \in \mathbb{R}$, where $T^{t} \mathcal{O}^{\prime}$ is the image of $\mathcal{O}^{\prime}$;
attr2) there exists a neighborhood $\mathcal{O} \supset A$ such that $T^{t} m \rightarrow A$ when $t \rightarrow \infty$ for $m \in \mathcal{O}$.

Theorem 4.1.3.1 If $\Omega\left(T^{\bullet}\right)=M$, then $\left(T^{\bullet}, M\right)$ is chain recurrent; if $\left(T^{\bullet}, M\right)$ has an attractor $A \neq M$, it is not chain recurrent.

Proof. The property $\Omega\left(T^{\bullet}\right)=M$ obviously implies the chain recurrence for $m=$ 1.Suppose there exists an attractor $A \neq M$. Take a point $m_{0}$ that does not belong to $A$ and choose a neighborhood $\mathcal{O} \supset A$ such that $d\left(m_{0}, \cos \mathcal{O}\right)=2 \epsilon>0$. This is possible because an attractor is closed. Let $\mathcal{O}^{\prime}$ be chosen by attr1) and $s$ be such that $T^{s} m \in \mathcal{O}^{\prime}$. Then there does not exist any $(\epsilon, s)$-chain from a small neighborhood of $m_{0}$ itself. By definition $\left(T^{t}, M\right)$ is not chain recurrent.

Let us give examples of dynamical systems on connected compacts that are chain recurrent.

Theorem 4.1.3.2 Let $M$ be a connected compact and let $T^{t}(-\infty<t<\infty)$ be the identity map. Then $\left(T^{\bullet}, M\right)$ is chain recurrent.

This theorem, of course, is trivial. However, if $M$ consists of a single point this dynamical system determines an important class of subharmonic and entire functions of completely regular growth (see [Le, Ch. III]).

Let $m \in M$. Set

$$
\begin{equation*}
\mathbb{C}(m):=\operatorname{clos}\left\{T^{t} m:-\infty<t<\infty\right\} . \tag{4.1.3.1}
\end{equation*}
$$

It is closed, connected and invariant.
Exercise 4.1.3.1 Prove this.
Let us denote as $\Omega(m)$ the set of all limits of the form

$$
\begin{equation*}
\Omega(m):=\left\{m^{\prime} \in M:\left(\exists t_{k} \rightarrow \infty\right)\left(m^{\prime}=\lim _{k \rightarrow \infty} T^{t_{k}} m\right\}\right. \tag{4.1.3.2}
\end{equation*}
$$

This is a limit set as $t \rightarrow \infty$. It is the "tangle" at the end of the curve. Denote by $A(m)$ the analogous set for $t \rightarrow-\infty$.

Exercise 4.1.3.2 Prove that $A(m)$ and $\Omega(m)$ are invariant.
Theorem 4.1.3.3 $\left(T^{\bullet}, \mathbb{C}(m)\right)$ is chain recurrent iff

$$
\begin{equation*}
A(m) \cap \Omega(m) \neq \varnothing . \tag{4.1.3.3}
\end{equation*}
$$

Proof. Suppose $B:=A(m) \cap \Omega(m)=\varnothing$. Then $\Omega(m)$ is an attractor and $\left(T^{\bullet}, \mathbb{C}(m)\right)$ is not chain recurrent by Theorem 4.1.3.1.

Suppose $B \neq \varnothing$. We will use cr2) from Theorem 4.1.1.1.
Let $U$ be an open proper subset of $\mathbb{C}(m)$ satisfying (4.1.1.1). Consider two cases:
i) $B$ contains a point of $U$. Thus $U$ contains a sequence of form $T^{t_{k}} m, t_{k} \rightarrow \infty$. From (4.1.1.1) we obtain that $U$ contains $T^{t} m$ for all $t \in(-\infty, \infty)$. Thus $U \supset \mathbb{C}(m)$ and $\operatorname{clos} U=\mathbb{C}(m)$. Set $K=\mathbb{C}(m) \backslash U$. One can show that $K$ satisfies (4.1.1.2)(see the beginning of proof of Theorem 4.1.1.1). Hence $K$ contains the set

$$
\begin{equation*}
K^{*}:=\bigcap_{t \geq 0} T^{t} K \tag{4.1.3.4}
\end{equation*}
$$

that is invariant (Exercise 4.1.3.3).
Therefore $K^{*} \subset K \subset \operatorname{clos} U \backslash U=\partial U$. By cr2) $\left(T^{t}, \mathbb{C}(m)\right)$ is chain recurrent.
ii) $B$ contains no point of $U$. Then $B \subset A(m) \subset \partial U$. By cr2) $\left(T^{\bullet}, \mathbb{C}(m)\right)$ is chain recurrent.

Exercise 4.1.3.3 Let $U$ satisfy (4.1.1.1) and $K:=M \backslash U$. Prove that $K^{*}$ from (4.1.3.4) is invariant.
4.1.4 The connectedness of $M$ is a necessary condition for a dynamical system to be chain recurrent.

Let $M$ be a subset of a linear space. The set $M$ is called polygonally connected if every pair of points $m_{1}, m_{2}$ can be connected by a polygonal path.

Of course, polygonal connectedness implies connectedness and even arcwise connectedness.

Theorem 4.1.4.1 Let $\left(T^{\bullet}, M\right)$ be a dynamical system such that $M$ is a polygonally connected set. Then $\left(T^{\bullet}, M\right)$ is chain recurrent.

Proof. Let $U$ be an open proper subset of $M$, satisfying (4.1.1.1). We choose $m_{1} \in U$ and $m_{2}$ in an invariant subset $K^{*}$ of $K:=M \backslash U$. Then there exists a polygonal path from $m_{1}$ to $m_{2}$ :

$$
\begin{aligned}
m_{\theta}:= & (j+1-\theta) m_{j}^{\prime}+(\theta-j) m_{j+1}^{\prime}, \text { for } \theta \in[j, j+1], \\
& j=0,1, \ldots, l-1 ; m_{0}^{\prime}:=m_{1}, m_{l}^{\prime}:=m_{2} .
\end{aligned}
$$

Now $M$ is invariant, so for each $t$ the continuous path $\theta \mapsto T^{t} m_{\theta}$ lies in $M$.
If $t \in(-\infty, 0)$ its initial point $T^{t} m_{1}$ belongs to $U$ and its endpoint $T^{t} m_{2}$ belongs to $K^{*} \subset K$.

For each $t \in(-\infty, 0)$ we set

$$
\theta(t):=\min \left[\theta \in[0 ; l]: T^{t} m_{\theta} \in K\right]
$$

Then $\theta(t)>0, T^{t} m_{\theta(t)} \in \partial U$ and (4.1.1.1) implies that $t \mapsto \theta(t)$ is a decreasing function. Hence the limit

$$
\theta(-\infty):=\lim _{t \rightarrow-\infty} \theta(t)
$$

exists and is positive.
Set $m_{3}:=m_{\theta(-\infty)}$. We claim that $A\left(m_{3}\right) \subset \partial U(A(\cdot)$ is a set defined before Theorem 4.1.3.3). If $\theta(-\infty) \in(j, j+1]$ for some $j \in[0, l]$ then $\theta(t) \in(j, j+1]$ for $t$ that is near to $-\infty$, and

$$
T^{t} m_{3}=T^{t} m_{\theta(t)}+(\theta(t)-\theta(-\infty)) T^{t} m_{j}^{\prime}+(\theta(-\infty)-\theta(t)) T^{t} m_{j+1}^{\prime}
$$

The first term in the right-hand side lies in $\partial U$. The set $M$ is compact and invariant so the other terms tend to zero as $t \rightarrow-\infty$. Hence $A\left(m_{3}\right) \subset \partial U$.

Thus $\partial U$ contains this invariant subset and $\left(T^{\bullet}, M\right)$ is chain recurrent by cr2), Theorem 4.1.1.1.

We have the obvious
Corollary 4.1.4.2 Let $\left(T^{\bullet}, M\right)$ be a dynamical system such that $M$ is a compact convex set. Then $\left(T^{\bullet}, M\right)$ is chain recurrent.

This is because the polygonal path can be taken as a line segment connecting every pair of points.
4.1.5 Let $U[\rho, \sigma]$ be a set of subharmonic functions defined in (3.1.2.4). It is invariant with respect to the transformation $(\bullet)_{[t]}$ defined in (3.1.2.4a).

Set (subindex!)

$$
\begin{equation*}
T_{t} v:=v_{\left[e^{t}\right]} \tag{4.1.5.1}
\end{equation*}
$$

Since $(\bullet)_{[t]}$ has the property (3.1.2.4b)

$$
\begin{equation*}
T_{t+\tau} v=\left(T_{t} \circ T_{\tau}\right) v, \forall t, \tau \in \mathbb{R} \tag{4.1.5.2}
\end{equation*}
$$

By Theorem 3.1.2.3, $T_{\bullet}$ is continuous in the appropriate topology and hence $\left(T_{\bullet}, U[\rho, \sigma]\right)$ is a dynamical system.
Theorem 4.1.5.1 (Universality of $U[\rho, \sigma])$ Let $\left(T^{\bullet}, M\right)$ be a chain recurrent dynamical system on a compact set $M$. Then for any $\rho, \sigma$ there exists $U \subset U[\rho, \sigma]$ and a homeomorphism $\mathrm{imb}: M \mapsto U$ such that $\mathrm{imb} \circ T^{t}=T_{t} \circ \mathrm{imb}, t \in(-\infty, \infty)$.
I.e., any dynamical system can be imbedded in $\left(T_{\bullet}, U[\rho, \sigma]\right)$.

This plot is developed in $[\mathrm{Az}(2008)]$.
It is sufficient to prove the theorem by supposition $P_{t} x=t x$ because $\left(T_{\bullet}^{P}, U[\rho, \sigma]\right)$ is a dynamical system for any $P_{t}$ and

$$
\operatorname{imb}:\left(T_{\bullet}, U[\rho, \sigma]\right) \mapsto\left(T_{\bullet}^{P}, U[\rho, \sigma]\right)
$$

where imb : $u(x) \mapsto T_{-t} T_{t}^{P} u(x)$ is also a homeomorphism of dynamical systems.
Exercise 4.1.5.1 Consider Theorem 3.1.6.1 from this point of view.

We need some auxiliary definitions and results. Let us denote as $\mathcal{M}\left(S^{m-1}\right)$ the set of measures $\nu$ with bounded full variation on the unit sphere $S^{m-1}$. Introduce the metric $d(\nu, 0):=\operatorname{Var} \nu$ and consider the set

$$
K:=\{\nu: \nu>0, d(\nu, 0) \leq 1\}
$$

i.e., the intersection of the cone of positive measures with the unit ball.

The following assertion is a corollary of Keller's theorem (see, e.g., [BP, Thm. 3.1, p. 100]).

Theorem 4.1.5.2 (Imbedding) Every metric compact set can be homeomorphically imbedded to $K$.

Thus we can assume below that for any $m \in M$ there exists a positive measure

$$
Y(\bullet, m)=Y\left(d x^{0}, m\right) \in K
$$

such that

$$
\begin{equation*}
\left(Y\left(\bullet, m_{1}\right)=Y\left(\bullet, m_{2}\right)\right) \Longrightarrow\left(m_{1}=m_{2}\right) \tag{4.1.5.3}
\end{equation*}
$$

and $Y(\bullet, m)$ is continuous with respect to the metrics.
We also introduce a new coordinate system. For $x:=e^{y} x^{0} \in \mathbb{R}^{m} \backslash 0$ set $\operatorname{Pol}(x)=\left(y, x^{0}\right)$. This formula gives a one-to-one map from $\mathbb{R}^{m} \backslash 0$ onto the cylinder $C y l:=(-\infty, \infty) \times S^{m-1}$. Thus, for any $\left(y, x^{0}\right) \in C y l, \operatorname{Pol}^{-1}\left(y, x^{0}\right)=e^{y} x^{0}$.

For $m=2$ this is a common cylinder.

### 4.1. 6

Proof of Theorem 4.1.5.1. We consider separately the cases of integer and noninteger $\rho$.

Let $\rho$ be non-integer and $\sigma>0$. For any $v \in U[\rho, \sigma]$, one has the representation of Theorem 3.1.4.4 (*Hadamard),

$$
\begin{equation*}
v(x)=\Pi(x, \mu, \rho) \tag{4.1.6.1}
\end{equation*}
$$

where $\mu \in \mathcal{M}[\rho, \Delta]$ and $\Delta$ depends only on $\sigma$ (Theorem 2.8.3.3).
Vice versa, every $\mu \in \mathcal{M}[\rho, \Delta]$ generates $v$ by (4.1.6.1) and

$$
v_{[t]}(x)=\Pi\left(x, \mu_{[t]}, \rho\right) .
$$

Let us "transplant" $\mu$ in $C y l$. For $\mu$ that has a dense $f_{\mu}\left(r x^{0}\right)$, we set

$$
\nu\left(d y \otimes d x^{0}\right):=f_{\mu}\left(e^{y} x^{0}\right) e^{(-\rho-2) y}\left(d y \otimes d x^{0}\right)
$$

i.e., the density $f_{\nu}$ of $\nu$ is defined by

$$
f_{\nu}\left(x^{0}, y\right):=f_{\mu}\left(e^{y} x^{0}\right) e^{(-\rho-2) y}
$$

Respectively,

$$
f_{\mu}\left(x^{0}, r\right)=f_{\nu}\left(x^{0}, \log r\right) r^{\rho+2}
$$

We can extend this equality for all $\mu \in \mathcal{M}[\rho, \Delta]$ by using a limit process in $\mathcal{D}^{\prime}$ topology.

Exercise 4.1.6.3 Do that using, for example, Theorem 2.3.4.5 (Properties of Regularization).

We can also define $\nu$ as a distribution in $\mathcal{D}^{\prime}(C y l)$. Namely, for $\psi \in \mathcal{D}(C y l)$ we set

$$
\psi^{*}\left(x^{0}, r\right):=\psi\left(\operatorname{Pol}^{-1}\left(x^{0}, \log r\right)\right) r^{-\rho-m+2}
$$

and

$$
\langle\nu, \psi\rangle:=\int \psi^{*}\left(x^{0}, r\right) \mu\left(d x^{0} \otimes r^{m-1} d r\right)
$$

Exercise 4.1.6.4 Check that this definition gives the same $\nu$.
The transformation $P_{t} x=\left(x^{0}, t r\right), r x^{0} \in \mathbb{R}^{m} \backslash 0$ passes to

$$
\operatorname{Pol} \circ P_{t} \circ \operatorname{Pol}^{-1}\left(x^{0}, y\right)=\left(x^{0}, y+\log t\right) .
$$

Thus $T_{e^{\tau}} \mu$ gives a transformation $S_{\tau} \nu$ defined by

$$
S_{t} f_{\nu}\left(x^{0}, y\right):=f_{\nu}\left(x^{0}, y+t\right)
$$

for densities or by

$$
\begin{equation*}
\left\langle S_{t} \nu, \psi\right\rangle:=\int \psi\left(x^{0}, y-t\right) \nu\left(d x^{0} \otimes d y\right) \tag{4.1.6.2}
\end{equation*}
$$

for distributions $(\psi \in \mathcal{D}(C y l)$.)
Exercise 4.1.6.5 Check the equivalence.
From $\mu \in \mathcal{M}[\rho, \Delta]$ we obtain

$$
\begin{equation*}
\int_{y \leq 0} e^{(\rho+m-2) y} S_{t} \nu\left(d y \otimes d x^{0}\right) \leq \Delta, t \in \mathbb{R} \tag{4.1.6.3}
\end{equation*}
$$

Exercise 4.1.6.6 Check this.
Let $X(t)$ be a positive function satisfying the condition

$$
\int_{-\infty}^{\infty} X(t) d t=1
$$

and such that the linear hull of its translations are dense in $L^{1}(-\infty, \infty)$. We can choose, for example, the function

$$
X(t):=\frac{1}{\sqrt{2 \pi}} e^{-\frac{t^{2}}{2}}
$$

because its Fourier transformation does not vanish in $\mathbb{R}$ (it is $e^{-\frac{s^{2}}{2}}$ ).

Exercise 4.1.6.7 Check these properties.
Let us define $\nu(\bullet, m)$ by

$$
\begin{equation*}
\langle\nu(\bullet, m), \psi\rangle:=\int_{\left(x^{0}, y\right) \in C y l} \psi\left(x^{0}, y\right)\left(\rho \int_{-\infty}^{\infty} Y\left(d x^{0}, T^{y-t} m\right) X(t) d t\right) d y \tag{4.1.6.4}
\end{equation*}
$$

Now we check the property

$$
S_{\tau} \nu(\bullet, m)=\nu\left(\bullet, T^{\tau} m\right)
$$

Using (4.1.6.2), we obtain

$$
\begin{aligned}
\left\langle S_{\tau} \nu(\bullet, m), \psi\right\rangle & =\int \psi\left(x^{0}, y\right)\left(\rho \int_{-\infty}^{\infty} Y\left(d x^{0}, T^{y+\tau-t} m\right) X(t) d t\right) d y \\
& =\int \psi\left(x^{0}, y\right)\left(\rho \int_{-\infty}^{\infty} Y\left(d x^{0}, T^{y-t}\left(T^{\tau} m\right)\right) X(t) d t\right) d y \\
& =\left\langle\nu(\bullet, m), T^{\tau} m\right\rangle
\end{aligned}
$$

We also check the condition (4.1.6.3).

$$
\begin{aligned}
\int_{y \leq 0} e^{\rho y} S_{t} \nu\left(d y \otimes d x^{0}\right) & \left.=\int_{-\infty}^{\infty} X(t) d t \int_{y \leq 0} e^{\rho y}\right) \rho d y \int_{S^{m-1}} Y\left(d x^{0}, T^{y+\tau-t} m\right. \\
& \leq \sup _{\tau \in \mathbb{R}} Y\left(S^{m-1}, T^{\tau} m\right) \int_{-\infty}^{\infty} X(t) d t \leq 1
\end{aligned}
$$

since $Y(\bullet, \bullet) \in K$.
Now we should "transplant" $\nu$ back to $\mathbb{R}^{m} \backslash 0$ such that $S_{\tau}$ passes to $(\bullet)_{\left[e^{\tau}\right]}$. Define $\mu(\bullet, m)$ by

$$
\begin{equation*}
\left\langle\mu(\bullet, m), \psi^{*}\right\rangle:=\langle\nu(\bullet, m), \psi\rangle, \tag{4.1.6.5}
\end{equation*}
$$

where $\psi^{*}\left(r x^{0}\right) \in \mathcal{D}\left(\mathbb{R}^{m} \backslash 0\right)$ and

$$
\psi\left(x^{0}, y\right):=\psi^{*}\left(e^{y} x^{0}\right) e^{-(\rho-m+2) y} \in \mathcal{D}(C y l)
$$

Then

$$
\left\langle(\mu)_{\left[e^{\tau}\right]}, \psi^{*}\right\rangle=\left\langle(\mu), T_{-\tau} \psi^{*}\right\rangle=\langle\nu, S-\tau \psi\rangle=\left\langle S_{\tau} \nu, \psi\right\rangle
$$

The condition $\mu(\bullet, m) \in \mathcal{M}[\rho, \sigma]$ is also satisfied.

Exercise 4.1.6.8 Check these properties.
Now we use (4.1.6.1) to transplant the dynamical system to $U[\rho, \sigma]$. This completes a construction of a homomorphism $\left(T^{t}, M\right) \mapsto\left(T_{t}, U[\rho, \sigma]\right)$.

Let us check that it is an imbedding, i.e., we must check the one-to-one correspondence. One-to-one correspondence of $v(\bullet, m)$ and $\mu(\bullet, m)$ is known (Theorem 3.1.4.4). One-to-one correspondence of $\mu(\bullet, m)$ and $\nu(\bullet, m)$ can be also checked easily.

Exercise 4.1.6.9 Check this in detail.
So we should check the one-to-one correspondence of $\nu(\bullet, m)$ and $Y(\bullet, m)$.
Suppose

$$
\nu\left(\bullet, m_{1}\right)=\nu\left(\bullet, m_{2}\right) .
$$

Then

$$
\left\langle\nu\left(\bullet, m_{1}\right), \psi\right\rangle=\left\langle\nu\left(\bullet, m_{2}\right), \psi\right\rangle \forall \psi \in \mathcal{D}(C y l) .
$$

In particular, set

$$
\psi\left(x^{0}, y\right)=\phi\left(x^{0}\right) R(y), \phi \in \mathcal{D}\left(S^{m-1}\right), R \in \mathcal{D}(-\infty, \infty)
$$

Then

$$
\begin{align*}
\left\langle\nu\left(\bullet, m_{1}\right), \psi\right\rangle & =\int R(y) d y \int_{-\infty}^{\infty}\left\langle Y\left(\bullet, T^{y-t} m_{1}\right), \phi\right\rangle_{S^{m-1}} X(t) d t  \tag{4.1.6.6}\\
& =\left\langle\nu\left(\bullet, m_{2}\right), \psi\right\rangle=\int R(y) d y \int_{-\infty}^{\infty}\left\langle Y\left(\bullet, T^{y-t} m_{2}\right), \phi\right\rangle_{S^{m-1}}
\end{align*}
$$

where

$$
\langle Y(\bullet), \phi\rangle_{S^{m-1}}:=\int_{S^{m-1}} \phi\left(x^{0}\right) Y\left(d x^{0}\right) .
$$

Set

$$
F_{j}(y):=\left\langle Y\left(\bullet, T^{y} m_{j}\right), \phi\right\rangle_{S^{m-1}}, j=1,2 .
$$

From (4.1.6.6) we obtain for the convolutions

$$
\left(F_{1} * X\right)(y) \equiv\left(F_{2} * X\right)(y), y \in(-\infty, \infty)
$$

Thus

$$
F_{1}(y) \equiv F_{2}(y), y \in(-\infty, \infty)
$$

because of the property of $X$.
Hence

$$
Y\left(\bullet, T^{y} m_{1}\right) \equiv Y\left(\bullet, T^{y} m_{2}\right), y \in(-\infty, \infty)
$$

In particular, for $y=0$ we have

$$
Y\left(\bullet, m_{1}\right)=Y\left(\bullet, m_{2}\right) .
$$

Hence $m_{1}=m_{2}$ because of (4.1.5.3), and this completes the proof of one-to-one correspondence.

Consider the case of an integer $\rho$. For this case we can use $v \in U[\rho, \sigma]$ of the form

$$
v(x)=\Pi_{<}(x, \mu, \rho)+\Pi_{>}(x, \mu, \rho)
$$

instead of (4.1.6.1).

Exercise 4.1.6.10 Check this.
4.1.7 The most simple set satisfying the conditions of Theorem 4.1.3.3 is the set that is generated by a function $v \in U[\rho]$ that has the property

$$
v_{\left[t e^{P}\right]}=v_{[t]}, t \in(0, \infty)
$$

for some $P$.
Then

$$
T_{t+P} v=T_{t} v, t \in(-\infty, \infty)
$$

i.e., the dynamical system $T_{\bullet}$ is periodic with the period $P$ on the set

$$
\mathbb{C}(v)=\left\{T_{t} v: 0 \leq t \leq P\right\}
$$

Theorem 4.1.7.1 (Periodic Limit Set) For all $P>0, \rho>0, \sigma>0$, there exists $v \in U[\rho, \sigma]$ such that the dynamical system $\left(T_{\bullet}, \mathbb{C}(v)\right)$ is periodic with the period $P$.

Proof. Suppose $\rho$ is non-integer. Let us take $\mu \in \mathcal{M}[\rho, \Delta]$ such that the canonical potential $\Pi(x, \mu,[\rho])$ belongs to $U[\rho, \sigma]$. This is possible because of Theorem 3.1.4.2 (*Brelot-Borel).

Denote as $\mu_{P}^{*}$ the restriction of $\mu$ on the spherical ring $\left\{x: 1<|x|<e^{P}\right\}$ and set

$$
\mu_{P}:=\sum_{k=-\infty}^{\infty} T_{k P} \mu_{P}^{*}
$$

We have $\mu_{P} \in \mathcal{M}[\rho, \Delta]$ and

$$
T_{t+P} \mu_{P}=T_{t}\left(\sum_{k=-\infty}^{\infty} T_{(k+1) P} \mu_{P}^{*}\right)=T_{t} \mu_{P}
$$

Then $v:=\Pi\left(x, \mu_{P},[\rho]\right) \in U[\rho, \sigma]$ and $T_{t+P} v=T_{t} v$ because of (3.1.5.0).
For an integer $\rho$ we use the function

$$
v(x):=\Pi_{<}\left(x, \mu_{P}, \rho\right)+\Pi_{>}\left(x, \mu_{P}, \rho\right) .
$$

### 4.2 Subharmonic function with prescribed limit set

4.2.1 The following two theorems describe structure of limit sets in terms of dynamical systems.

Theorem 4.2.1.1 (Necessity) Let $u \in S H\left(\mathbb{R}^{m}, \rho, \rho(r)\right)$. Then the dynamical system $\left(T_{\bullet}, \operatorname{Fr}[u, \bullet]\right)$ is chain recurrent.

The chain recurrence is also sufficient.
Theorem 4.2.1.2 (Sufficiency) Let $U$ be a compact connected and $T_{\bullet}$-invariant subset of $U[\rho, \sigma]$ for some $\sigma>0$, such that the dynamical system $\left(T_{\bullet}, U\right)$ is chain recurrent. Then for any proximate order $\rho(r) \rightarrow \rho$ there exists $u \in S H\left(\mathbb{R}^{m}, \rho, \rho(r)\right)$ such that

$$
\operatorname{Fr}\left[u, \rho(r), V_{t}, \mathbb{R}^{m}\right]=U
$$

Proof of Theorem 4.2.1.1. We need the curve $u_{t}, t \geq 1$, and $\operatorname{Fr}[u, \bullet]$ to be contained in a common metric space $X$. Thus we set

$$
X:=\left\{v \in S H\left(\mathbb{R}^{m}\right): \sup _{r \geq 1} M(r, v) r^{-\rho-1} \leq \sup _{r \geq 1} M(r, u) r^{-\rho-1}\right\}
$$

We want to use Theorem 4.1.1.1 cr 2). Let $U$ be an open proper subset of $\operatorname{Fr}[u, \bullet]$ satisfying (4.1.1.1) and let $F$ be a $T_{\bullet}$-invariant subset of $K:=\operatorname{Fr}[u, \bullet] \backslash U$.

Such $F$ exists. Indeed, $K$ is closed and $T_{t} K \subset K$ for $t>0$ (see proof of Theorem 4.1.1.1, cr2) $\Longleftrightarrow \operatorname{cr} 3)$ ). Thus $\Omega(K) \subset K$ where $\Omega(\bullet)$ was defined in (4.1.3.2). The set $\Omega(K)$ is invariant with respect to $T_{t}$ (see Exercise 4.1.3.1). So the set of such sets $F$ is not empty.

If $F$ intersects $\partial U$ at a point $v$, then $A(v) \subset F \cap \partial U$. Since $A(v)$ is invariant (Exercise 4.1.3.2) $\partial U$ contains a nonempty $T_{\bullet}$-invariant set. So we obtain the assertion of the theorem using Theorem 4.1.1.1, cr2).

Suppose $F$ does not intersect $\partial U$. Let $U_{0}$ be an open set in $X$ such that

$$
\begin{equation*}
U_{0} \cap \mathbf{F r}[u, \bullet]=U, \operatorname{clos} U_{0} \cap \mathbf{F r}[u, \bullet]=\operatorname{clos} U \tag{4.2.1.1}
\end{equation*}
$$

(see Exercise 4.2.1.1). Since $\operatorname{clos} U_{0} \cap F=\varnothing$ we can take a sequence of open neighborhoods $U_{1}, U_{2}, \ldots$ of $F$ in $X$ such that all sets $\operatorname{clos} U_{j}, j=1,2, \ldots$ do not intersect $\operatorname{clos} U_{0}$ and $U_{j} \downarrow F$.

By definition of $\mathbf{F r}[u, \bullet]$ we can find intervals $a_{j} \leq t \leq b_{j}$ with $a_{j} \rightarrow \infty$ such that $u_{e^{a_{j}}} \in \partial U_{j}, u_{e^{b_{j}}} \in \partial U_{0}$, and $u_{e^{t}} \notin \operatorname{clos} U_{0} \cup \operatorname{clos} U_{j}$ for $a_{j}<t<b_{j}$. We can pass to a subsequence and assume that

$$
\begin{equation*}
u_{e^{a_{j}}} \rightarrow w \in F . \tag{4.2.1.2}
\end{equation*}
$$

Let us use the identity

$$
\begin{equation*}
u_{e^{t+a_{j}}}=\left(u_{e^{a_{j}}}\right)_{e^{t}} \frac{\rho\left(e^{t}\right) \rho\left(e^{a_{j}}\right)}{\rho\left(e^{t+a_{j}}\right)} . \tag{4.2.1.3}
\end{equation*}
$$

By (4.2.1.2), (4.2.1.3) and the property (3.1.2.2) of a proximate order we obtain

$$
u_{e^{t+a_{j}}} \rightarrow T_{t} w \in F
$$

uniformly for any bounded set of $t$. Thus $b_{j}-a_{j} \rightarrow \infty$.
Passing to a subsequence we may assume that $u_{e^{b_{j}}} \rightarrow v \in \operatorname{Fr}[u, \bullet] \cap \partial U_{0}=$ $\partial U$. Since $u_{e^{t+b_{j}}} \rightarrow T_{t} v$ and $u_{e^{t+b_{j}}} \notin U_{0}$ when $a_{j}-b_{j}<t<0$ we obtain that $T_{t} v \notin U$ when $t<0$.

Hence the whole backward orbit $\left\{T_{t} v: t<0\right\}$ lies in $\partial U$, which must therefore contain the $T_{\bullet}$-invariant set $A(v)$.

Exercise 4.2.1.1 Prove that the set

$$
U_{0}:=\bigcup_{v \in U}\{w \in X: \operatorname{dist}(v, w)<\operatorname{dist}(v, K) / 2\}
$$

satisfies the conditions (4.2.1.1).

Proof. We have

$$
U_{0} \supset U \Longrightarrow U_{0} \cap \mathbf{F r}[u, \bullet] \supset U \cap \mathbf{F r}[u, \bullet]=U
$$

Thus

$$
\begin{equation*}
U_{0} \cap \operatorname{Fr}[u, \bullet] \supset U \tag{4.2.1.4}
\end{equation*}
$$

From (4.2.1.4) we have

$$
\begin{equation*}
\operatorname{clos} U_{0} \cap \mathbf{F r}[u, \bullet]=\operatorname{clos} U_{0} \cap \operatorname{clos} \mathbf{F r}[u, \bullet]=\operatorname{clos}\left(U_{0} \cap \mathbf{F r}[u, \bullet]\right) \supset \operatorname{clos} U \tag{4.2.1.5}
\end{equation*}
$$

Finally $(4.2 .1 .4) \wedge(4.2 .1 .5) \Longrightarrow(4.2 .1 .1)$.
4.2.2 To prove Theorem 4.2.1.2 we need some preparation. Theorems of the next Sections form the basis of the construction that we will use in the proof.

Let $\beta$ be an infinitely differentiable function on $\mathbb{R}$ such that $0 \leq \beta(x) \leq$ $1, \beta(x)=0$ for $x \leq 0$ and $\beta(x)=1$ for $x \geq 1$. We can set, for example,

$$
\beta(x):=A \int_{-\infty}^{x} \alpha(y+1) d y
$$

where $\alpha$ is taken from (2.3.1.1) and

$$
A=\int_{-\infty}^{\infty} \alpha(y+1) d y
$$

Suppose that the sequences $\left\{r_{k}, \sigma_{k}, k=0,1, \ldots\right\}$ satisfy the following conditions:

$$
\begin{gather*}
r_{0}=1 ; r_{k}<r_{k} \sigma_{k}<r_{k+1} / \sigma_{k+1}<r_{k+1}, k=1,2, \ldots,  \tag{4.2.2.1}\\
\sigma_{k} \uparrow \infty ; \frac{r_{k+1}}{\sigma_{k+1} r_{k} \sigma_{k}} \uparrow \infty \tag{4.2.2.2}
\end{gather*}
$$

Set

$$
\begin{aligned}
& \psi_{k}(r):=\beta\left(\frac{\log r-\log \left(r_{k} / \sigma_{k}\right)}{\log \left(\sigma_{k} r_{k}\right)-\log \left(r_{k} / \sigma_{k}\right)}\right)-\beta\left(\frac{\log r-\log \left(r_{k+1} / \sigma_{k+1}\right)}{\log \left(\sigma_{k+1} r_{k+1}\right)-\log \left(r_{k+1} / \sigma_{k+1}\right)}\right) \\
& \psi_{0}(r):=1-\beta\left(\frac{\log r-\log \left(r_{1} / \sigma_{1}\right)}{\log \left(\sigma_{1} r_{1}\right)-\log \left(r_{1} / \sigma_{1}\right)}\right)
\end{aligned}
$$

The sequence $\left\{\psi_{k}\right\}, k=0,1, \ldots$ forms a partition of unity with the following properties:

Theorem 4.2.2.1 (Partition of Unity) One has

$$
\begin{gather*}
\sum_{k=0}^{\infty} \psi_{k}=1  \tag{prtu1}\\
\operatorname{supp} \psi_{k} \subset\left(r_{k} / \sigma_{k}, r_{k+1} \sigma_{k+1}\right)  \tag{prtu2}\\
\psi_{k}(r)=1, \text { for } r \in\left(r_{k} \sigma_{k}, r_{k+1} / \sigma_{k+1}\right)  \tag{prtu3}\\
\operatorname{supp} \psi_{k} \cap \operatorname{supp} \psi_{l}=\varnothing \text { for }|k-l|>1  \tag{prtu4}\\
\lim _{k \rightarrow \infty} \max _{r} \psi_{k}^{\prime}(r) r=\lim _{k \rightarrow \infty} \max _{r} \psi_{k}^{\prime \prime}(r) r^{2}=0 \tag{prtu5}
\end{gather*}
$$

Moreover

$$
\begin{equation*}
\max _{r}\left|\psi_{k}^{\prime}(r) r\right|, \max _{r}\left|\psi_{k}^{\prime \prime}(r) r^{2}\right| \leq \gamma_{k} \tag{prtu6}
\end{equation*}
$$

where $\gamma_{k}$ can be made to tend to zero arbitrarily fast by choosing the sequences $\left\{\sigma_{k}\right\}$ and $\left\{r_{k}\right\}$.

Proof. Set

$$
\beta_{k}(r):=\beta\left(\frac{\log r-\log \left(r_{k} / \sigma_{k}\right)}{\log \left(\sigma_{k} r_{k}\right)-\log \left(r_{k} / \sigma_{k}\right)}\right)
$$

The functions $\beta_{k}(r)$ and $\beta_{k+1}(r)$ vanish for $r<r_{k} / \sigma_{k}$ because $\beta(x)=0$ for $x \leq 0$, and both of them are equal to 1 for $r \geq \sigma_{k} r_{k}$ because $\beta(x)=1$ for $x \geq 1$. Hence, (prtu2) holds.

One has for any $r \in(0, \infty)$,

$$
\sum_{k=0}^{n} \psi_{k}=1-\beta_{n+1}(r)
$$

As mentioned, $\beta_{n+1}(r)=0$ for $n$ such that $r_{n+1} / \sigma_{n+1}>r$. Thus (prtu1) holds.

Counting derivatives of $\psi_{k}$, we have:

$$
\begin{gathered}
\max _{r}\left|r \psi_{k}^{\prime}(r)\right| \leq \\
{\left[\left(\log \left(\sigma_{k} r_{k}\right)-\log \left(r_{k} / \sigma_{k}\right)\right)^{-1}+\left(\log \left(\sigma_{k+1} r_{k+1}\right)-\log \left(r_{k+1} / \sigma_{k+1}\right)\right)^{-1}\right] \max _{x}\left|\beta^{\prime}\right|(x)}
\end{gathered}
$$

Thus we can take the right side of the inequality as $\gamma_{k}$ and regulate its vanishing by choice of the ratio in (4.2.2.2). The same holds for $r^{2} \psi^{\prime \prime}(r)$. Hence (prtu5) and (prtu6) are proved.

Exercise 4.2.2.1 Check (prtu4).
4.2.3 Now we construct a function which is of zero type but has a "maximal possible" mass density.

Theorem 4.2.3.1 (Maximal Mass Density Function) Let $\rho(r) \rightarrow \rho, \rho>0$ be $a$ smooth proximate order (i.e., having properties (2.8.1.8)), and let $\gamma(r), r \in[0, \infty)$, satisfy the conditions: $\gamma(r)>0$ and $\gamma(r) \rightarrow 0$, as $r \rightarrow \infty$.

Then there exists an infinitely differentiable subharmonic function $\Phi(x)$ such that

$$
\begin{equation*}
\Delta \Phi(x) \geq \gamma(x)|x|^{\rho(r)-2} \tag{4.2.3.1}
\end{equation*}
$$

and

$$
\begin{equation*}
(\Phi)_{t} \rightarrow 0 \tag{4.2.3.2}
\end{equation*}
$$

in $\mathcal{D}^{\prime}$ as $t \rightarrow \infty$.
To prove Theorem 4.2.3.1 we need an elementary lemma.
Theorem 4.2.3.2 (Convex Majorization) Let $a(s), s \in\left[s_{0}, \infty\right)$ be a function such that $a(s) \rightarrow-\infty$ as $s \rightarrow \infty$. Then there exists an infinitely differentiable, convex function $k(s)$ such that:
k1) $\quad k(s) \geq a(s)$;
k2) $\quad k(s) \downarrow-\infty$ as $s \rightarrow \infty$;
$\mathrm{k} 3) \quad k^{(n)}(s) \rightarrow 0$ for all $n=1,2, \ldots$.
Proof. Set

$$
a^{*}(s):=\sup \{a(t): t \geq s\}
$$

Then $a^{*}(s) \downarrow-\infty$ as $s \rightarrow \infty$.
Set $b_{0}:=-a^{*}\left(s_{0}\right)$ and denote as $s(b), b \in\left[b_{0},+\infty\right)$ the function inverse to the function $-a^{*}(s)$. Let us construct a convex function that majorates $s(b)$ and tends to infinity monotonically with all its derivatives. It can be done in the following way.

First we construct a piecewise linear convex function. Set

$$
s_{1}(b):=s_{0}+1+\alpha_{0}\left(b-b_{0}\right), b \in\left[b_{0}, b_{0}+1\right],
$$

and choose $\alpha_{0}$ such that the inequality $s_{1}(b)>s(b)$ holds for $b \in\left[b_{0}, b_{0}+1\right]$.
For this we choose

$$
\alpha_{0} \geq \sup _{b \in\left[b_{0}, b_{0}+1\right]} \frac{s(b)-s_{0}-1}{b-b_{0}}
$$

Since $s(b)-s_{0}-1<0$ the right side is finite.
For all the following intervals we set

$$
s_{1}(b):=s_{1}\left(b_{0}+j\right)+\alpha_{j}\left(b-b_{0}-j\right), b \in\left[b_{0}+j, b_{0}+j+1\right],
$$

where $\alpha_{j} \geq \alpha_{j-1}$ and satisfies the condition

$$
\alpha_{j} \geq \sup _{b \in\left[b_{0}+j, b_{0}+j+1\right]} \frac{s_{1}(b)-s_{1}\left(b_{0}+j\right)}{b-b_{0}-j}
$$

To obtain a smooth function, set

$$
\begin{equation*}
s_{2}(b):=\int \alpha(b-x) s_{1}(x) d x \tag{4.2.3.3}
\end{equation*}
$$

where $\alpha(x)$ is defined by (2.3.1.1). Then $s_{2}(b)$ is infinitely differentiable, monotonic and convex.

## Exercise 4.2.3.1 Check this.

Set

$$
\begin{equation*}
k(s):=-s_{2}^{-1}(s) \tag{4.2.3.4}
\end{equation*}
$$

where $s_{2}^{-1}(s)$ is the inverse function to $s_{2}$. One can check that $k(s)$ satisfies the properties k1), k2), k3).

Exercise 4.2.3.2 Check that $k(s)$ satisfies k 1$), \mathrm{k} 2), \mathrm{k} 3)$.
Proof of Theorem 4.2.3.1. We are going to show that $\Phi$ can be taken in the form

$$
\begin{equation*}
\Phi(x):=c e^{k\left(\log |x|^{2}\right)}|x|^{\rho(|x|)} \tag{4.2.3.5}
\end{equation*}
$$

where $c$ and $k(s)$ will be chosen later.
Note that $\Phi(x)=\Phi(|x|)$ depends only on $r=|x|$ and pass to the variable $s:=\log r^{2}$. Then for $\phi(s):=\Phi\left(e^{s / 2}\right)$ we have

$$
\begin{align*}
\Delta \Phi(x) & =r^{1-m} \frac{\partial}{\partial r} r^{m-1} \frac{\partial}{\partial r} c e^{k\left(\log r^{2}\right)} r^{\rho(r)} \\
& =c e^{-s}\left(\frac{\partial^{2}}{\partial s^{2}}+\frac{m-2}{2} \frac{\partial}{\partial s}\right) \phi(s) \geq c m e^{-s} \min \left[\phi^{\prime \prime}(s), \phi^{\prime}(s)\right] \tag{4.2.3.6}
\end{align*}
$$

Let us chose $k$ as in Theorem 4.2.3.2 with $a(s):=\log \gamma(r)=\log \gamma\left(e^{\frac{s}{2}}\right)$. Now we estimate the derivatives from below.

$$
\phi^{\prime}(s)=\phi(s)\left[k^{\prime}(s)+\frac{1}{4} s \rho^{\prime}\left(e^{\frac{s}{2}}\right)+\frac{1}{2} \rho\left(e^{\frac{s}{2}}\right)\right] .
$$

By k3) and k1), $k^{\prime}(s) \rightarrow 0$ and $k(s) \geq a(s)$. Also $s \rho^{\prime}\left(e^{\frac{s}{2}}\right) \rightarrow 0$ and $\rho\left(e^{\frac{s}{2}}\right) \rightarrow \rho$ by properties of proximate order (Theorem 2.8.1.4). Thus we can chose $c$ such that

$$
\begin{equation*}
\phi^{\prime}(s)>\frac{1}{m} e^{\log \gamma\left(e^{\frac{s}{2}}\right)+\frac{s}{2} \rho\left(e^{\frac{s}{2}}\right)} . \tag{4.2.3.7}
\end{equation*}
$$

Differentiating once again, we obtain

$$
\phi^{\prime \prime}(s)=\phi(s)\left[k^{\prime}(s)+\frac{1}{4} s \rho^{\prime}\left(e^{\frac{s}{2}}\right)+\frac{1}{2} \rho\left(e^{\frac{s}{2}}\right)\right]^{2}+\left[k^{\prime \prime}(s)+\frac{1}{2} \rho^{\prime}\left(e^{\frac{s}{2}}\right)+\frac{1}{8} s \rho^{\prime \prime}\left(e^{\frac{s}{2}}\right)\right] .
$$

From here we obtain by choosing $c$ :

$$
\begin{equation*}
\phi^{\prime \prime}(s)>\frac{1}{m} e^{\log \gamma\left(e^{\frac{s}{2}}\right)+\frac{s}{2} \rho\left(e^{\frac{s}{2}}\right)} . \tag{4.2.3.8}
\end{equation*}
$$

Using (4.2.3.6), (4.2.3.7) and (4.2.3.8) we obtain:

$$
\Delta \Phi(s)>e^{\log \gamma\left(e^{\frac{s}{2}}\right)+\frac{s}{2} \rho\left(e^{\frac{s}{2}}\right)}
$$

Returning to the variable $r$ we obtain (4.2.3.1). Correctness of (4.2.3.2) can be checked directly using k 2 ) and properties of the proximate order (Theorem 2.8.1.3).

Exercise 4.2.3.3 Check this.
4.2.4 We have already approximated distributions and subharmonic functions by infinitely differentiable functions (Theorems 2.3.4.5 and 2.6.2.3). Now we need to make more precise this approximation. Namely, we are going to make it uniform with respect to $v \in U[\rho, \sigma]$. We will denote

$$
\begin{equation*}
\partial^{l}:=\frac{\partial^{|l|}}{\left(\partial x_{1}\right)^{l_{1}}\left(\partial x_{2}\right)^{l_{2}} \cdots\left(\partial x_{m}\right)^{l_{m}}} \tag{4.2.4.1}
\end{equation*}
$$

where $l=\left(l_{1}, l_{2}, \ldots, l_{m}\right),|l|=l_{1}+l_{2}+\cdots+l_{m}$.
Set for $v \in U[\rho, \sigma]$

$$
\begin{equation*}
R_{\epsilon} v(x):=\int \alpha_{\epsilon}(x-y) v(y) d y \tag{4.2.4.2}
\end{equation*}
$$

where $\alpha_{\epsilon}$ is taken from (2.3.1.3).
We have changed the notation from 2.3.1 and 2.6.2 because a subindex of $v$ was already engaged for $t$.

For a fixed $0<\delta \leq 0.5$, set

$$
\begin{equation*}
\operatorname{Str}(\delta):=\left\{x: \delta \leq|x| \leq \delta^{-1}\right\} \tag{4.2.4.3}
\end{equation*}
$$

Theorem 4.2.4.1 (Estimation of $R_{\epsilon}$ ) Let $v \in U[\rho, \sigma]$. Then
R1. for a fixed $g \in \mathcal{D}\left(\mathbb{R}^{m} \backslash 0\right)$ with $\operatorname{supp} g \subset \operatorname{Str}(\delta)$,

$$
\begin{equation*}
\left|\left\langle R_{\epsilon} v-v, g\right\rangle\right| \leq o(1, g) 2 \sigma \delta^{-\rho} \tag{4.2.4.4}
\end{equation*}
$$

where $o(1, g) \rightarrow 0$ as $\epsilon \rightarrow 0 ;$
R 2 . the inequality

$$
\begin{equation*}
\left|\partial^{l} R_{\epsilon} v(x)\right| \leq A(m) \sigma \epsilon^{-|l|-m+1}|x|^{-|l|+\rho} \tag{4.2.4.5}
\end{equation*}
$$

with $A(m)$ depending only on $m$, holds for $\epsilon<|x| / 2$.
Proof. One has

$$
\begin{equation*}
\left\langle R_{\epsilon} v, g\right\rangle=\left\langle v, R_{\epsilon} g\right\rangle \tag{4.2.4.6}
\end{equation*}
$$

Thus

$$
\begin{equation*}
\left\langle R_{\epsilon} v-v, g\right\rangle=\left\langle v, R_{\epsilon} g-g\right\rangle \tag{4.2.4.7}
\end{equation*}
$$

Exercise 4.2.4.1 Check (4.2.4.6) and (4.2.4.7).
Now

$$
\begin{equation*}
\left|\left\langle v, R_{\epsilon} g-g\right\rangle\right| \leq \max _{\operatorname{Str}(\delta)}\left|R_{\epsilon} g-g\right|(x) \int_{\operatorname{Str}(\delta)}|v|(x) d x \tag{4.2.4.8}
\end{equation*}
$$

The first factor is $o(1)$ because $g$ is smooth. For the second one we have

$$
\begin{equation*}
\int_{\operatorname{Str}(\delta)}|v|(x) d x \leq 2 \int_{\operatorname{Str}(\delta)} v^{+}(x) d x \leq 2 \sigma \delta^{-\rho} \tag{4.2.4.9}
\end{equation*}
$$

This and (4.2.4.8) imply R1).
Differentiating the equality

$$
R_{\epsilon} v(x):=C_{m} \int \epsilon^{-m} \alpha(|x-y| / \epsilon) v(y) d y
$$

we have

$$
\left|\partial^{l} R_{\epsilon} v(x)\right| \leq C_{m} \epsilon^{-|l|-m} \max _{\{|y|<\epsilon\}}\left|\partial^{l} \alpha(|y|)\right| \int_{\{|y|<\epsilon\}}|v(x-y)| d y
$$

Suppose $|x|=1$. Then for $0<\epsilon \leq 0.5$, we have

$$
\begin{aligned}
\int_{|y|<\epsilon}|v|(x-y) d y & \leq \int_{1-\epsilon<|x|<1+\epsilon}|v|(x) \leq 2 \int_{1-\epsilon<|x|<1+\epsilon} v^{+}(x) d x \\
& \leq \sigma_{m} 2 \cdot 2 \epsilon \sigma(1+\epsilon)^{\rho} \leq \sigma_{m} 6 \sigma \epsilon
\end{aligned}
$$

where $\sigma_{m}$ is the square of the unit sphere. Hence for $|x|=1$

$$
\begin{equation*}
\left|\partial^{l} R_{\epsilon} v(x)\right| \leq A(m) \sigma \epsilon^{-|l|-m+1} \tag{4.2.4.10}
\end{equation*}
$$

with

$$
A(m)=6 \sigma_{m} \max _{y \in \mathbb{R}^{m}}\left|\partial^{l} \alpha(|y|)\right| .
$$

Set $t=|x|$. Apply the inequality (4.2.4.10) to $v:=v_{[t]}(y)$ with $y:=x /|x|$. Then

$$
\left|\partial^{l} R_{\epsilon} v_{[t]}(y)\right| \leq A(m) \sigma \epsilon^{-|l|-m+1}
$$

Computing the derivatives, we obtain

$$
\partial^{l} R_{\epsilon} v_{[t]}(x)=\left.t^{-\rho} t^{|l|} \partial^{l} R_{\epsilon} v(x)\right|_{x=t y}
$$

Thus one has R2.
4.2.5 In this section we describe the main part of a construction that will be used in the proof of Theorem 4.2.1.2.

Let $\left\{v_{j} \in U[\rho, \sigma], j=1,2, \ldots\right\}$ and $\left\{\psi_{j}, j=1,2 \ldots\right\}$ be the partition of unity from Theorem 4.2.2.1. Let us chose $\epsilon_{j} \downarrow 0$ such that the condition

$$
\begin{equation*}
\gamma_{j} \epsilon_{j}^{-m} \rightarrow \infty \tag{4.2.5.1}
\end{equation*}
$$

holds for $\gamma_{j}$ taken from Theorem 4.2.2.1, (prtu 6). Set

$$
\begin{equation*}
v(x \mid t):=\sum_{j=0}^{\infty} \psi_{j}(t)\left(v_{j}\right)_{[t]}(x), \tag{4.2.5.2}
\end{equation*}
$$

where $(\cdot)_{[t]}$ is defined by (3.1.2.4a).
One can see that $v(x \mid t) \in U[\rho, 3 \sigma]$ for all $t$.
Exercise 4.2.5.1 Show this, using properties of $\left\{\psi_{j}\right\}$ and invariance of $U[\rho, \sigma]$ with respect to $(\cdot)_{[t]}$.

We can consider $v(x \mid t)$ as a curve (a pseudo-trajectory) in $U[\rho, 3 \sigma]$.
Set

$$
\begin{equation*}
u(x):=\sum_{j=0}^{\infty} \psi_{j}(|x|) R_{\epsilon_{j}}\left(v_{j}\right)(x)|x|^{\rho(|x|)-\rho} . \tag{4.2.5.3}
\end{equation*}
$$

where $R_{\epsilon}$ is defined by (4.2.4.2).
It is an infinitely differentiable function in $\mathbb{R}^{m}$.
Theorem 4.2.5.1 (Construction) One has

$$
\begin{equation*}
u_{t}-v(\bullet \mid t) \rightarrow 0 \tag{4.2.5.4}
\end{equation*}
$$

in $\mathcal{D}^{\prime}\left(\mathbb{R}^{m}\right)$, and

$$
\begin{equation*}
\Delta u(x)=f(x)+\gamma(x)|x|^{\rho(|x|)-2} \tag{4.2.5.5}
\end{equation*}
$$

with $f(x) \geq 0$ and $\gamma(x)=o(1)$ as $|x| \rightarrow \infty$.

Let us note that the function $u(x)$ is "almost-subharmonic" and can be made subharmonic by summing with the function $\Phi$ from Theorem 4.2.3.1.

Exercise 4.2.5.2 Prove this.
So we have
Theorem 4.2.5.2 (Pseudo-Trajectory Asymptotics) For any $v(x \mid t)$ of the form (4.2.5.2) there exists an infinitely differentiable function $u \in S H(\rho(r))$ that satisfies (4.2.5.4).

Proof of Theorem 4.2.5.1. One has

$$
u_{t}(x):=\sum_{j=0}^{\infty} \psi_{j}(t|x|)\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}(x) a(x, t)
$$

where

$$
a(x, t):=\frac{|t x|^{\rho(t|x|)-\rho}}{t^{\rho(t)-\rho}} .
$$

For any $0<\delta<0.5$ and $x \in \operatorname{Str}(\delta), a(x, t) \rightarrow 1$ uniformly in $|x|$ as $t \rightarrow \infty$. This follows from Theorem 2.8.1.3, ppo3).

Exercise 4.2.5.3 Check this in detail.
We have

$$
\begin{equation*}
u_{t}(x)-v(x \mid t)=\sum_{j=0}^{\infty}\left[\psi_{j}(t|x|)\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}(x) a(x, t)-\psi_{j}(t)\left(v_{j}\right)_{[t]}(x)\right] \tag{4.2.5.6}
\end{equation*}
$$

and there are no more than three summands in the sum for sufficiently large $t=t(\delta)$ because of Theorem 4.2.2.1, prtu4. Let us estimate every summand. One has

$$
\begin{aligned}
b_{j}(x, t): & =\left[\psi_{j}(t|x|)\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}(x) a(x, t)-\psi_{j}(t)\left(v_{j}\right)_{[t]}(x)\right] \\
= & {\left.\left[\psi_{j}(t|x|)-\psi_{j}(t)\right]\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)\right)_{[t]} a(x, t)+\psi_{j}(t)\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}(x)[a(x, t)-1] } \\
& +\psi_{j}(t)\left[\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}(x)-\left(v_{j}\right)_{[t]}(x)\right] \\
:= & \left(a_{1}+a_{2}+a_{3}\right)(x, t) .
\end{aligned}
$$

Let us estimate $\left\langle b_{j}(\bullet, t), g\right\rangle$ for every $g \in \mathcal{D}\left(\mathbb{R}^{m} \backslash 0\right)$.
We can assume that supp $g \subset \operatorname{Str}(\delta)$. Set

$$
M(g):=\max _{x \in \operatorname{Str}(\delta)}|g|(x)
$$

We have

$$
\left|\left\langle a_{1}(\bullet, t), g\right\rangle\right| \leq M(g) \max _{r \in(0, \infty)}\left|r \psi_{j}^{\prime}(r)\right| \delta^{-1} \int_{\operatorname{Str}(\delta)}\left|\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}\right|(x) d x
$$

One can check that

$$
\int_{\operatorname{Str}(\delta)}\left|\left(R_{\epsilon_{j}}\left(v_{j}\right)\right)_{[t]}\right|(x) d x \leq 3 \sigma \delta^{-\rho}
$$

Exercise 4.2.5.4 Check this using (4.2.4.9) and the invariance of $U[\rho, \sigma]$ with respect to $(\bullet)_{[t]}($ see (3.1.2.4)).

Hence

$$
\begin{equation*}
\left|\left\langle a_{1}(\bullet, t), g\right\rangle\right| \leq C_{1}(g) \gamma_{j} \tag{4.2.5.7}
\end{equation*}
$$

Let us estimate $a_{2}(x, t)$. We have

$$
\begin{equation*}
\left\langle a_{2}(\bullet, t), g\right\rangle \leq \max _{\operatorname{Str}(\delta)}|a(x, t)-1| \psi_{j}(t) M(g) 3 \sigma \delta^{-\rho}=C_{2}(g) o(1) \tag{4.2.5.8}
\end{equation*}
$$

where $o(1) \rightarrow 0$ as $t \rightarrow \infty$.
For estimating $a_{3}(x, t)$, we use Theorem 4.2.4.1 (Estimation of $R_{\epsilon}$ ), (4.2.4.4):

$$
\begin{equation*}
\left|\left\langle a_{3}(\bullet, t), g\right\rangle\right| \leq o\left(\epsilon_{j}, g\right) 2 \sigma \delta^{-\rho} \tag{4.2.5.9}
\end{equation*}
$$

where $o\left(\epsilon_{j}, g\right) \rightarrow 0$ as $j \rightarrow \infty$.
Hence (4.2.5.7), (4.2.5.8) and (4.5.5.9) imply

$$
\begin{equation*}
\left\langle b_{j}(\bullet, t), g\right\rangle \rightarrow 0 \tag{4.2.5.10}
\end{equation*}
$$

as $t \rightarrow \infty$ and $j \rightarrow \infty$.
Suppose, for a large fixed $t$, the sum (4.2.5.6) contains $b_{j}(x, t)$ for $j=j(t), j=$ $j(t)+1$ and $j=j(t)+2$. This implies that $j(t) \rightarrow \infty$ as $t \rightarrow \infty$.

Since

$$
\left\langle u_{t}(\bullet)-v(\bullet \mid t), g\right\rangle=\left\langle b_{j(t)}(\bullet, t), g\right\rangle+\left\langle b_{j(t)+1}(\bullet, t), g\right\rangle+\left\langle b_{j(t)+2}(\bullet, t), g\right\rangle
$$

we obtain from (4.2.5.10) that $\left\langle u_{t}(\bullet)-v(\bullet \mid t), g\right\rangle \rightarrow 0$ as $t \rightarrow \infty$ for any $g \in$ $\mathcal{D}\left(\mathbb{R}^{m} \backslash 0\right)$. This is (4.2.5.4).

Let us prove (4.2.5.5). We have

$$
\begin{equation*}
\Delta u=\sum_{j=0}^{\infty}\left[\Delta\left(R_{\epsilon_{j}} v_{j}\right)(x) \psi_{j}(x)|x|^{\rho(|x|)-\rho}+\sum_{l, m \cdot k} \partial^{l}\left(R_{\epsilon_{j}} v_{j}\right)(x) \partial^{n} \psi_{j}(x) \partial^{k}|x|^{(\rho|x|)-\rho}\right] \tag{4.2.5.11}
\end{equation*}
$$

where $l, m, k$ are multi-indexes that satisfy the condition: in any summand there are derivatives in the same variable, the derivatives of $\psi_{j}$ and $|x|^{\rho(|x|)-\rho}$ have no more than second order and the derivatives of $R_{\epsilon_{j}} v_{j}(x)$ have no more than first order.

Exercise 4.2.5.5 Check this.

As usual, the derivative of zero order is the function itself. For any $x \in$ $\operatorname{Str}(\delta)$, the outside sum contains no more then three summands. First we consider only the terms in the square brackets. The first term is nonnegative because of subharmonicity of $R_{\epsilon_{j}} v_{j}$ and non-negativity of all the other factors. Set

$$
\begin{equation*}
f(x):=\sum_{j=0^{\infty}}\left[\Delta\left(R_{\epsilon_{j}} v_{j}\right)(x) \psi_{j}(x)|x|^{\rho(|x|)-\rho} \geq 0\right. \tag{4.5.2.12}
\end{equation*}
$$

Using Theorem 4.2.4.1, R2) we obtain

$$
\begin{equation*}
\left|\partial^{l}\left(R_{\epsilon_{j}} v_{j}\right)(x)\right| \leq A(m) \sigma \epsilon_{j}^{-|l|-m+1}|x|^{-|l|+\rho} \tag{4.2.5.13}
\end{equation*}
$$

for $|l|=0$ or $|l|=1$.
From Theorem 4.2.2.1, prtu6), and inequality $\left|\partial_{x_{i}}\right| x|\mid \leq 1$ we obtain

$$
\begin{equation*}
\left|\partial^{n} \psi_{j}(|x|) \leq\left|\psi^{(n)}(r) \|_{r=|x|} \leq \gamma_{j}\right| x\right|^{-|n|} \tag{4.2.5.14}
\end{equation*}
$$

for $|n|=1,2$.
Using properties of the smooth proximate order (Theorem 2.8.1.4), one can obtain

$$
\begin{equation*}
\left.\left|\partial^{|k|}\right| x\right|^{\rho(|x|)-\rho}\left|=\left(|x|^{\rho(|x|)-\rho-|k|}\right)\right|_{r=|x|}(1+o(1)), \tag{4.2.5.15}
\end{equation*}
$$

as $|x| \rightarrow \infty$.
Exercise 4.2.5.6 Check in detail (4.2.5.13), (4.2.5.14) and (4.2.5.15).
Thus, for every term of the inner sum, we have

$$
\begin{align*}
& \left.\left|\partial^{l}\left(R_{\epsilon_{j}} v_{j}\right)(x) \partial^{n} \psi_{j}(x) \partial^{k}\right| x\right|^{\rho(|x|)-\rho} \mid \\
& \quad \leq A(m) \sigma \gamma_{j} \epsilon_{j}^{-|l|-m+1}|x|^{-2+\rho}|x|^{\rho(|x|)-\rho} \\
& \quad \leq \beta_{j}|x|^{\rho(|x|)-2} \tag{4.2.5.16}
\end{align*}
$$

where $\beta_{j} \rightarrow 0$ because of the condition (4.2.5.1).
Recall that for every large $x$ the outside sum contains no more than three summands, say, $j=j(x), j=j(x)+1$ and $j=j(x)+2$. Thus $j(x) \rightarrow \infty$ as $|x| \rightarrow \infty$. Hence (4.2.5.12) and (4.2.5.16) imply (4.2.5.5).

### 4.2.6

Proof of Theorem 4.2.1.2. Let $v(\bullet \mid t)$ have the form (4.2.5.2). We denote as $\Omega(v)$ a set of the $\mathcal{D}^{\prime}$-limits of the form

$$
w:=\lim _{t_{k} \rightarrow \infty} v\left(\bullet \mid t_{k}\right)
$$

We are going to construct some $v(\bullet \mid t)$ for which

$$
\begin{equation*}
\Omega(v)=U \tag{4.2.6.1}
\end{equation*}
$$

and at the next step to use Theorem 4.2.5.2 to obtain a subharmonic function with the same limit set.

First we describe the construction of the function $v(\bullet \mid t)$. Let $\left\{r_{k}, t_{k}, k=\right.$ $1,2, \ldots\}$ be an alternating sequence $r_{0}=1, r_{k}<t_{k}<r_{k+1}$ such that

$$
\lim _{k \rightarrow \infty} \frac{t_{k}}{r_{k}}=\lim _{k \rightarrow \infty} \frac{r_{k+1}}{t_{k}}=\infty
$$

Let us chose in $U$ a countable, dense set $\left\{g_{j}\right\}$ and form from it a sequence $\left\{w_{k}\right\}$ such that every element $g_{j}$ is repeated infinitely often. For example,

$$
w_{1}:=g_{1}, w_{2}:=g_{1}, w_{3}:=g_{2}, w_{4}:=g_{1}, w_{5}:=g_{2}, w_{6}:=g_{3}, \ldots
$$

Set

$$
q_{k}:=\left(w_{k}\right)_{\left[1 / t_{k}\right]}=T_{-\log t_{k}} w_{k}
$$

in the notation (4.1.5.1).
Now we use that $\left(T_{\bullet}, U\right)$ is chain recurrent. Set

$$
\alpha_{k}:=\log \frac{r_{k}}{t_{k}} ; \quad \omega_{k}:=\log \frac{r_{k+1}}{t_{k}}
$$

and find, by Theorem 4.1.2.1, a sequence $\left\{v_{j}\right\} \supset\left\{q_{k}\right\}$ such that the condition (4.1.2.1) holds, i.e.,

$$
\begin{equation*}
T_{\omega_{k}} v_{k}-T_{\alpha_{k+1}} v_{k+1} \rightarrow 0 \tag{4.2.6.2}
\end{equation*}
$$

as $k \rightarrow \infty$.
Set in Theorem 4.1.2.3

$$
p_{k}:=T_{\alpha_{k+1}} v_{k+1}, q_{k}:=T_{\omega_{k}} v_{k}
$$

and find $\gamma_{k}$ such that the condition

$$
\begin{equation*}
T_{\tau} \circ T_{\omega_{k}} v_{k}-T_{\tau} \circ T_{\alpha_{k+1}} v_{k+1} \rightarrow 0 \tag{4.2.6.3}
\end{equation*}
$$

holds uniformly for $\tau \in\left[-\gamma_{k+1}, \gamma_{k}\right]$.
Set

$$
\sigma_{k}:=\min \left[e^{\gamma_{k}}, \sqrt{\frac{t_{k}}{r_{k}}}\right] .
$$

These $\sigma_{k}$ satisfy the conditions (4.2.2.1) and (4.2.2.2).
Exercise 4.2.6.1 Check this.
We define $v(\bullet \mid t)$ by (4.2.5.2) with described $v_{j}$ and with $\psi_{j}$ from Theorem 4.2.2.1, corresponding to the chosen $r_{j}$ and $\sigma_{j}$. Let us prove (4.2.6.1).

Consider for fixed $k$ the following three cases.

1. $t \in\left[r_{k} \sigma_{k}, r_{k+1} / \sigma_{k+1}\right)$;
2. $t \in\left[r_{k+1} / \sigma_{k+1}, r_{k+1}\right)$;
3. $t \in\left[r_{k}, r_{k} \sigma_{k}\right)$.

For the first case we have

$$
v(\bullet \mid t)=\left(v_{k}\right)_{\left[t / t_{k}\right]}=T_{\log \left(t / t_{k}\right)} v_{k}
$$

For the second one

$$
\begin{aligned}
v(\bullet \mid t) & =\psi_{k}(t)\left(v_{k}\right)_{\left[t / t_{k}\right]}+\psi_{k+1}(t)\left(v_{k+1}\right)_{\left[t / t_{k+1}\right]} \\
& =\left(v_{k}\right)_{\left[t / t_{k}\right]}+\psi_{k+1}(t)\left[\left(v_{k+1}\right)_{\left[t / t_{k+1}\right]}-\left(v_{k}\right)_{\left[t / t_{k}\right]}\right.
\end{aligned}
$$

We transform the expression in the square brackets

$$
\begin{aligned}
\left(v_{k+1}\right)_{\left[t / t_{k+1}\right]} & =T_{\log \left(t / t_{k+1}\right)} v_{k+1}=T_{\log \left(t / r_{k+1}\right)} \circ T_{\log \left(r_{k+1} / t_{k+1}\right)} v_{k+1} \\
& =T_{\log \left(t / r_{k+1}\right)} \circ T_{\alpha_{k+1}} v_{k+1} .
\end{aligned}
$$

For the second term, we obtain

$$
\left(v_{k}\right)_{\left[t / t_{k}\right]}=T_{\log \left(t / r_{k+1}\right)} \circ T_{\omega_{k}} v_{k}
$$

Exercise 4.2.6.2 Check this.
Setting $\tau:=\log \left(t / r_{k+1}\right)$, we have

$$
\begin{equation*}
v(\bullet \mid t)=\left(v_{k}\right)_{\left[t / t_{k}\right]}+\psi_{k+1}(t)\left[T_{\tau} \circ T_{\alpha_{k+1}} v_{k+1}-T_{\tau} \circ T_{\omega_{k}} v_{k}\right] \tag{4.2.6.4}
\end{equation*}
$$

where $\tau \in\left[-\log \sigma_{k+1}, 0\right) \subset\left[-\gamma_{k+1}, \gamma_{k}\right]$. For the third case, set $\tau:=\log \left(t / r_{k}\right)$. Then

$$
\begin{equation*}
v(\bullet \mid t)=\left(v_{k}\right)_{\left[t / t_{k}\right]}+\psi_{k}(t)\left[T_{\tau} \circ T_{\omega_{k-1}} v_{k-1}-T_{\tau} \circ T_{\alpha_{k}} v_{k}\right] \tag{4.2.6.5}
\end{equation*}
$$

where $\tau \in\left[0, \log \sigma_{k}\right) \subset\left[-\gamma_{k}, \gamma_{k-1}\right]$.
Let $t_{N} \rightarrow \infty$ be an arbitrary sequence. Choosing a subsequence, we may suppose that there exist the limits $\left(v_{k\left(t_{N}\right)}\right)_{\left[t_{N} / t_{k\left(t_{N}\right)}\right]} \rightarrow v^{*} \in U$ and $v\left(\bullet \mid t_{N}\right) \rightarrow v_{\infty}$.

Choosing a subsequence, we may suppose that $t_{N}$ satisfies either 1 or 2 or 3 For case 1 , we obtain at once $v_{\infty}=v^{*} \in U$.

For case 2, from (4.2.6.4), (4.2.6.2) and Theorem 4.1.1.3 we obtain that the superfluous addends tend to zero, and hence $v_{\infty} \in U$.

The same holds for case 3 . Hence $\Omega(v) \subset U$.
Further, for $t=t_{k}$, we have $v(\bullet \mid t)=w_{k}$. The sequence $\left\{w_{k}\right\}$ contains the set $\left\{g_{j}\right\}$ that is dense in $U$. Thus $\Omega(v) \supset U$. Thus equality (4.2.6.1) has been proved.

As already said, the application of Theorem 4.2.5.2 concludes the proof.

### 4.3 Further properties of limit sets

4.3.1 Let as mark the following property of the pseudo-trajectory $v(\bullet \mid t)$ defined in (4.2.5.2):

Theorem 4.3.1.1 One has

$$
\begin{equation*}
T_{\tau} v\left(\bullet \mid e^{t}\right)-v\left(\bullet \mid e^{t+\tau}\right) \rightarrow 0 \tag{4.3.1.1}
\end{equation*}
$$

as $t \rightarrow \infty$ uniformly with respect to $\tau \in[a, b]$ for any $[a, b] \subset(-\infty, \infty)$.
Proof. Using the definition of $(\bullet)_{t}$ (see (3.1.2.1)) and (4.2.5.4), the remainder in (4.3.1.1) can be represented in the form

$$
b(t, \tau, \bullet):=T_{\tau} v\left(\bullet \mid e^{t}\right)-v\left(\bullet \mid e^{t+\tau}\right)=T_{\tau}\left(u_{e^{t}}\right)-u_{e^{t+\tau}}+o(1)
$$

where $o(1) \rightarrow 0$ uniformly with respect to $\tau \in[a, b]$ for any $[a, b] \subset(-\infty, \infty)$.
Exercise 4.3.1.1 Check this in detail.
Then we obtain

$$
b(t, \tau, \bullet)=u_{e^{t+\tau}}\left[e^{\rho\left(e^{t}\right)-\rho\left(e^{t+\tau}\right)}-1\right]+o(1) \rightarrow 0
$$

uniformly in the same sense due to precompactness of the family $\left\{u_{e^{t}}\right\}$ and properties of the proximate order.

Exercise 4.3.1.2 Check this in detail.
The property (4.3.1.1) shows that the pseudo-trajectory $v(\bullet \mid t)$ behaves asymptotically like the dynamical system $T_{0}$. Thus a pseudo-trajectory with this property is called an asymptotically dynamical pseudo-trajectory with dynamical asymptotics $T_{\bullet}$ (a.d.p.t.).

Theorem 4.2.5.1 shows that for any a.d.p.t. of the form (4.2.5.2) there exists $u \in S H(\rho(r))$ that satisfies the condition

$$
\begin{equation*}
u_{e^{t}}-v\left(\bullet \mid e^{t}\right) \rightarrow 0 \tag{4.3.1.2}
\end{equation*}
$$

as $t \rightarrow \infty$.
The following assertion shows that we can suppose $v(\bullet \mid \bullet)$ to be an arbitrary, in some sense, a.d.p.t.

We call a pseudo-trajectory $w(\bullet \mid \bullet)$ piecewise continuous if the property

$$
\begin{equation*}
w(\bullet \mid t+h)-w(\bullet \mid t) \rightarrow 0 \tag{4.3.1.3}
\end{equation*}
$$

as $h \rightarrow 0$ holds for all $t$ except perhaps a countable set without points of condensation.

Let $U \subset U[\rho, \sigma]$ for some $\sigma>0$. A pseudo-trajectory $w(\bullet \mid \bullet)$ is called $\omega$-dense in $U$ if $\Omega(w)=U$ (see (4.1.3.2)), i.e.,

$$
\begin{equation*}
\left\{v \in U[\rho]:\left(\exists t_{j} \rightarrow \infty\right) v=\mathcal{D}^{\prime}-\lim w\left(\bullet \mid e^{t_{j}}\right)\right\}=U \tag{4.3.1.4}
\end{equation*}
$$

We have proved already that $v(\bullet \bullet \bullet)$ defined by (4.2.5.2) has this property (see (4.2.6.1)).

Now we consider again the dynamical system $\left(T_{\bullet}, U\right)$ where $U \subset U[\rho, \sigma]$ for some $\sigma>0$ and $T_{t}$ is defined by (4.2.1.1).

Theorem 4.3.1.2 (A.D.P.T. and Chain Recurrence) $\left(T_{\bullet}, U\right)$ is chain recurrent iff there exists an a.d.p.t. that is piecewise continuous and $\omega$-dense in $U$.

Necessity has been proved already, because the pseudo-trajectory (4.2.6.2) possesses this property. Sufficiency will be proved later.

The claim of piecewise continuity can be justified by
Theorem 4.3.1.3 For any $u \in S H(\rho(r))$ there exists a piecewise continuous pseudotrajectory $w(\bullet \mid \bullet)$ such that

$$
\begin{equation*}
u_{t}-w(\bullet \mid t) \rightarrow 0 \tag{4.3.1.5}
\end{equation*}
$$

as $t \rightarrow \infty$.
Of course, $w(\bullet \mid \bullet)$ is a.d.p.t.
Exercise 4.3.1.2 Check this.

### 4.3.2

Proof of Theorem 4.3.1.3. Let $\left\{t_{n}\right\}$ be any sequence such that

$$
\begin{equation*}
t_{n} \rightarrow \infty, t_{n+1} / t_{n} \rightarrow 1 \tag{4.3.2.1}
\end{equation*}
$$

for example, $t_{n}=n$.
There exists a sequence $\left\{v_{n}\right\} \subset \mathbf{F r}[u]$ such that

$$
\begin{equation*}
u_{t_{n}}-v_{n} \rightarrow 0 . \tag{4.3.2.2}
\end{equation*}
$$

Set

$$
\begin{equation*}
w(\bullet \mid t):=v_{n}, \text { for } t_{n}<t \leq t_{n+1} . \tag{4.3.2.3}
\end{equation*}
$$

This is a piecewise continuous function.
Let us prove that

$$
\begin{equation*}
u_{t}-w(\bullet \mid t) \rightarrow 0 . \tag{4.3.2.4}
\end{equation*}
$$

Assume the opposite; i.e., there exists a sequence $\left\{t_{k}^{\prime}\right\}$ such that it is not true. We can suppose that

$$
\begin{equation*}
u_{t_{k}^{\prime}} \rightarrow w_{1} \in \mathbf{F r}[u], w\left(\bullet \mid t_{k}^{\prime}\right) \rightarrow w_{2} \in \mathbf{F r}[u], w_{1} \neq w_{2} \tag{4.3.2.5}
\end{equation*}
$$

Let us find a sequence $\left\{n_{k}\right\}$ such that $t_{n_{k}}<t_{k}^{\prime}<t_{n_{k}+1}$. Then

$$
\begin{equation*}
t_{n_{k}} / t_{k}^{\prime} \rightarrow 1 \tag{4.3.2.6}
\end{equation*}
$$

From (4.3.2.5), (4.3.2.3) and (4.3.2.2) we have

$$
\begin{equation*}
u_{t_{n_{k}}} \rightarrow w_{2} \tag{4.3.2.7}
\end{equation*}
$$

Then we have, using properties of $(\bullet)_{t}$ and the proximate order,

$$
\begin{equation*}
u_{t_{k}^{\prime}}=\left(u_{t_{n_{k}}}\right)_{\left[t_{k}^{\prime} / t_{n_{k}}\right]}\left(1+o\left(\log \left(t_{k}^{\prime} / t_{n_{k}}\right)\right) \rightarrow w_{2}\right. \tag{4.3.2.8}
\end{equation*}
$$

because of (4.3.2.6) and the continuity of $u_{[t]}$ on $(u, t)$.
However (4.3.2.8) contradicts (4.3.2.5). Thus (4.3.2.4) holds.
4.3.3 Now we will prepare the proof of Theorem 4.3.1.2.

Let $\left\{v_{k}, k=1,2, \ldots\right\} \subset U[\rho, \sigma]$ for some $\sigma$ be a sequence of functions and $\left\{r_{k}, k=1,2, \ldots\right\},\left\{t_{k} k=1,2, \ldots\right\}$ be two sequences such that

$$
\begin{equation*}
0<r_{k}<t_{k}<r_{k+1}, k=1,2, \ldots \tag{4.3.3.1}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{k \rightarrow \infty} t_{k} / r_{k}=\lim _{k \rightarrow \infty} r_{k+1} / t_{k}=\infty \tag{4.3.3.2}
\end{equation*}
$$

Set

$$
\begin{equation*}
w^{*}(\bullet \mid t):=\left(v_{k}\right)_{t / t_{k}}, \text { for } t \in\left[r_{k}, r_{k+1}\right) \tag{4.3.3.3}
\end{equation*}
$$

where $k=1,2, \ldots$.
Theorem 4.3.3.1 Let $w(\bullet \mid \bullet) \subset U$ be an arbitrary $\omega$-dense a.d.p.t. and $\left\{p_{j}, j=\right.$ $1,2, \ldots\} \subset U$ an arbitrary sequence. Then there exists a sequence $\left\{v_{k}, k=\right.$ $1,2, \ldots\} \supset\left\{p_{j}, j=1,2, \ldots\right\}$ and sequences $\left\{r_{k}, k=1,2, \ldots\right\}$ and $\left\{t_{k}, k=\right.$ $1,2, \ldots\}$ satisfying (4.3.3.1) and (4.3.3.2) such that for $w^{*}(\bullet \mid \bullet)$ determined by (4.3.3.3) the condition

$$
\begin{equation*}
w^{*}(\bullet \mid t)-w(\bullet \mid t) \rightarrow 0 \tag{4.3.3.4}
\end{equation*}
$$

as $t \rightarrow \infty$ is fulfilled.
This proposition shows that any $\omega$-dense a.d.p.t. is equivalent to one constructed of long pieces of trajectories of the dynamical system $T_{\bullet}$.

Proof of Theorem 4.3.3.1. We can take sequences $\left\{\epsilon_{j} \downarrow 0, j=1,2, \ldots\right\}$ and $\left\{b_{j} \uparrow \infty, j=1,2, \ldots\right\}$ and choose a sequence $\left\{\tau_{j}, j=1,2, \ldots\right\}$ such that the inequalities

$$
\begin{equation*}
d\left(T_{\tau} p_{j}-T_{\tau} w\left(\bullet \mid \tau_{j}\right)\right)<\epsilon_{j} / 2 \tag{4.3.3.5}
\end{equation*}
$$

and

$$
\begin{equation*}
d\left(T_{\tau} w(\bullet \mid t)-w\left(\bullet \mid e^{\tau} t\right)\right)<\epsilon_{j} / 2, t>\tau_{j} \tag{4.3.3.6}
\end{equation*}
$$

are fulfilled uniformly with respect to $\tau \in\left[b_{j+1}^{-2}, b_{j}^{2}\right]$.

Indeed, $w(\bullet \mid \bullet)$ is $\omega$-dense in $U$, and hence we can find $\tau_{n} \rightarrow \infty$ such that

$$
p_{n}-w\left(\bullet \mid \tau_{n}\right) \rightarrow 0
$$

Set in Lemma 4.1.2.2,

$$
p_{n}:=p_{n}, q_{n}:=w\left(\bullet \mid \tau_{n}\right), \gamma_{n}:=2 \log b_{j} .
$$

Then for any $\epsilon_{j}$ we can find $\tau_{j}:=\tau_{n_{j}}$ such that (4.3.3.5) holds uniformly with respect to $\tau \in\left[b_{j+1}^{-2}, b_{j}^{2}\right]$.

The inequality (4.3.3.6) holds, because $w(\bullet \mid \bullet)$ is asymptotically dynamical (see (4.3.1.1)).

We can also suppose without loss of generality that $\tau_{j}>\tau_{j-1} b_{j-1}^{2}$, i.e., that the sequence $\left\{\tau_{j}\right\}$ is rather thin.

The inequality (4.3.3.5) shows that for intervals of $t$ that are determined by the inequality $b_{j+1}^{-2} \leq t / \tau_{j} \leq b_{j}^{2}$ our pseudo-trajectory is already close to some trajectories.

Now we divide the spaces between such intervals into equal parts in the $\operatorname{logarithmic}$ scale such that their logarithmic lengths would be between $\log b_{j}$ and $\log b_{j+1}$, so that they tend to infinity.

To this end, set

$$
n_{j}:=\left[\frac{\log \tau_{j+1}-\log \tau_{j}}{b_{j}}\right]
$$

where [•] means the entire part, and

$$
\gamma_{j}:=\left(\tau_{j+1} / \tau_{j}\right)^{\frac{1}{2 n_{j}}}
$$

It is clear that $b_{j}<\gamma_{j}<b_{j}^{2}$. As centers of new intervals we take the points

$$
\tau_{j, l}:=\tau_{j} \gamma_{j}^{2 l}, l=0,1, \ldots, n_{j} .
$$

Thus $\tau_{j, 0}=\tau_{j}$ and $\tau_{j, n_{j}}=\tau_{j+1}$. The ends of the intervals are $\tau_{j, l} / \gamma_{j}$ and $\tau_{j, l} \gamma_{j}$. Now we complete the sequence $\left\{p_{j}\right\}$ by the values of the pseudo-trajectory $w(\bullet \mid t)$ in the centers of the intervals, i.e., we set

$$
p_{j, l}:=w\left(\bullet \mid \tau_{j, l}\right), l=1, \ldots, n_{j}-1
$$

For $t \in\left(\tau_{j, l} / \gamma_{j}, \tau_{j, l} \gamma_{j}\right), l=1, \ldots, n_{j}-1$ we have

$$
\begin{equation*}
d\left(\left(p_{j, l}\right)_{t / \tau_{j, l}}-w(\bullet \mid t)\right)<\epsilon_{j} / 2 \tag{4.3.3.7}
\end{equation*}
$$

because of (4.3.3.6).
For $l=0$ and $l=n_{j}$ we set accordingly

$$
p_{j, 0}:=p_{j} ; p_{j, n_{j}}:=p_{j+1} .
$$

Using (4.3.3.5) and (4.3.3.6) we have an inequality like (4.3.3.7) for $l=0, l=n_{j}$ but with $\epsilon_{j}$ instead of $\epsilon_{j} / 2$.

We complete the proof, re-denoting all the centers $\tau_{j, l}$ as $t_{k}$, all the ends as $r_{k}$ and all the $p_{j, l}$ as $v_{k}$.

### 4.3.4

Proof of sufficiency in Theorem 4.3.1.2. A direct corollary of the previous Theorem 4.3.3.1 is

$$
\begin{equation*}
w^{*}\left(\bullet \mid r_{k}-0\right)-w^{*}\left(\bullet \mid r_{k}\right) \rightarrow 0 \tag{4.3.4.1}
\end{equation*}
$$

as $k \rightarrow \infty$.
Really, $w^{*}(\bullet \mid \bullet)$ is an a.d.p.t.
Exercise 4.3.4.1 Check this as in Theorem 4.3.1.1 using that $w(\bullet \mid \bullet)$ is asymptotically dynamical.

For $\tau \in[-\epsilon, 0]$ and $t=r_{k}$ we have uniformly on $\tau$,

$$
T_{\tau} w^{*}(\bullet \mid t)-w^{*}(\bullet \mid t)=T_{\tau}\left(v_{k}\right)_{r_{k} / t_{k}}-\left(v_{k+1}\right)_{r_{k} / t_{k}} \rightarrow 0
$$

Setting $\tau=0$ we obtain (4.3.4.1).
Let $V \subset U$ be an arbitrary open set, $\epsilon>0$ arbitrarily small and $s>0$ arbitrarily large. We show that there exists an $(\epsilon, s)$-chain from $V$ to $V$.

Choose $s_{1}$ such that
i. for $r_{k}>s_{1}, d\left(w^{*}\left(\bullet \mid r_{k}-0\right), w^{*}\left(\bullet \mid r_{k}\right)\right)<\epsilon$. This is possible by virtue of (4.3.4.1).
ii. $w\left(\bullet \mid s_{1}\right) \in V$. This is possible because $w(\bullet \mid \bullet)$ is $\omega$-dense.
iii. $d\left(w^{*}(\bullet \mid t), w(\bullet \mid t)\right)<d\left(w\left(\bullet \mid s_{1}\right), \partial V\right)$ for $t>s_{1}$. This is possible because of Theorem 4.3.3.1.

Choose $s_{2}>s_{1}$ such that $w\left(\bullet \mid s_{2}\right) \in V$. This is possible because $w(\bullet \mid \bullet)$ is $\omega$-dense. Then the pseudo-trajectory $w^{*}\left(\bullet \mid e^{t}\right)$ for $s_{1} \leq e^{t} \leq s_{2}$ is an $(\epsilon, s)$-chain connecting $w^{*}\left(\bullet \mid s_{1}\right)$ and $w^{*}\left(\bullet \mid s_{2}\right)$ that belong to $V$.

Exercise 4.3.4.2 Check this in detail.
Hence $\left(T_{\bullet}, U\right)$ is chain recurrent.
4.3.5 We will prove one more existence theorem that is a corollary of Theorem 4.2.1.2.

Theorem 4.3.5.1 Let $\Lambda \subset U[\rho]$ be a compact connected and $T_{\bullet}$-invariant subset of $U[\rho]$. Then for any proximate order $\rho(r) \rightarrow \rho$ there exists $u \in S H\left(\mathbb{R}^{m}, \rho, \rho(r)\right)$ such that

$$
\begin{align*}
h(x, u) & =\sup \{v(x): v \in \Lambda\},  \tag{4.3.5.1}\\
\underline{h}(x, u) & =\inf \{v(x): v \in \Lambda\} . \tag{4.3.5.2}
\end{align*}
$$

Proof. Let $U:=\operatorname{Conv} \Lambda$ be the convex hull of $\Lambda$. It is linearly connected and hence polygonally connected (see 4.1.4). By Theorem 4.1.4.1 it is chain recurrent and by Theorem 4.2.1.2 for any proximate order $\rho(r) \rightarrow \rho$ there exists $u \in S H\left(\mathbb{R}^{m}, \rho, \rho(r)\right)$
such that

$$
\operatorname{Fr}\left[u, \rho(r), V_{t}, \mathbb{R}^{m}\right]=U
$$

Since every $v \in U$ can be represented in the form $v=a v_{1}+(1-a) v_{2}$ for $0 \leq a \leq$ $1, v_{1}, v_{2} \in \Lambda$ we obtain (4.3.5.1) and (4.3.5.2) from Theorem 3.2.1.1 (Properties of Indicators), h2).

Exercise 4.3.5.1 Check this.
4.3.6 In applications we need the following

Theorem 4.3.6.1 Let $p \in P \subset \mathbb{R}^{m}$ and let $P$ be a connected closed set. Let $U_{P}:=$ $\left\{v(z, p): p \in P \subset \mathbb{R}^{m}\right\}$ be a family of functions with parameter $p$ such that for every $p \in P, v(\bullet, p) \in U[\rho]$ and satisfy the condition (4.1.3.3). Then there exists $u \in S H(\rho(r))$ such that $\mathbf{F r}[u]=U_{P}$.

This is a direct corollary of Theorems 4.1.1.2, 4.1.3.3 and 4.2.1.2.
Exercise 4.3.6.1 Explain this in detail.
4.3.7 In the next three sections we return to the periodic limit sets (see Theorem 4.1.7.1). We show that the limit set $\operatorname{Fr}\left[u, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right]$ of every subharmonic function $u \in S H\left(\rho(r), \mathbb{R}^{m}\right), \rho(r) \rightarrow \rho$ for non-integer $\rho$ can be approximated in some sense by periodic limit sets ([Gi], [GLO, Ch. 3, § 2, Thm. 10]).

Here we give some definitions. Let $X_{n} \subset U[\rho], n=1,2, \ldots$ be a sequence of compact sets. We say that $X_{n}$ converges to a compact set $Y \subset U[\rho]$, i.e.,

$$
\begin{equation*}
\mathcal{D}^{\prime}-\lim _{n \rightarrow \infty} X_{n}=Y \tag{4.3.7.1}
\end{equation*}
$$

if the following two conditions hold:
converg1) $\quad \forall x_{n} \in X_{n}, n=1,2, \ldots \exists x_{n_{j}} \in X_{n_{j}}, j=1,2, \ldots$ and $y \in Y$ such that $\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} x_{n_{j}}=y$;
converg2) $\quad \forall y \in Y \exists x_{n} \in X_{n}, n=1,2, \ldots$, such that $x_{n} \rightarrow y$.
On every compact set $K$ in $\mathcal{D}^{\prime}$-topology one can introduce a metric $d(\bullet, \bullet)$ such that the topology generated by this metric is equivalent to $\mathcal{D}^{\prime}$-topology (see, e.g., $[A G(1982)])$.

Denote by

$$
X_{\epsilon}:=\{y \in K: \exists x \in X \text { such that } d(y, x)<\epsilon\}
$$

the $\epsilon$-neighborhood of $X$.
Let $X, Y$ be two compact sets. Set

$$
d(X, Y):=\inf \left\{\epsilon: X \subset Y_{\epsilon}, Y \subset X_{\epsilon}\right\} .
$$

Exercise 4.3.7.1 Prove the assertion

$$
\begin{equation*}
(4.3 .7 .1) \Longleftrightarrow\left\{d\left(X_{n}, Y\right) \rightarrow 0\right\} \tag{4.3.7.2}
\end{equation*}
$$

We prove the following
Theorem 4.3.7.1 (Approximation by Periodic Limit Sets) Let $u \in S H\left(\rho(r), \mathbb{R}^{m}\right)$, $\rho(r) \rightarrow \rho$ for non-integer $\rho$. Then for every $V_{\bullet}$ there exists a sequence $u_{n} \in$ $S H\left(\rho(r), \mathbb{R}^{m}\right)$ with periodic limit sets $\operatorname{Fr}\left[u_{n}, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right]$ such that $\operatorname{Fr}\left[u_{n}, \bullet\right] \rightarrow$ $\operatorname{Fr}[u, \bullet]$.

This theorem is a corollary of the following
Theorem 4.3.7.2 Let $\mu \in \mathcal{M}\left(\rho(r), \mathbb{R}^{m}\right), \rho(r) \rightarrow \rho$ for non-integer $\rho$. Then for every $V_{\bullet}$ there exists a sequence $\mu_{n} \in S H\left(\rho(r), \mathbb{R}^{m}\right)$ with periodic limit sets $\operatorname{Fr}\left[\mu_{n}, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right]$ such that $\operatorname{Fr}\left[\mu_{n}, \bullet\right] \rightarrow \operatorname{Fr}[\mu, \bullet]$.

Proof of Theorem 4.3.7.1. The canonical potential $u(x):=\Pi(x, \mu, p)$ (see (2.9.2.1)) of a measure $\mu \in \mathcal{M}\left(\rho(r), \mathbb{R}^{m}\right)$ belongs to $S H\left(\rho(r), \mathbb{R}^{m}\right)$ by Theorem 2.9.3.3 and has a limit set

$$
\mathbf{F r}[u, \bullet]=\left\{\Pi(\bullet, \nu, p): \nu \in \mathbf{F r}\left[\mu_{u}\right]\right\}
$$

by Theorem 3.1.4.4 (*Hadamard). The potentials $u_{n}(x):=\Pi\left(x, \mu_{n}, p\right)$ have periodic limit sets

$$
\mathbf{F r}\left[u_{n}, \bullet\right]=\left\{\Pi(\bullet, \nu, p): \nu \in \mathbf{F r}\left[\mu_{u_{n}}\right]\right\}
$$

by Theorem 3.1.5.0. Let us prove that

$$
\operatorname{Fr}\left[u_{n}, \bullet\right]=: X_{n} \rightarrow Y:=\operatorname{Fr}[u, \bullet] .
$$

If $v_{n} \in \operatorname{Fr}\left[u_{n}, \bullet\right]$ then from the corresponding sequence of $\nu_{n}:=\nu_{v_{n}} \in \operatorname{Fr}\left[\mu_{n}, \bullet\right]$ we can find a subsequence $\nu_{n_{j}}$ and $\nu \in \operatorname{Fr}[\mu, \bullet]$ such that $\nu_{n} \rightarrow \nu$ (by Theorem 2.2.3.2 (Helly)). It is easy to check, using Theorem 3.1.4.3 (*Liouville), that $v^{*}=\mathcal{D}^{\prime}-$ $\lim _{j \rightarrow \infty} \Pi\left(\bullet, \nu_{n_{j}}\right)$ exists and coincides with $v=\Pi(\bullet, \nu, p) \in \mathbf{F r}[u, \bullet]$.

So the condition converg1) is verified. In the same way one can check converg2).

Exercise 4.3.7.2 Prove this in detail.
4.3.8 Now we are going to prove Theorem 4.3.7.2. We begin from

Proposition 4.3.8.1 For any $\mu \in \mathcal{M}(\rho(r), \bullet)$ there exists $\hat{\mu} \in \mathcal{M}(\rho, \bullet)$ such that

$$
\begin{equation*}
\operatorname{Fr}[\hat{\mu}, \rho, \bullet]=\mathbf{F r}[\mu, \rho(r), \bullet] . \tag{4.3.8.0}
\end{equation*}
$$

In other words we can suppose further that $\rho(r) \equiv \rho$.

Proof. Set $L(r)=r^{\rho(r)-\rho}$ and

$$
\begin{equation*}
\hat{\mu}(d x):=L^{-1}(|x|) \mu(d x) . \tag{4.3.8.1}
\end{equation*}
$$

Using properties of proximate order (Section 2.8.1., po1)-po4)), it is easy to check that

$$
\begin{equation*}
\left[L^{-1}(r)\right]^{\prime}=L^{-1}(r) o(1) \text { and }[L(r)]^{\prime}=L(r) o(1), \text { as } \rightarrow 0 \tag{4.3.8.2}
\end{equation*}
$$

Exercise 4.3.8.1 Prove this.
Let us show that $\hat{\mu} \in \mathcal{M}[\rho, \Delta]$ for some $\Delta$. Indeed

$$
\frac{\hat{\mu}(R)}{R^{\rho+m-2}}=R^{-\rho-m+2} \int_{0}^{R} \frac{\mu(d r)}{L(r)}=\left.\frac{\mu(r)}{R^{\rho+m-2} L(r)}\right|_{0} ^{R}+R^{-\rho-m+2} \int_{0}^{R} \mu(r)\left(L^{-1}\right)^{\prime} d r
$$

We suppose that $\mu(r)=0$ in some neighborhood of zero. Using (4.3.8.2) we obtain further for the last expression,

$$
\mu(R) R^{-\rho(r)}+R^{-\rho-m+2} \int_{0}^{R} \mu(r)\left(L^{-1}\right) o(1) d r
$$

Using the l'Hospital rule, we obtain
$\lim _{R \rightarrow \infty} R^{-\rho-m+2} \int_{0}^{R} \mu(r)\left(L^{-1}(r)\right) o(1 / r) d r=(-\rho-m+2) \lim _{R \rightarrow \infty} \mu(R) R^{-\rho(R)} o(1 / R)$.
Thus

$$
\limsup _{R \rightarrow \infty} \frac{\hat{\mu}(R)}{R^{\rho+m-2}} \leq \limsup _{R \rightarrow \infty} \mu(R)\left(L^{-1}(R)\right)[1+o(1 / R)]=\bar{\Delta}[\mu, \rho(r)]<\infty
$$

Let us note that $\mu_{t}=L(t) \hat{\mu}_{[t]}$. This implies equality (4.3.8.0) because $L(t) \rightarrow 1$ as $t \rightarrow \infty$.

Exercise 4.3.8.2. Prove this in detail.

Proof of Theorem 4.3.7.2. As we already said we can also suppose that $\mu \in \mathcal{M}[\rho]$. Let $\nu \in \mathbf{F r}[\mu]$. We can suppose that

$$
\begin{equation*}
\nu(\{|x|=1\})=0 . \tag{4.3.8.2a}
\end{equation*}
$$

Otherwise we can find $\tau$ such that $\nu_{[\tau]}(\{|x|=1\})=0$ and if $\nu_{n} \rightarrow \nu_{\tau}$ and are periodic, then $\left(\nu_{n}\right)_{[1 / \tau]}$ are also periodic and $\left(\nu_{n}\right)_{[1 / \tau]} \rightarrow \nu$.

Let $r_{n} \rightarrow \infty$ be such that $\mu_{\left[r_{n}\right]} \rightarrow \nu$. By passing to subsequences we can make $r_{n+1} / r_{n}>r_{n}$.

Denote $K_{n}:=\left\{x: r_{n} \leq|x|<r_{n+1}\right\}$. Set for every $E \subset K_{n}$

$$
\left.\mu_{n}\right|_{E}:=\mu(E)
$$

and continue it periodically with the period $P_{n}=r_{n+1} / r_{n}$ by the equality

$$
\begin{equation*}
\mu_{n}\left(P_{n}^{k} E\right)=P_{n}^{k \rho} \mu(E), k= \pm 1, \pm 2, \ldots \tag{4.3.8.3}
\end{equation*}
$$

Since every $X \in \mathbb{R}^{m}$ can be represented in the form

$$
X=\bigcup_{k=-\infty}^{\infty}\left\{X \cap P_{n}^{k} K_{n}\right\}
$$

we can define

$$
\mu_{n}(X):=\sum_{k=-\infty}^{\infty} \mu_{n}\left(\left\{X \cap P_{k}^{n} K_{n}\right\}\right)
$$

It is easy to check that $\mu_{n}$ is periodic with period $P_{n}$ and $\mu_{n} \in \mathcal{M}[\rho, \Delta]$ with $\Delta$ independent of $n$.

Exercise 4.3.8.3 Check this.
Let us prove that

$$
\begin{equation*}
\operatorname{Fr}[\mu]=\lim _{n \rightarrow \infty} \operatorname{Fr}\left[\mu_{n}\right] . \tag{4.3.8.4}
\end{equation*}
$$

Check the condition converg1). Let $\nu_{n_{j}} \in \operatorname{Fr}\left[\mu_{n_{j}}\right]$ and suppose $\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} \nu_{n_{j}}:=\nu$. Let us prove that $\nu \in \operatorname{Fr}[\mu]$.

Since $\operatorname{Fr}\left[\mu_{n_{j}}\right]$ is a periodic limit set,

$$
\nu_{n_{j}}=\left(\mu_{n_{j}}\right)_{\left[\tau_{j}\right]}
$$

Take $k_{j}$ such that

$$
\tau_{j}^{\prime}:=\tau_{j} P_{n_{j}}^{k_{j}} \in\left[r_{n_{j}}, r_{n_{j}+1}\right)
$$

From periodicity $\mu_{n_{j}}$ we obtain

$$
\nu_{n_{j}}=\left(\mu_{n_{j}}\right)_{\left[\tau_{j}^{\prime}\right]} .
$$

Passing to a subsequence if necessary, we can consider three cases:
i) $\lim _{j \rightarrow \infty} \tau_{j}^{\prime} / r_{n_{j}}=\infty, \lim _{j \rightarrow \infty} \tau_{j}^{\prime} / r_{n_{j}+1}=0$;
ii) $\lim _{j \rightarrow \infty} \tau_{j}^{\prime} / r_{n_{j}}=\tau ; 1 \leq \tau<\infty$; In this case we have also $\lim _{j \rightarrow \infty} \tau_{j}^{\prime} / r_{n_{j}+1}=0$.
iii) $\lim _{j \rightarrow \infty} \tau_{j}^{\prime} / r_{n_{j}+1}=\tau ; 0<\tau \leq 1$; In this case we have also $\lim _{j \rightarrow \infty} \tau_{j}^{\prime} / r_{n_{j}}=\infty$.

Consider the case i). Let $\phi \in \mathcal{D}\left(\mathbb{R}^{m} \backslash O\right)$. Then $\operatorname{supp} \phi\left(x / \tau_{j}^{\prime}\right) \subset\left(r_{n_{j}}, r_{n_{j}+1}\right)$ for $j \geq j_{0}$. It is easy to see that, for $j \geq j_{0}$,

$$
\left\langle\left(\mu_{n_{j}}\right)_{\left[\tau_{j}^{\prime}\right]}, \phi\right\rangle=\left\langle\mu_{\left[\tau_{j}^{\prime}\right]}, \phi\right\rangle .
$$

Exercise 4.3.8.4 Check this.
Since $\mu_{\left[\tau_{j}^{\prime}\right]} \rightarrow \nu \in \operatorname{Fr}[\mu]$ by definition the condition converg1) holds for the case i).

Consider the case ii). Recall that $O \notin \operatorname{supp} \phi$. Then there exists $1 \leq c<\infty$ such that

$$
\operatorname{supp} \phi \subset\{x:|x| \in(1 / c, c)\}
$$

Define

$$
\phi_{t}(x):=\phi(x / t)(1 / t)^{\rho} .
$$

Represent $\tau_{j}^{\prime}$ in the form

$$
\tau_{j}^{\prime}:=e_{j} \tau r_{n_{j}} \quad \text { where } \quad e_{j}:=\frac{\tau_{j}^{\prime}}{r_{n_{j}} \tau}
$$

The condition ii) means that

$$
\begin{equation*}
e_{j} \rightarrow 1 \tag{4.3.8.4a}
\end{equation*}
$$

Compute

$$
\left\langle\nu_{j}, \phi\right\rangle:=\left\langle\left(\mu_{n_{j}}\right)_{\left[\tau_{j}^{\prime}\right]}, \phi\right\rangle=\left\langle\mu_{n_{j}},\left(\left(\phi_{\tau}\right)_{e_{j}}\right)_{r_{n_{j}}}\right\rangle
$$

Note that

$$
\operatorname{supp} \phi_{\tau} \subset\{x:|x| \in(\tau / c, \tau c)\}
$$

We can increase $c$ so that $1 \in(\tau / c, \tau c)$.
Consider the following partition of unity. Choose the functions $\eta_{k} \in \mathcal{D}\left(\mathbb{R}^{m}\right)$, $k=1,2,3$ so that

$$
\eta_{1}(t)+\eta_{2}(t)+\eta_{3}(t)=1
$$

for $t \geq 1$ and

$$
\begin{aligned}
& \operatorname{supp} \eta_{1} \subset\{x:|x|<1-\epsilon\} \\
& \operatorname{supp} \eta_{2} \subset\{x:|x| \in(1-2 \epsilon, 1+2 \epsilon)\} \\
& \operatorname{supp} \eta_{3} \subset\{x:|x|>1+\epsilon\}
\end{aligned}
$$

where $\epsilon$ is an arbitrary number, satisfying

$$
\tau / c<1-2 \epsilon<1+2 \epsilon<\tau c
$$

Represent $\phi_{\tau}$ in the form

$$
\phi_{\tau}=\psi_{1}+\psi_{2}+\psi_{3}, \quad \text { where } \quad \psi_{k}=\phi_{\tau} \eta_{k}, k=1,2,3
$$

In this notation

$$
\begin{equation*}
\left\langle\nu_{j}, \phi\right\rangle=\sum_{k=1}^{3}\left\langle\mu_{n_{j}},\left(\psi_{k}\right)_{e_{j} r_{n_{j}}}\right\rangle \tag{4.3.8.5}
\end{equation*}
$$

Choose $j_{\epsilon}$ such that for $j \geq j_{\epsilon}$ the following inclusions hold:

$$
\begin{aligned}
& \operatorname{supp}\left(\psi_{1}\right)_{e_{j} r_{n_{j}}} \subset\left\{x:|x| \in\left(r_{n_{j}} \tau / c, r_{n_{j}}(1-\epsilon)\right)\right\} \\
& \operatorname{supp}\left(\psi_{2}\right)_{e_{j} r_{n_{j}}} \subset\left\{x:|x| \in\left((1-\epsilon) r_{n_{j}}, r_{n_{j}}(1+\epsilon)\right)\right\} \\
& \operatorname{supp}\left(\psi_{3}\right)_{e_{j} r_{n_{j}}} \subset\left\{x:|x| \in\left((1+\epsilon) r_{n_{j}}, \tau r_{n_{j}}\right\}\right.
\end{aligned}
$$

Thus for $\psi_{3}$ we have

$$
\begin{aligned}
\left\langle\mu_{n_{j}},\left(\psi_{3}\right)_{e_{j} r_{n_{j}}}\right\rangle & =\int\left(e_{j} r_{n_{j}}\right)^{-\rho} \psi_{3}\left(|x| /\left(e_{j} r_{n_{j}}\right)\right) \mu_{n_{j}}(d x) \\
& \left.=\int\left(e_{j} r_{n_{j}}\right)^{-\rho} \psi_{3}\left(|x| /\left(e_{j} r_{n_{j}}\right)\right) \mu(d x)=\left\langle\mu_{\left[r_{n_{j}}\right.}\right],\left(\psi_{3}\right)_{e_{j}}\right\rangle
\end{aligned}
$$

Since $\mu_{\left[r_{n_{j}}\right]} \xrightarrow{\mathcal{D}^{\prime}} \nu$ and $\left(\psi_{3}\right)_{e_{j}} \xrightarrow{\mathcal{D}} \psi_{3}$ we have (see Theorem 2.3.4.6)

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left\langle\mu_{n_{j}},\left(\psi_{3}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\nu, \psi_{3}\right\rangle \tag{4.3.8.6}
\end{equation*}
$$

Consider the addend with $\psi_{1}$. Because of periodicity $\mu_{n_{j}}$ we have

$$
\left\langle\mu_{n_{j}},\left(\psi_{1}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\left(\mu_{n_{j}}\right)_{\left[P_{n_{j}}\right]},\left(\psi_{1}\right)_{e_{j} r_{n_{j}}}\right\rangle
$$

Transforming the RHS we obtain

$$
\left\langle\left(\mu_{n_{j}}\right)_{\left[P_{n_{j}}\right]},\left(\psi_{1}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\mu_{n_{j}},\left(\left(\psi_{1}\right)_{P_{n_{j}}}\right)_{e_{j} r_{n_{j}}}\right\rangle .
$$

Since $P_{n_{j}}=r_{n_{j}+1} / r_{n_{j}}$ the following inclusion holds for $j \geq j_{\epsilon}$ :

$$
\begin{aligned}
& \operatorname{supp}\left(\psi_{1}\right)_{P_{n_{j}} e_{j} r_{n_{j}}} \subset\left\{x:|x| \in\left(P_{n_{j}} r_{n_{j}} \frac{\tau}{c}, P_{n_{j}} r_{n_{j}}(1-\epsilon)\right)\right\} \\
& \quad=\left\{x:|x| \in\left(r_{n_{j}+1} \frac{\tau}{c}, r_{n_{j}+1}(1-\epsilon)\right)\right\} \subset\left\{x:|x| \in\left(r_{n_{j}}, r_{n_{j}+1}\right)\right\}
\end{aligned}
$$

Thus

$$
\left\langle\mu_{n_{j}},\left(\left(\psi_{1}\right)_{P_{n_{j}}}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\mu,\left(\psi_{1}\right)_{P_{n_{j}} e_{j} r_{n_{j}}}\right\rangle=\left\langle\mu,\left(\psi_{1}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\mu_{\left[r_{n_{j}}\right]},\left(\psi_{1}\right)_{e_{j}}\right\rangle .
$$

Hence

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left\langle\mu_{n_{j}},\left(\psi_{1}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\nu, \psi_{1}\right\rangle \tag{4.3.8.7}
\end{equation*}
$$

because $e_{j} \rightarrow 1$ and $\mu_{\left[r_{n_{j}}\right]} \xrightarrow{\mathcal{D}^{\prime}} \nu$ (see Theorem 2.3.4.6).
From (4.3.8.5), (4.3.8.6) and (4.3.8.7) we obtain

$$
\lim _{j \rightarrow \infty}\left\langle\nu_{j}, \phi\right\rangle=\left\langle\nu, \psi_{1}+\psi_{3}\right\rangle+\lim _{j \rightarrow \infty}\left\langle\mu_{n_{j}},\left(\left(\psi_{2}\right)_{e_{j} r_{n_{j}}}\right\rangle .\right.
$$

Let us estimate the last limit. We have

$$
\lim _{j \rightarrow \infty}\left\langle\mu_{n_{j}},\left(\psi_{2}\right)_{e_{j} r_{n_{j}}}\right\rangle=\left\langle\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}, \psi_{2}\right\rangle
$$

Define

$$
E_{1}(\epsilon):=\{x:|x| \in(1-2 \epsilon, 1)\} ; E_{2}(\epsilon):=\{x:|x| \in[1,1+2 \epsilon)\} .
$$

Suppose $\epsilon$ is chosen so that

$$
\begin{equation*}
\nu\left(\partial E_{k}\right)=0, k=1,2 \tag{4.3.8.7a}
\end{equation*}
$$

Recall that $\nu$ satisfies the condition (4.3.8.2a), hence $E_{1}, E_{2}$ are $\nu$-squarable and hence (see Theorem 2.2.3.7)

$$
\lim _{n \rightarrow \infty} \mu_{\left[r_{n}\right]}\left(E_{k}(\epsilon)\right)=\nu\left(E_{k}(\epsilon)\right), k=1,2
$$

Define

$$
C_{\phi}:=\max \left\{\phi(x): x \in \mathbb{R}^{m}\right\}
$$

Then for $j \geq j_{\epsilon}$,

$$
\begin{aligned}
\left|\left\langle\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}, \psi_{2}\right\rangle\right| & \leq C_{\phi}\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}\left(E_{1}(\epsilon) \cup E_{2}(\epsilon)\right) \\
& =C_{\phi}\left(\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}\left(E_{1}(\epsilon)\right)+\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}\left(E_{2}(\epsilon)\right)\right.
\end{aligned}
$$

By definition

$$
\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}\left(E_{2}(\epsilon)\right)=\mu_{\left[e_{j} r_{n_{j}}\right]}\left(E_{2}(\epsilon)\right) .
$$

Because of (4.3.8.4a) we obtain

$$
\lim _{j \rightarrow \infty} \mu_{\left[e_{j} r_{n_{j}}\right]}\left(E_{2}(\epsilon)\right)=\nu\left(E_{2}\right)
$$

Exercise 4.3.8.5 Check in detail.
To compute the limit of the first addend we use periodicity of $\mu_{n_{j}}$ :

$$
\mu_{n_{j}}\left(E_{1}(\epsilon)\right)=P_{n_{j}}^{-\rho} \mu_{n_{j}}\left(P_{n_{j}} E_{1}(\epsilon)\right)=\left(\mu_{n_{j}}\right)_{\left[P_{n_{j}}\right]}\left(E_{1}(\epsilon)\right),
$$

where

$$
r_{n_{j}} P_{n_{j}} E_{1}(\epsilon)=\left\{x:|x| \in\left(r_{n_{j}+1}(1-2 \epsilon), r_{n_{j}+1}\right)\right\} .
$$

Thus

$$
\begin{aligned}
\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}\left(E_{1}(\epsilon)\right) & =\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}} P_{n_{j}}\right]}\left(E_{1}(\epsilon)\right) \\
& =\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}+1\right]}\left(E_{1}(\epsilon)\right)=\mu_{\left[e_{j} r_{n_{j}}+1\right]}\left(E_{1}(\epsilon)\right) .
\end{aligned}
$$

From this we obtain

$$
\lim _{j \rightarrow \infty}\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}\left(E_{1}(\epsilon)\right)=\nu\left(E_{1}(\epsilon)\right) .
$$

Therefore

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left|\left\langle\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}, \psi_{2}\right\rangle\right| \leq C_{\phi} \nu\left(E_{1}(\epsilon) \cup E_{2}(\epsilon)\right) . \tag{4.3.8.8}
\end{equation*}
$$

Because of (4.3.8.2a) we have

$$
\begin{equation*}
\nu(\{x:|x| \in(1-2 \epsilon, 1+2 \epsilon)\}) \rightarrow 0 \tag{4.3.8.9}
\end{equation*}
$$

as $\epsilon \rightarrow 0$ over the set of $\epsilon$ satisfying (4.3.8.7a). From (4.3.8.8) and (4.3.8.9) we obtain

$$
\begin{equation*}
\lim _{j \rightarrow \infty}\left\langle\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}, \psi_{2}\right\rangle \rightarrow 0 \tag{4.3.8.10}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle\nu, \psi_{2}\right\rangle \rightarrow 0 \tag{4.3.8.11}
\end{equation*}
$$

when $\epsilon \rightarrow 0$. Hence if $\epsilon \rightarrow 0$ satisfying (4.3.8.7a) we have

$$
\begin{aligned}
& \lim _{j \rightarrow \infty}\left\langle\nu_{j}, \phi\right\rangle-\left\langle\nu_{[\tau]}, \phi\right\rangle \\
& \quad=\lim _{\epsilon \rightarrow 0}\left[\left\langle\nu, \psi_{1}+\psi_{3}\right\rangle+\lim _{j \rightarrow \infty}\left\langle\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}, \psi_{2}\right\rangle-\left\langle\nu, \psi_{1}+\psi_{2}+\psi_{3}\right\rangle\right] \\
& \quad=\lim _{\epsilon \rightarrow 0}\left[\lim _{j \rightarrow \infty}\left\langle\left(\mu_{n_{j}}\right)_{\left[e_{j} r_{n_{j}}\right]}, \psi_{2}\right\rangle-\left\langle\nu, \psi_{2}\right\rangle\right]=0 .
\end{aligned}
$$

The last equality holds because every addend tends to zero.
The case iii) can be considered in an analogous way.
Exercise 4.3.8.6 Consider it.
Thus the condition converg1) was checked.
4.3.9 Now we should check the condition converg2). We need

Lemma 4.3.9.1 Let $\mu \in \mathcal{M}[\rho], \nu \in \operatorname{Fr}[\mu]$ and $r_{n} \rightarrow \infty, n=1,2, \ldots$ be a sequence such that

$$
\begin{equation*}
\mathcal{D}^{\prime}-\lim _{n \rightarrow \infty} \mu\left[r_{n}\right]=\nu_{0} \tag{4.3.9.1}
\end{equation*}
$$

Then passing if necessary to a subsequence, we can find $\left\{r_{n}\right\}$ such that for arbitrarily $\nu \in \operatorname{Fr}[\mu]$ a sequence $t_{j} \rightarrow \infty$ exists such that

$$
\begin{equation*}
\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} \mu\left[t_{j}\right]=\nu \tag{4.3.9.2}
\end{equation*}
$$

and for every $n$ we can find $t_{j} \in\left[r_{n}, r_{n+1}\right]$.
Proof. Note that if the assertion of the lemma is satisfied for the sequence $\left\{r_{n}, n=\right.$ $1,2, \ldots\}$ it is satisfied for every subsequence of $\left\{r_{n}, n=1,2, \ldots\right\}$.

Let $M$ be a countable set that is dense in $\operatorname{Fr}[\mu]$. Since reduction $\mathcal{D}^{\prime}$-topology on $\mathcal{M}[\rho]$ is metrizable, it is sufficient to prove that we can choose a subsequence $r_{n}$ for which assertion of the lemma is satisfied for all $\nu \in M$. We can do it using a diagonal process.

Let $r_{n}^{0} \rightarrow \infty$ be an arbitrary sequence such that

$$
\mathcal{D}^{\prime}-\lim _{n \rightarrow \infty} \mu_{\left[r_{n}^{0}\right]}=\nu_{0}
$$

and let a sequence $t_{j}^{1}$ satisfy the condition

$$
\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} \mu_{\left[t_{j}^{1}\right]}=\nu_{1}
$$

Omitting in the sequence $\left\{r_{n}^{0}\right\}$ the ends of the segments $\left[r_{n}^{0}, r_{n+1}^{0}\right]$ that do not contain elements $\left\{t_{j}\right\}$ we obtain a subsequence $\left\{r_{n}^{1}\right\}$. Continuing in such a way we obtain a subsequence $\left\{r_{n}^{m}\right\}$ satisfying (4.3.9.1) and the sequences $\left\{t_{j}^{1}\right\},\left\{t_{j}^{2}\right\}, \ldots,\left\{t_{j}^{m}\right\}$, $j=1,2 \ldots$ satisfying

$$
\begin{equation*}
\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} \mu_{\left[t_{j}^{l}\right]}=\nu_{l}, l=1,2, \ldots m \tag{4.3.9.3}
\end{equation*}
$$

Taking a diagonal sequence $\left\{r_{n}^{n}\right\}, n=1,2, \ldots$ we observe that it is a subsequence of every subsequence $\left\{r_{n}^{m}\right\}$ and hence satisfies the assertion of the lemma.

Proof of converg2). We can suppose that $\left\{r_{n}\right\}$ from the construction of $\mu_{n}$ with periodic limit sets satisfies the assertion of Lemma 4.3.9.1. Let $\nu \in \operatorname{Fr}[\mu]$ and $\mu_{t_{j}} \rightarrow \nu$ under condition $t_{j} \in\left[r_{n}, r_{n+1}\right]$. We should consider as in the proof of converg1) three cases i), ii) and iii). But all these cases were already considered and hence it was proved that

$$
\nu_{n}:=\left(\mu_{n}\right)_{r_{n}} \rightarrow \nu
$$

Exercise 4.3.9.1 Check this.

### 4.4 Subharmonic curves. Curves with prescribed limit sets

4.4.1 In this paragraph we consider subharmonic functions $u \in S H(\rho(r))$ in the plane of finite type with respect to some proximity order $\rho(r) \rightarrow \rho$.

The pair $\boldsymbol{u}:=\left(u_{1}, u_{2}\right), u_{1}, u_{2} \in S H(\rho(r))$ is called a subharmonic curve (which for brevity we will refer to simply as a curve).

The family

$$
(\boldsymbol{u})_{t}:=\left(\left(u_{1}\right)_{t},\left(u_{2}\right)_{t}\right)
$$

is precompact in the topology of convergence in $\mathcal{D}^{\prime}$-topology on every component. The set of all limits

$$
\operatorname{Fr}[\boldsymbol{u}]:=\left\{\boldsymbol{v}=\left(v_{1}, v_{2}\right): \exists t_{j} \rightarrow \infty, \boldsymbol{v}=\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} \boldsymbol{u}_{t_{j}}\right\}
$$

is called the limit set of the curve $\boldsymbol{u}$.
Actually this set describes coordinated asymptotic behavior of pairs of subharmonic functions.

Theorem 4.4.1.1 $\operatorname{Fr}[\boldsymbol{u}]$ is closed, connected, invariant with respect to $(\bullet)_{[t]}$ (see 3.1.2.4a) and is contained in the set

$$
\boldsymbol{U}[\rho, \boldsymbol{\sigma}]:=\left\{\boldsymbol{v}=\left(v_{1}, v_{2}\right): v_{n}(z) \leq \sigma_{n}|z|^{\rho}, v_{n}(0)=0, n=1,2 .\right\}
$$

where $\boldsymbol{\sigma}:=\left(\sigma_{1}, \sigma_{2}\right)$

Exercise 4.4.1.1 Prove this by using Theorem 3.1.2.2.
Let us define $\boldsymbol{\sigma}>0$ as $\sigma_{n}>0, n=1,2$. Set

$$
\boldsymbol{U}[\rho]:=\bigcup_{\boldsymbol{\sigma}>0} \boldsymbol{U}[\rho, \boldsymbol{\sigma}] .
$$

We will write $\boldsymbol{U} \subset \boldsymbol{U}[\rho]$ if $\boldsymbol{U} \subset \boldsymbol{U}[\rho, \boldsymbol{\sigma}]$ for some $\boldsymbol{\sigma}$.
Since $\left(T_{\bullet}, \boldsymbol{U}[\rho]\right)$ is a dynamical system we have two theorems analogous to Theorems 4.2.1.1 and 4.2.1.2.

Exercise 4.4.1.2 Formulate and prove these theorems.
All the other assertions and definitions of Sections 4.2,4.3 can be repeated for subharmonic curves.

Let $\boldsymbol{U} \subset \boldsymbol{U}[\rho]$. Set

$$
U^{\prime}:=\left\{v^{\prime}: \exists v^{\prime \prime}:\left(v^{\prime}, v^{\prime \prime}\right) \in \boldsymbol{U}\right\}
$$

This is a projection of $\boldsymbol{U}$. Set for $v^{\prime} \in U^{\prime}$,

$$
U^{\prime \prime}\left(v^{\prime}\right):=\left\{v^{\prime \prime}:\left(v^{\prime}, v^{\prime \prime}\right) \in \boldsymbol{U}\right\}
$$

This is the fibre over $v^{\prime}$.
Theorem 4.4.1.2 Let $\boldsymbol{U} \Subset \boldsymbol{U}[\rho]$ be closed and invariant and assume that every fiber $U^{\prime \prime}\left(v^{\prime}\right)$ is convex. Let $U^{\prime}=\mathbf{F r}\left[u^{\prime}\right]$ for some $u^{\prime} \in U(\rho(r))$. Then there exists $u^{\prime \prime} \in U(\rho(r))$ such that $\operatorname{Fr}\left(u^{\prime}, u^{\prime \prime}\right)=\boldsymbol{U}$.

We construct a pseudo -trajectory asymptotics in the form (4.2.5.2) replacing $u$ with $\boldsymbol{u}$ and $v$ with $\boldsymbol{v}$. We can directly check that this curve satisfies the assertion of the theorem.

Exercise 4.4.1.3 Check this.

Theorem 4.4.1.3 (Concordance Theorem) Let $u \in U(\rho(r))$ and $v^{0} \in \mathbf{F r}[u]$, and suppose $v \in U[\rho]$ has the property

$$
\lim _{\tau \rightarrow-\infty} T_{\tau} v=\lim _{\tau \rightarrow+\infty} T_{\tau} v=\tilde{v}
$$

Then there exists a function $w \in U(\rho(r))$ such that the limit set of the curve $\boldsymbol{u}=(u, w) \operatorname{Fr}[\boldsymbol{u}]=(\mathbf{F r}[u], \mathbb{C}(v))$ and for every sequence $t_{n} \rightarrow \infty$ such that $\lim _{n \rightarrow \infty} w_{t_{n}}=v$,

$$
\begin{equation*}
\lim _{n \rightarrow \infty} \boldsymbol{u}_{t_{n}}=\left(v^{0}, v\right) \tag{4.4.1.3.}
\end{equation*}
$$

For proving this theorem we should use a.d.p.t. (4.2.5.2). If $v_{j}=v^{0}$ we replace $v_{j}$ by $\boldsymbol{v}_{j}:=\left(v^{0}, v\right)$. If $v_{j} \neq v^{0}$ we replace $v_{j}$ by $\boldsymbol{v}_{j}:=\left(v_{j}, \tilde{v}\right)$.

Exercise 4.4.1.3 Do that and exploit Theorem 4.3.1.2 and Theorem 4.2.1.2.
Corollary 4.4.1.4 Under conditions of Theorem 4.4.1.3, if $\lim _{n \rightarrow \infty} w_{t_{n}}=T_{\tau} v$, then $\lim _{n \rightarrow \infty} u_{t_{n}}=T_{\tau} v^{0}$.

We should apply $T_{\tau}$ to (4.4.1.3) and use its continuity in $\mathcal{D}^{\prime}$-topology.

## Chapter 5

## Applications to Entire Functions

### 5.1 Growth characteristics of entire functions

5.1.1 Let $f(z)$ be an entire function. The function $u(z):=\log |f(z)|$ is subharmonic in $\mathbb{R}^{2}(=\mathbb{C})$. Hence the scale of growth subharmonic functions considered in Section 2.8 is transferred completely to entire functions. We will mark passing to entire function by changing index $u$ for index $f$. For example,

$$
M(r, f):=M(r, \log |f|), T(r, f):=T(r, \log |f|)
$$

If $u(z):=\log |f(z)|$ has order $\rho[u]=\rho$, then $f(z)$ has order $\rho[f]:=\rho$ and so on.
We will write $f \in A(\rho, \rho(r))$ and say " $f$ is an entire function of order $\rho$ and normal type with respect to proximate order $\rho(r) "$ if $\log |f|$ is a subharmonic function of order $\rho$ and normal type with respect to the same proximate order. Shortly, if $\log |f| \in S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$, then $f \in A(\rho, \rho(r))$.

Exercise 5.1.1.1 Give definitions of

$$
T(r, f), M(r, f), \rho_{T}[f], \rho_{M}[f], \sigma_{T}[f, \rho(r)], \sigma_{M}[f, \rho(r)]
$$

and reformulate all the assertions of Section 2.8 in terms of entire and meromorphic functions.
5.1.2 A divisor of zeros of an entire function can be represented as an integer mass distribution $n$ on a discrete set $\left\{z_{j}\right\} \subset \mathbb{C}$. The multiplicity of a zero $z_{j}$ is the mass concentrated at the point $z_{j}$.

The notation for characteristics of the behavior of zeros will mimic that of the behavior of masses, replacing $\mu$ for $n$. For example, $n\left(K_{r}\right), n(r)$ is the number of zeros (with multiplicities) in the disk $K_{r}, \rho[n]$ is the convergence exponent, $\bar{\Delta}[n]$ is the upper density and so on.

Exercise 5.1.1.2 Give definitions of $N(r, n), \rho_{N}[n], \bar{\Delta}_{N}[n], p[n]$.
5.1.3 The limit set $\mathbf{F r}[f]$ of an entire function $f \in A(\rho, \rho(r))$ is defined as the limit set of the subharmonic function $u(z):=\log |f(z)| \in S H\left(\rho, \rho(r), \mathbb{R}^{2}\right.$ ) (see Section 3.1), i.e.,

$$
\begin{equation*}
\operatorname{Fr}[f]:=\mathbf{F r}[\log |f|] . \tag{5.1.3.1}
\end{equation*}
$$

It possesses, of course, all the properties described in Chapters 3, 4 but it is not clear now if there exists an entire function with prescribed limit set, i.e., whether the subharmonic function in Theorem 4.2.1.2 can be chosen to be $\log |f(z)|$ where $f \in A(\rho, \rho(r))$. It turns out that this is possible and we prove this in Section 5.3.

As it was mentioned in 3.1.1 the general form of $V_{\bullet}$ for the case of the plane is

$$
V_{t} z=z e^{i \gamma \log t}
$$

where $\gamma$ is real.
The limit set $\operatorname{Fr}[n]$ of a divisor $n$ is the limit set of the corresponding mass distribution $n$ (see 3.1.3).

Of course generally speaking $n_{t}$ (see (3.1.3.2)) is not an integer mass distribution.

Exercise 5.1.3.1 Give a complete definition of $\mathbf{F r}[f]$ and $\operatorname{Fr}[n]$, and reformulate all the theorems of Sections 3.1.2, 3.1.3 in terms of entire functions and their zeros.

The connection between $\mathbf{F r}[f]$ and $\mathbf{F r}[n]$ is preserved completely (see Section 3.1.5).

Exercise 5.1.3.2 Reformulate the theorems of Section 3.1.5 for entire functions.
5.1.4 Let $f=f_{1} / f_{2}$ be a meromorphic function, where $f_{1}, f_{2}$ have no common zeros. If $f_{2}(0)=1, f_{1}(0) \neq 0$ and $f_{1}, f_{2} \in A(\rho, \rho(r))$, then $u:=\log \left|f_{1}\right|-\log \left|f_{2}\right| \in$ $\delta S H(\rho, \rho(r))$, and we write $f \in \operatorname{Mer}(\rho, \rho(r))$ and say " $f$ is a meromorphic function of order $\rho$ and normal type with respect to the proximate order $\rho(r)$ ". For $f \in$ $\operatorname{Mer}(\rho, \rho(r))$ we use the following characteristics: $T(r, f), \rho_{T}[f], \sigma_{T}[f, \rho(r)]$. The charge of $\log |f|$ consists of integer positive and negative masses.

## $5.2 \mathcal{D}^{\prime}$-topology and topology of exceptional sets

5.2.1 Let $\alpha$-mes be the Carleson measure defined in Section 2.5.4. Set for $C \subset \mathbb{R}^{2}$,

$$
\begin{equation*}
\alpha-\overline{\mathrm{mes}} C:=\limsup _{R \rightarrow \infty}\left[\alpha-\operatorname{mes}\left(C \cap K_{R}\right)\right] R^{-\alpha} . \tag{5.2.1.1}
\end{equation*}
$$

It is called the relative Carleson $\alpha$-measure.

Theorem 5.2.1.1 (Properties of the Relative Carleson Measure) One has
$\mathrm{rCm} 1)$ If $C$ is bounded $\alpha-\overline{\mathrm{mes}} C=0$;
rCm2)

$$
\alpha-\overline{\operatorname{mes}}\left(C_{1} \cap C_{2}\right) \leq \alpha-\overline{\mathrm{mes}} C_{1}+\alpha-\overline{\mathrm{mes}} C_{2}
$$

i.e., the relative Carleson measure is sub-additive;
rCm3)

$$
C_{1} \subset C_{2} \Rightarrow \alpha-\overline{\mathrm{mes}}\left(C_{1}\right) \leq \alpha-\overline{\mathrm{mes}} C_{2}
$$

i.e., the relative Carleson measure is monotonic with respect to sets;
rCm4)

$$
\alpha_{1}>\alpha_{2} \Rightarrow \alpha_{1}-\overline{\mathrm{mes}} C \leq \alpha_{2}-\overline{\mathrm{mes}} C,
$$

i.e., the relative Carleson measure is monotonic with respect to $\alpha$.

Exercise 5.2.1.1 Prove this.
A set $C \subset \mathbb{R}^{2}$ for which $\alpha-\overline{\mathrm{mes}} C=0$ is called a $C_{0}^{\alpha}-$ set. If $\alpha-\overline{\mathrm{mes}} C=0$ for all $\alpha>0, C$ is called a $C_{0}^{0}-$ set.

Let us recall that if $u_{1}, u_{2} \in S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$, then $u=u_{1}-u_{2} \in \delta S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$ (see Section 2.8.2).

Theorem 5.2.1.2 ( $\mathcal{D}^{\prime}$-topology and Exceptional sets) Let $u \in \delta S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$. In order that

$$
\begin{equation*}
u_{t} \rightarrow 0 \tag{5.2.1.2}
\end{equation*}
$$

in $\mathcal{D}^{\prime}$ as $t \rightarrow \infty$ it is sufficient that

$$
\begin{equation*}
u(z)|z|^{-\rho(|z|)} \rightarrow 0 \tag{5.2.1.3}
\end{equation*}
$$

as $z \rightarrow \infty$ outside some $C_{0}^{2}$-set.
If (5.2.1.2) holds, then (5.2.1.3) holds outside some $C_{0}^{0}$-set.
5.2.2 To prove Theorem 5.2.1.2 we need some auxiliary assertions. Recall that $d z$ is an element of area following the notation of the previous chapters.
Proposition 5.2.2.1 Let $u \in S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$, and $C_{0, R}^{2}:=C_{0}^{2} \cap K_{R}$. Then

$$
\begin{equation*}
\int_{C_{0, R}^{2}}|u|(z) d z=o\left(R^{\rho(R)+2}\right) \tag{5.2.2.1}
\end{equation*}
$$

as $R \rightarrow \infty$.
Proof. Suppose (5.2.2.1) does not hold. Then there exists a sequence $R_{j} \rightarrow \infty$ such that

$$
\begin{equation*}
\lim _{R_{j} \rightarrow \infty} R_{j}^{-\rho\left(R_{j}\right)-2} \int_{C_{0, R_{j}}^{2}}|u|(z) d z=A>0 . \tag{5.2.2.2}
\end{equation*}
$$

Consider the following family of $\delta$-subharmonic functions:

$$
\begin{equation*}
u_{j}(\zeta):=R_{j}^{-\rho\left(R_{j}\right)} u\left(\zeta R_{j}\right) \tag{5.2.2.2a}
\end{equation*}
$$

It can be represented as a difference $u_{j}=u_{1, j}-u_{2, j}$ of subharmonic functions of the same form.

Thus it is precompact in $L_{\text {loc }}$ (Theorem 2.7.1.3). Let us choose a convergent subsequence for which we keep the same notation. Its limit $v$ is a locally summable function.

Now let $\chi_{j}$ be the characteristic functions of the sets

$$
E_{j}:=R_{j}^{-1} C_{0, R_{j}}^{2}
$$

Since mes $E_{j} \rightarrow 0$ it is possible to choose a sequence (for which we keep the same notation) such that $\chi_{j} \rightarrow 0$ almost everywhere. We will also suppose that $R_{j}$ are the same for $\chi_{j}$ and $u_{j}$. Thus

$$
\int_{|\zeta| \leq 1}\left|\chi_{j}(\zeta) u_{j}(\zeta)-0 \cdot v(\zeta)\right| d \zeta=\int_{|\zeta| \leq 1}\left|\chi_{j}(\zeta) u_{j}(\zeta)\right| d \zeta \rightarrow 0
$$

By change of variables $z=R_{j} \zeta$ we obtain that

$$
R_{j}^{-\rho\left(R_{j}\right)-2} \int_{C_{0, R_{j}}^{2}}|u|(z) d z=\int_{|\zeta| \leq 1}\left|\chi_{j}(\zeta) u_{j}(\zeta)\right| d \zeta \rightarrow 0
$$

Hence the limit in (5.2.2.2) is equal to zero. Contradiction.
Proposition 5.2.2.2 Under condition (5.2.1.2) the set

$$
C:=\left\{z:|u(z)||z|^{-\rho(|z|)}>\epsilon\right\}
$$

is a $C_{0}^{0}$-set for arbitrary $\epsilon$.
Proof. Assume the contrary; that is , $\exists \alpha>0$ such that

$$
\begin{equation*}
\alpha-\overline{\mathrm{mes}} C=2 \delta>0 \tag{5.2.2.3}
\end{equation*}
$$

One can see that for some $\eta>0$,

$$
\begin{equation*}
\limsup _{R \rightarrow \infty}\left(\alpha-\operatorname{mes} K_{\eta R}\right) R^{-\alpha} \leq \delta / 2 \tag{5.2.2.4}
\end{equation*}
$$

Exercise 5.2.2.1 Check this.
(5.2.2.3) and (5.2.2.4) imply that there exists a sequence $R_{j} \rightarrow \infty$ such that

$$
\lim _{R_{j} \rightarrow \infty} \alpha-\operatorname{mes}\left[C \cap\left(K_{R_{j}} \backslash K_{\eta R_{j}}\right)\right] R_{j}^{-\alpha} \geq \frac{3}{2} \delta
$$

Set

$$
E_{j}:=R_{j}^{-1} C \cap\left(K_{R_{j}} \backslash K_{\eta R_{j}}\right) .
$$

It is clear that $E_{j} \subset K_{1} \backslash K_{\eta}$ and for sufficiently large $j$,

$$
\begin{equation*}
\alpha-\operatorname{mes} E_{j} \geq \delta \tag{5.2.2.5}
\end{equation*}
$$

Set $u_{j}$ as in (5.2.2.2a). We claim that for large $j$ and $\zeta \in E_{j}$,

$$
\begin{equation*}
\left|u_{j}\right|(\zeta) \geq \frac{\epsilon}{2}|\zeta|^{\rho} \tag{5.2.2.6}
\end{equation*}
$$

Indeed,

$$
\left|u_{j}\right|(\zeta)=\frac{|u|\left(R_{j} \zeta\right)}{R_{j}^{\rho\left(R_{j}\right)}}=\frac{|u|(z)}{|z|^{\rho(|z|)}}(1+o(1))|\zeta|^{\rho} \geq \frac{\epsilon}{2}|\zeta|^{\rho} .
$$

We used here properties of the proximate order and the equivalence

$$
z=R_{j} \zeta \in C \cap\left(K_{R_{j}} \backslash K_{\eta R_{j}}\right) \Leftrightarrow \zeta \in E_{j} .
$$

Exercise 5.2.2.2 Check this in detail.
Now we will show that the condition (5.2.1.2) contradicts (5.2.2.6). Since $u \in$ $\delta S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$ it is a difference of $u_{1}, u_{2} \in S H\left(\rho, \rho(r), \mathbb{R}^{2}\right)$. The corresponding sequences $u_{1, j}$ and $u_{2, j}$ are precompact in $\mathcal{D}^{\prime}$ and there exist subsequences (with the same notation) that converge to $v_{1}$ and $v_{2}$, respectively.

By Theorem 2.7.5.1 these sequences converge to $v_{1}$ and $v_{2}$ with respect to $\alpha$ - mes on $K_{1} \backslash K_{\eta}$. Since $u_{t} \rightarrow 0$ in $\mathcal{D}^{\prime}$, it follows that $v_{1}=v_{2}$. Thus $u_{j} \rightarrow 0$ with respect to $\alpha-$ mes on $K_{1} \backslash K_{\eta}$. However, this contradicts (5.2.2.5) and (5.2.2.6).

Proposition 5.2.2.3 Let $\left\{C_{j}\right\}_{1}^{\infty}$ be a sequence of $C_{0}^{0}$-sets. There exists a sequence $R_{j} \rightarrow \infty$ such that the set

$$
\begin{equation*}
C=\bigcup_{j=1}^{\infty}\left\{C_{j} \cap\left(K_{R_{j+1}} \backslash K_{R_{j}}\right)\right\} \tag{5.2.2.7}
\end{equation*}
$$

is a $C_{0}^{0}$-set.

Proof. Choose $\epsilon_{j} \downarrow 0$ and $\alpha_{j} \downarrow 0$. Set $R_{0}:=1$. Suppose $R_{j-1}$ was already chosen. Take $R_{j}$ such that

$$
\begin{equation*}
\alpha_{j}-\operatorname{mes}\left[C \cap K_{R_{j-1}}\right]<\epsilon_{j} R^{\alpha_{j}} \tag{5.2.2.8}
\end{equation*}
$$

for $R>R_{j}$.
It is possible because of property rC 1 ) Theorem 5.2.1.1. We can also increase $R_{j}$ so that

$$
\begin{equation*}
\alpha_{j}-\operatorname{mes}\left[C_{j} \cap K_{R}\right]<\epsilon_{j} R^{\alpha_{j}} \tag{5.2.2.9}
\end{equation*}
$$

and

$$
\begin{equation*}
\alpha_{j}-\operatorname{mes}\left[C_{j+1} \cap K_{R}\right]<\epsilon_{j} R^{\alpha_{j}} \tag{5.2.2.10}
\end{equation*}
$$

for $R>R_{j}$.
It is possible because $C_{j}$ and $C_{j+1}$ are $C_{0}^{0}$-sets.
Let us estimate $\alpha_{j}-\operatorname{mes}\left[C \cap K_{R}\right]$ for $R_{j} \leq R<R_{j+1}$. From (5.2.2.8), (5.2.2.9) and (5.2.2.10) we obtain

$$
\begin{equation*}
\alpha_{j}-\operatorname{mes}\left[C \cap K_{R}\right] \leq 3 \epsilon_{j} R^{\alpha_{j}} . \tag{5.2.2.11}
\end{equation*}
$$

Let $\alpha>0$ be arbitrarily small. Find $\alpha_{j}<\alpha$. For $R_{j+1} \geq R>R_{j}$ we have

$$
\alpha-\operatorname{mes}\left[C \cap K_{R}\right] R^{-\alpha} \leq \alpha_{j}-\operatorname{mes}\left[C \cap K_{R}\right] R^{-\alpha_{j}} \leq 3 \epsilon_{j} .
$$

Hence $\alpha-\overline{\mathrm{mes}} C=0$.

### 5.2.3

Proof of Theorem 5.2.1.2. Let $\phi \in \mathcal{D}(\mathbb{C})$ and $\operatorname{supp} \phi \subset K_{R}$. Then for any $\epsilon>0$,

$$
\begin{equation*}
J(t):=\int \phi(z) u_{t}(z) d z=\left(\int_{K_{R} \backslash K_{\epsilon}}+\int_{K_{\epsilon}}\right) \phi(z) u_{t}(z) d z:=J_{1}(t)+J_{2}(t) \tag{5.2.3.1}
\end{equation*}
$$

We have for $J_{2}$ (see 2.8.2.3):

$$
\begin{equation*}
\left|J_{2}\right|(t) \leq \max _{|z| \leq \epsilon}|\phi(z)| \times \text { const } \int_{0}^{\epsilon} T\left(r,\left|u_{t}\right|\right) r d r \leq \operatorname{const} T\left(\epsilon,\left|u_{t}\right|\right) \epsilon^{2} \tag{5.2.3.2}
\end{equation*}
$$

Further (see Theorem 2.8.2.1)

$$
\begin{align*}
T\left(r,\left|u_{t}\right|\right) & \leq 2 T\left(r, u_{t}\right)+O\left(t^{-\rho(t)}\right) \leq 2\left[T\left(r, u_{1, t}\right)+T\left(r, u_{2, t}\right)\right]+O\left(t^{-\rho(t)}\right) \\
& \leq 2\left[M\left(r, u_{1, t}\right)+M\left(r, u_{2, t}\right)\right]+O\left(t^{-\rho(t)}\right) \tag{5.2.3.3}
\end{align*}
$$

Using (5.2.3.2), (5.2.3.3) and (3.1.2.3) we obtain

$$
\begin{equation*}
\limsup _{t \rightarrow \infty}\left|J_{2}(t)\right| \leq \operatorname{const} \epsilon^{\rho+2} \tag{5.2.3.4}
\end{equation*}
$$

Exercise 5.2.3.1 Check this using the change of variable $z=t \zeta$.
To estimate $J_{1}(t)$ write

$$
\begin{equation*}
\left|J_{1}(t)\right| \leq \mathrm{const}\left(\int_{\tilde{K}_{t} \backslash C_{0, R t}^{2}}|u(z)| d z+\int_{C_{0, R t}^{2}}|u(z)| d z\right) t^{-\rho(t)-2}:=J_{1,1}(t)+J_{1,2}(t) \tag{5.2.3.5}
\end{equation*}
$$

where $\tilde{K}_{t}:=\{z: \epsilon t \leq|z| \leq R t\}$.
The summand $J_{1,1}$ is $o(1)$ as $t \rightarrow \infty$ by (5.2.1.3).
Exercise 5.2.3.2 Check this using the properties of the proximate order (Theorem 2.8.1.3, ppo3).

The summand $J_{1,2}$ is $o(1)$ by Theorem 5.2.2.1. Thus

$$
\limsup _{t \rightarrow \infty}|J(t)| \leq \operatorname{const} \epsilon^{\rho+2}
$$

for any $\epsilon$. Hence it is equal to zero and the sufficiency of (5.2.1.3) has been proved.
Let us prove sufficiency of (5.2.1.2). Let $\epsilon_{j} \downarrow 0$. By Theorem 5.2.2.2 we choose a $C_{0}^{0}$-set $C_{j}$ outside which $|u(z) \| z|^{-\rho(|z|)}<\epsilon_{j}$.

We construct the set $C$ by (5.2.2.7). Outside $C$ we have (5.2.1.3). And by Theorem 5.2.2.3 it is a $C_{0}^{0}$-set.

### 5.3 Asymptotic approximation of subharmonic functions

5.3.1 One of the widely applied methods of constructing entire functions with a prescribed asymptotic behavior is the following: First construct a subharmonic function behaving asymptotically as the logarithm of modulus of the entire function, and then approximate it in some sense by the logarithm of modulus of entire function such that the asymptotic is preserved.

Various queries about the a precision of the approximation and about the metric in which it was implemented generated a spectrum of theorems of such kind that we will demonstrate.

Historically the first theorems of this kind were proved for concrete functions, the masses of which were concentrated on sufficiently smooth curves (in particular, on lines, see, e.g., $[\mathrm{BM}, \mathrm{Ev}, \mathrm{Kj}, \mathrm{Ar}], \ldots$ )

In such cases the approximation was very precise and exceptional sets where the approximation failed were small and determined.

The first general case was proved in $[\mathrm{Az}(1969)]$. Next it was developed in $[\mathrm{Yu}(1982)]$, and vastly improved in $[\mathrm{Yu}(1985)]$. It is the following

Theorem 5.3.1.1 (Yulmukhametov) Let $u \in S H(\rho)$. Then there exists an entire function $f$ such that for every $\alpha \geq \rho$,

$$
|u(z)-\log | f(z)\left|\left|<C_{\alpha} \log \right| z\right|
$$

for $z \notin E_{\alpha}$, where $E_{\alpha}$ is an exceptional set that can be covered by discs $D_{z_{j}, r_{j}}:=$ $\left\{z:\left|z-z_{j}\right|<r_{j}\right\}$ satisfying the condition

$$
\sum_{\left|z_{j}\right|>R} r_{j}=o\left(R^{\rho-\alpha}\right), R \rightarrow \infty .
$$

This theorem is precise in the following sense: If

$$
\|z|-\log | f(z)\|=o(\log |z|), z \rightarrow \infty, z \notin E
$$

then for every covering of $E$ by discs $D_{z_{j}, r_{j}}$ and every $\epsilon>0$

$$
\sum_{\left|z_{j}\right|<R} r_{j} \geq R^{1-\epsilon}, R \rightarrow \infty
$$

i.e., in any case this sum is not even bounded.

However it is necessary to remark that the construction from $[\mathrm{Yu}(1985)]$ "rigidly" fastens zeros of the entire function, whereas the construction of $[\mathrm{Az}(1969)]$ and $[\mathrm{Yu}(1982)]$ gives some possibilities to move them, which is needed in some constructions.

Let us also mention that such approximation generates an approximation of a plurisubharmonic function by the logarithm of the modulus of an entire function in $\mathbb{C}^{p}($ see $[\mathrm{Yu}(1996)])$.

It is also useful to approximate subharmonic functions in an integral metric, for example $L^{p}$, as was done in [GG].

Set

$$
\|g\|_{p}:=\left(\int_{0}^{2 \pi}|g(t)|^{p} d t\right)^{1 / p}
$$

Denote by $Q(r, u)$ a function that satisfies the conditions:

1) if $u$ is of finite order, then $Q(r, u)=O(\log r)$;
2) if $u$ is of infinite order, then $Q(r, u)=O\left(\log r+\log \mu_{u}(r)\right)$.

Theorem 5.3.1.2 (Girnyk, Gol'dberg) For every subharmonic function in $\mathbb{C} u$ there exists an entire function $f$ such that $\left\|u\left(r e^{i \cdot}\right)-\log |f|\left(r e^{i \cdot}\right)\right\|_{p}=Q(r, u)$.

This theorem also considers functions of infinite order. In this case, it is possible to replace $\mu_{u}(r)$ by $T(r, u)$ or $M(r, u)$ in $Q(r, u)$ outside an exceptional set $E \subset \mathbb{R}^{+}$of finite measure. This theorem is also unimprovable for subharmonic
functions of finite order, because, for example, $u=\frac{1}{2} \log |z|$ gives, as it is possible to prove:

$$
\liminf _{r \rightarrow \infty} \frac{\left\|u\left(r e^{i \bullet}\right)-\log |f|\left(r e^{i \bullet}\right)\right\|_{p}}{\log r}>0
$$

However it was found [LS], $[\mathrm{LM}]$ that the remainder term $O(\log |z|)$ that was regarded the best possible is not precise and in some "regular" cases can be replaced with $O(1)$ outside a bigger (but still "small") set .

Set, for $E \subset \mathbb{C}$ :

$$
\Delta(E):=\limsup _{r \rightarrow \infty} \frac{\operatorname{mes} E \cap D_{0, R}}{R^{2}}
$$

Theorem 5.3.1.3 (Lyubarskii, Malinnikova) Let u be a subharmonic function in $\mathbb{C}$ with $\mu_{u}$ satisfying the conditions: $\mu_{u}(\mathbb{C})=\infty$ and there exists $\alpha>0, q>1, R_{0}>$ 0 such that

$$
\mu_{u}\left(D_{0, q R} \backslash D_{0, R}\right)>\alpha
$$

for all $R>R_{0}$.
Then there exists an entire function $f$ such that for every $\epsilon>0$,

$$
|u(z)-\log | f(z)\left|\mid<C_{\epsilon}\right.
$$

for $z \in \mathbb{C} \backslash E_{\epsilon}$ with $\Delta\left(E_{\epsilon}\right)<\epsilon$.
So if $\mu_{u}$ has no "Hadamard's gaps" such approximation is possible.
In this book we restrict ourself to a weaker and simply proved theorem that is sufficient for our aim

Theorem 5.3.1.4 (Approximation Theorem) For every $u \in S H(\rho, \rho(r))$ there exists an entire function $f$ such that

$$
D^{\prime}-\lim _{t \rightarrow \infty}(u-\log |f|)_{t}=0
$$

Nevertheless this theorem has an important
Corollary 5.3.1.5 For every $u \in S H(\rho, \rho(r))$ there exists an entire function $f$ such that

$$
\mathbf{F r}[u]=\mathbf{F r}[f] .
$$

5.3.2 Now we prove Theorem 5.3.1.4. We can suppose, because of Theorem 3.1.6.1 (Dependence $\mathbf{F r}$ on $V_{\bullet}$ ), that in the definition of $(\bullet)_{t}\left(\right.$ see 3.1.2.1) $V_{t} \equiv I$

We prove this theorem for the case non-integer $\rho$. For proving this theorem we need

Lemma 5.3.2.1 Let $u \in \delta S H(\rho, \rho(r))$, for $\rho$ non-integer, and $\nu$ is its charge. Then $u_{t} \rightarrow 0$ iff $\nu_{t} \rightarrow 0$ in $\mathcal{D}^{\prime}$ as $t \rightarrow \infty$.

Proof. Sufficiency. Suppose $u_{t}:=\left(u_{1}\right)_{t}-\left(u_{2}\right)_{t} \nrightarrow 0$. There exists a subsequence $t_{j} \rightarrow \infty$ and subharmonic functions $v_{1}$ and $v_{2}$ such that

$$
\begin{equation*}
u_{t_{j}}=\left(u_{1}\right)_{t_{j}}-\left(u_{2}\right)_{t_{j}} \rightarrow v_{1}-v_{2}:=v \neq 0 . \tag{5.3.2.1}
\end{equation*}
$$

Applying to (5.3.2.1) the continuity of $\Delta$ in $\mathcal{D}^{\prime}$ and using the conditions of the theorem, we obtain

$$
\nu_{t_{j}} \rightarrow \frac{1}{2 \pi} \Delta v=0
$$

Hence $v$ is harmonic. Since $v_{1}, v_{2} \in U[\rho$,$] also v \in U[\rho]$ (see Theorem 2.8.2.1, t3), t4) and Theorem 2.8.2.3).

Exercise 5.3.2.1 Prove this in detail.
By Theorem 3.1.4.3 we obtain $v=0$. Contradiction.
Necessity. Since the Laplace operator is continuous in $\mathcal{D}^{\prime}$-topology, the assertion $u_{t} \rightarrow 0$ implies $\nu_{t}:=\frac{1}{2 \pi} \Delta u_{t} \rightarrow 0$.

Now we describe a construction of the zero distribution of the future entire function. Let $u \in S H(\rho)$ and $\mu$ be its mass distribution. Set

$$
\begin{equation*}
R_{j+1}:=R_{j}(j+1)^{4 / \kappa} \tag{5.3.2.2}
\end{equation*}
$$

where $\kappa:=\min (\rho-[\rho],[\rho]+1-\rho)$.
Let us divide all the plane by circles of the form $S_{R_{j}}:=\left\{|z|=R_{j}\right\}$ such that $R_{j+1} / R_{j} \rightarrow \infty$ and $\mu\left(S_{R_{j}}\right)=0$.

Exercise 5.3.2.2 Prove that it is possible.
Choose a sequence $\delta_{j} \downarrow 0$. Divide every annulus $K_{j}:=\left\{z: R_{j} \leq|z|<R_{j+1}\right\}$ by circles $S_{R_{j, n}}$ for

$$
R_{j, n}:=\left(\frac{1+\delta_{j}}{1-\delta_{j}}\right)^{n} R_{j}, n=0,1,2, \ldots, n_{j}
$$

where

$$
n_{j}:=\left[\frac{\log \frac{R_{j+1}}{R_{j}}}{\log \frac{1+\delta_{j}}{1-\delta_{j}}}\right],
$$

and by rays

$$
L_{k}:=z: \arg z=k \delta_{j}, k=0,1, \ldots,\left[2 \pi / \delta_{j}\right]
$$

They divide all the plane into sectors $K_{j, n, k}$. We can choose $\delta_{j}$ in such a way that $\mu\left(\partial K_{j, n, k}\right)=0$ because $\mu\left(K_{j, n, k}\right)$ is a monotonic function of $\delta_{j}$ and has only a countable set of jumps.

Exercise 5.3.2.3 Explain this in detail.

Choose a point $z_{j, n, k}$ in every sector $K_{j, n, k}$ and concentrate all the mass of the sector at this point. In other words we consider a new mass distribution $\hat{\mu}$ that has masses concentrated in the points $z_{j, n, k}$ and $\hat{\mu}\left(z_{j, n, k}\right)=\mu\left(K_{j, n, k}\right)$.

The next lemma shows that $\hat{\mu}$ is close to $\mu$.
Lemma 5.3.2.2 One has

$$
\hat{\mu}_{t}-\mu_{t} \rightarrow 0
$$

in $\mathcal{D}^{\prime}$ as $t \rightarrow \infty$.

Proof. Assume the contrary, i.e., $\hat{\mu}_{t}-\mu_{t} \nrightarrow 0$. Choose a sequence $t_{l} \rightarrow \infty$ such that $\hat{\mu}_{t_{l}} \rightarrow \hat{\nu}$ and $\mu_{t_{l}} \rightarrow \nu, \nu, \hat{\nu} \in \mathcal{M}[\rho], \nu \neq \hat{\nu}$. Then there exists a disc $K_{z_{0}, r_{0}}:=$ $\left\{z:\left|z-z_{0}\right|<r_{0}\right\}$ such that $\nu\left(K_{z_{0}, r_{0}}\right) \neq \hat{\nu}\left(K_{z_{0}, r_{0}}\right)$. We can assume that this disc does not contain zero since for all the $\nu \in \mathcal{M}[\rho]$ the condition $\nu\left(K_{r}\right) \leq \Delta r^{\rho}, \forall r>0$ is fulfilled.

Suppose, for example,

$$
\begin{equation*}
\nu\left(K_{z_{0}, r_{0}}\right)>\hat{\nu}\left(K_{z_{0}, r_{0}}\right) \tag{5.3.2.3}
\end{equation*}
$$

Set $a:=\nu\left(K_{z_{0}, r_{0}}\right)-\hat{\nu}\left(K_{z_{0}, r_{0}}\right)>0$. Choose $\epsilon$ such that

$$
\begin{equation*}
\nu\left(K_{z_{0}, r_{0}}\right)<\nu\left(\overline{K_{z_{0}, r_{0}-\epsilon}}\right)+a / 3 . \tag{5.3.2.4}
\end{equation*}
$$

This is possible because the countable additivity of $\hat{\nu}$ implies $\lim _{r^{\prime} \uparrow r} \nu\left(K_{z_{0}, r^{\prime}}\right)=$ $\nu\left(K_{z_{0}, r}\right)$.

Consider now the sets $t_{l} K_{z_{0}, r_{0}}, t_{l} K_{z_{0}, r_{0}-\epsilon}$. For sufficiently large $t_{l}$ they are contained in the union of the annuluses $K_{j_{l}} \cup K_{j_{l}+1}$.

As $j_{l} \rightarrow \infty$ the diameters of all the sectors $K_{j_{l}, n, k}$ are $o\left(R_{j_{l}}\right)$ uniformly. Thus they are $o\left(t_{l}\right)$. Hence for such $t_{l}$ 's we can find a union $\Gamma_{l}$ of sectors covering $t_{l} K_{z_{0}, r_{0}-\epsilon}$ that does not intersect the circle of $t_{l} K_{z_{0}, r_{0}}$.

We have $\hat{\mu}\left(\Gamma_{l}\right)=\mu\left(\Gamma_{l}\right)$ by definition of $\hat{\mu}$. Using the monotonicity of measures, we obtain $\mu\left(t_{l} K_{z_{0}, r_{0}-\epsilon}\right) \leq \hat{\mu}\left(t_{l} K_{z_{0}, r_{0}}\right)$, whence

$$
\mu_{t_{l}}\left(K_{z_{0}, r_{0}-\epsilon}\right) \leq \hat{\mu}_{t_{l}}\left(K_{z_{0}, r_{0}}\right)
$$

Passing to the limit as $l \rightarrow \infty$ and using Theorems 2.2.3.1 and 2.3.4.4, we obtain $\nu\left(\overline{K_{z_{0}, r_{0}-\epsilon}}\right) \leq \hat{\nu}\left(K_{z_{0}, r_{0}}\right)$. Using (5.3.2.4), we obtain $\nu\left(K_{z_{0}, r_{0}}\right)-1 / 3\left[\nu\left(K_{z_{0}, r_{0}}\right)-\right.$ $\left.\hat{\nu}\left(K_{z_{0}, r_{0}}\right)\right] \leq \hat{\nu}\left(K_{z_{0}, r_{0}}\right)$ and hence $\nu\left(K_{z_{0}, r_{0}}\right) \leq \hat{\nu}\left(K_{z_{0}, r_{0}}\right)$, that contradicts (5.3.2.3). Since $\nu$ and $\hat{\nu}$ are symmetric in this reasoning the lemma is proved.

Let us finish the proof of Theorem 5.3.1.4 for non-integer $\rho$.
We construct a distribution $n$ with integer masses concentrated at points $z_{j, k, n}$. Set

$$
n\left(z_{j, k, n}\right):=\left[\hat{\mu}\left(z_{j, k, n}\right)\right]
$$

and estimate the growth of the difference

$$
\delta \mu:=\hat{\mu}-n
$$

that is also a mass distribution concentrated at the same points.
Since

$$
\delta \mu\left(z_{j, k, n}\right) \leq 1
$$

it is sufficient to count the number of points in the disc $K_{R}$.
The number of points in the annulus $\left\{R_{j} \leq|z|<R\right\}$ is found from (5.3.2.2),

$$
\begin{aligned}
\delta \mu\left(\left\{R_{j} \leq|z|<R\right\}\right) & \leq\left[\log \left(\frac{1+\delta_{j}}{1-\delta_{j}}\right)\right]^{-1} \frac{2 \pi}{\delta_{j}} \log \frac{R}{R_{j}} \\
& \leq \operatorname{const} \times \frac{\log (j+1)}{\delta_{j}^{2}}=\mathrm{const} \times(j+1)^{4} \log (j+1)
\end{aligned}
$$

The mass of the disc $K_{R}$ is estimated by the inequality

$$
\begin{equation*}
\delta \mu\left(K_{R}\right) \leq \mathrm{const} \times \sum_{k=0}^{n-1}(k+1)^{4} \log (k+1)=o\left(n^{6}\right)=o\left(R^{\epsilon}\right) \tag{5.3.2.5}
\end{equation*}
$$

for any $\epsilon>0$ because $R>R_{n-1}=((n-1)!)^{4 / \kappa}$.
Exercise 5.3.2.4 Check this in detail.
The estimate (5.3.2.5) shows that

$$
\begin{equation*}
\delta \mu_{t} \rightarrow 0 \tag{5.3.2.6}
\end{equation*}
$$

as $t \rightarrow \infty$.
Lemma 5.3.2.2 and (5.3.2.6) imply that

$$
\begin{equation*}
\mu_{t}-n_{t} \rightarrow 0 \tag{5.3.2.7}
\end{equation*}
$$

Set

$$
u_{1}(z):=\Pi(z, n, p)
$$

(see (2.9.2.1)) where $\Pi$ is a canonical potential. This is a subharmonic function in the plane with integral masses. Thus it is the logarithm of the modulus of the entire function

$$
f(z)=\prod E\left(z / z_{j, k, n}\right)
$$

(5.3.2.7) implies by Lemma 5.3.2.1 that $u_{t}-\left(u_{1}\right)_{t} \rightarrow 0$ and this is the assertion of Theorem 5.3.1.4 for non-integer $\rho$.

### 5.4 Lower indicator of A.A. Gol'dberg. Description of lower indicator Description of the pair: indicator-lower indicator

5.4.1 Now we consider the lower indicator. For an entire function of finite order $\rho$ and normal type it can be defined in one of the following ways:

$$
\begin{equation*}
\underline{h}_{1}(\phi, f):=\sup _{C \in \mathcal{C}}\left\{\liminf _{r e^{i \phi} \rightarrow \infty, r e^{i \phi} \notin C} \log \left|f\left(r e^{i \phi}\right)\right| r^{-\rho(r)}\right\}, \tag{5.4.1.1}
\end{equation*}
$$

where $\mathcal{C}$ is the set of $C_{0}$-sets (see [Le, Ch. II, $\left.\S 1\right]$ ), i.e., the sets that can be covered by a union of discs $K_{\delta_{j}}\left(z_{j}\right):=\left\{z:\left|z-z_{j}\right|<\delta_{j}\right\}$ such that

$$
\lim _{R \rightarrow \infty} \frac{1}{R} \sum_{\left|z_{j}\right|<R} \delta_{j}=0
$$

The exclusion of $C_{0}$-sets is necessary because we must exclude from our consideration some neighborhoods of roots of $f(z)$ where $\log |f(z)|$ is near $-\infty$.

Similarly, define

$$
\begin{equation*}
\underline{h}_{2}(\phi, f):=\sup _{E(\phi) \in \mathcal{E}}\left\{\liminf _{r \rightarrow \infty, r \notin E(\phi)} \log \left|f\left(r e^{i \phi}\right)\right| r^{-\rho(r)}\right\} \tag{5.4.1.2}
\end{equation*}
$$

where $\mathcal{E}$ is the set of $E_{0}$-sets (see [Le, Ch. III]), i.e., sets $E \subset(0, \infty]$ satisfying the condition

$$
\lim _{R \rightarrow \infty} \operatorname{mes}\{E \cap(0, R)\} R^{-1}=0
$$

The definition (5.4.1.1) was introduced by A.A. Gol'dberg (see [Go(1967)]). We will use the definition (3.2.1.2)

$$
\begin{equation*}
\underline{h}(\phi, f)=\inf \left\{v\left(e^{i \phi}\right): v \in \mathbf{F r}[f]\right\} . \tag{5.4.1.3}
\end{equation*}
$$

It was proved in [AP, Thm. 1] that the definitions (5.4.1.1), (5.4.1.2) and (5.4.1.3) coincide.

Let us note that (5.4.1.3) uses the definition (3.2.1.2) only on the circle $\{|z|=1\}$. However, it is easy to check, by using Theorem 3.2.1.2 that for $\underline{h}(z)=$ $|z|^{\rho} \underline{h}(\arg z)$ properties h1) and h2), Theorem 3.2.1.1, are preserved.
Exercise 5.4.1.1 Check this.
We are going to prove
Theorem 5.4.1.1 Let $g(\phi)$ be a $2 \pi$-periodic function that is either semicontinuous from above or $\equiv-\infty$ and $\rho(r) \rightarrow \rho$ be an arbitrary approximate order. Then there exists an entire function $f \in A(\rho, \rho(r))$ such that

$$
\begin{equation*}
\underline{h}(\phi, f)=g(\phi) \tag{5.4.1.4}
\end{equation*}
$$

for all $\phi \in[0,2 \pi)$.
5.4.2 We will use the following assertion that is a corollary of Theorem 4.3.5.1 and Corollary 5.3.1.5:

Theorem 5.4.2.1 Let $\Lambda \subset U[\rho]$ be a compact, connected and $T_{\bullet}$-invariant subset. Then for any proximate order $\rho(r) \rightarrow \rho$ there exists $f \in A(\rho, \rho(r))$ such that

$$
\begin{align*}
h(\phi, f) & =\sup \left\{v\left(e^{i \phi}\right): v \in \Lambda\right\}  \tag{5.4.2.1}\\
\underline{h}(\phi, f) & =\inf \left\{v\left(e^{i \phi}\right): v \in \Lambda\right\} \tag{5.4.2.2}
\end{align*}
$$

Exercise 5.4.2.1 Prove Theorem 5.4.2.1.
For the sake of clarity let us restrict ourselves to non-integer $\rho$. We will construct a set $\Lambda$ such that

$$
\inf \left\{v\left(e^{i \phi}\right): v \in \Lambda\right\}=g(\phi)
$$

Denote

$$
\begin{aligned}
H(z, p) & :=\log |1-z|+\Re \sum_{k=1}^{p} \frac{z^{k}}{k} ; p=[\rho], \\
\gamma(z, K, \lambda) & :=-\lambda+K|z-1|, \lambda, K \geq 0 .
\end{aligned}
$$

Note the following properties of these functions:
a) $\min _{|z-1| \geq \delta} \delta H(z, p)|z|^{-\rho} \rightarrow 0$, as $\delta \rightarrow 0$;
b) $\delta H(z, p)|z|^{-\rho} \leq A \delta$, for all $z \in \mathbb{C}$, where $A$ depends only on $p$;
c) $\max _{|z-1| \leq 0.5 K} \gamma(z, K, \lambda) \leq-\frac{1}{2}$.

Exercise 5.4.2.2 Prove properties a), b), c).
Let us note that $H(1, p)=-\infty$. Consider the family:

$$
\Lambda_{\infty}=\left\{v_{\theta, \tau}(z):=H\left(z e^{-i \theta} \tau, p\right) \tau^{-\rho}: \theta \in[0,2 \pi), \tau \in(0, \infty)\right\} \cup 0
$$

This family is contained in $U[\rho]$ because of b) and closed in $\mathcal{D}^{\prime}$-topology. It is also $T_{\bullet}$-invariant, hence, satisfies the conditions of Theorem 5.4.2.1. For every $\phi \in[0,2 \pi)$ there exists $\theta_{0}(=\phi)$, and $\tau_{0}(=1)$ such that $v_{\theta_{0}, \tau_{0}}\left(e^{i \phi}\right)=H(1, p)=-\infty$ Hence

$$
\begin{equation*}
\inf \left\{v\left(e^{i \phi}\right): v \in \Lambda_{\infty}\right\}=-\infty \tag{5.4.2.4}
\end{equation*}
$$

For the general case this construction will be improved, cutting the "trunk" of the function $H\left(z e^{-i \theta}, p\right)$.

Take $\delta$ small enough so that the following conditions hold:

$$
\begin{align*}
\delta H(z, p)|z|^{-\rho} & \geq-\frac{1}{4}, \text { for }|z-1| \geq \delta,  \tag{5.4.2.5}\\
\delta H(z, p) & \geq-\frac{1}{4}, \text { for }|z-1|=\delta,  \tag{5.4.2.6}\\
\delta & \leq \frac{1}{2 K} . \tag{5.4.2.7}
\end{align*}
$$

Then

$$
\begin{equation*}
\delta H(z, p)>\gamma(z, K, \lambda), \text { for }|z-1|=\delta . \tag{5.4.2.8}
\end{equation*}
$$

Denote

$$
W(z, K, \delta, \lambda):= \begin{cases}\max \{\delta H(z, p), \gamma(z, K, \lambda)\}, & \text { for }|z-1|<\delta  \tag{5.4.2.9}\\ \delta H(z, p), & \text { for }|z-1| \geq \delta\end{cases}
$$

Lemma 5.4.2.2 The following holds:
aw) The function $W(z, K, \delta, \lambda)$ is subharmonic in $\mathbb{C}$;
bw) $\operatorname{supp} \mu_{W} \Subset\{|z-1|<\delta\}$;
cw) $\quad \sup _{z \in \mathbb{C}} W(z, \bullet, \delta, \lambda)|z|^{-\rho} \leq A \delta$,
where $A$ depends only on $p$.
Proof. For $|z-1|<\delta, W$ is subharmonic as the maximum of two subharmonic functions. For $|z-1| \geq \delta$ it is harmonic even in the neighborhood of the circle $|z-1|=\delta$, because of inequality (5.4.2.8). So aw) and bw) hold. The assertion $\mathrm{cw})$ follows from b) and c) (5.4.2.3) above.

Now we get to the proof of (5.4.1.4). Let $g_{n} \downarrow g$ be a sequence of continuously differentiable functions that converges to $g$ monotonically. This is possible, because $g$ is semicontinuous from above.

Exercise 5.4.2.3 Prove that Theorem 2.1.2.9 and the Weierstrass theorem of approximation of every periodic function by trigonometrical polynomials imply the last assertion.

We write

$$
M_{n}:=\max _{\phi} g_{n}^{+}(\phi)
$$

where as usual $a^{+}=\max (a, 0)$. Set

$$
v_{\theta, n}(z):=W\left(z e^{-i \theta}, K_{n}, \delta_{n}, M_{n}+1-g_{n}(\theta)\right)+\left(M_{n}+1\right)|z|^{\rho},
$$

where $\delta_{n}$ is chosen small and $K_{n}$ is chosen large. Set $z=\tau e^{i \phi}$. It is clear that

$$
\begin{equation*}
v_{\phi, n}\left(e^{i \phi}\right)=g_{n}(\phi) \tag{5.4.2.11}
\end{equation*}
$$

for all $K_{n}, \delta_{n}$.

We can choose $K_{n}$ so large and $\delta_{n}$ so small that

$$
\gamma\left(z, K_{n}, M_{n}+1-g(\theta)\right)|z|^{-\rho} \geq g_{n}(\phi)
$$

for $|z-1| \leq \delta_{n}$, because $g_{n}$ has bounded derivative.
After that we can make $\delta_{n}$ smaller so that for $|z-1| \geq \delta_{n}$ the inequality (5.4.2.8) would hold.

Exercise 5.4.2.4 Estimate exactly $K_{n}$ and $\delta_{n}$ via the derivative of $g_{n}$.
Then

$$
v_{\theta, n}(z)|z|^{-\rho} \geq g_{n}(\phi)
$$

for all $z=r e^{i \phi}$. Thus

$$
\min _{\theta, \tau} v_{\theta, n}\left(\tau e^{i \phi}\right) \tau^{-\rho}=g_{n}\left(e^{i \phi}\right)
$$

and the minimum is attained for $\tau=1, \theta=\phi$.
Let us note that from (5.4.2.10) we have

$$
\sup _{\theta} \sup _{z \in \mathbb{C}} v_{\theta, n}(z)|z|^{-\rho} \leq A \delta_{n}+M_{n}+1 \leq A+M_{1}+1
$$

Consider now the family of functions

$$
\Lambda_{0}:=\left\{v_{\theta, n}(z \tau)|\tau|^{-\rho}: \theta \in[0 ; 2 \pi), n=1,2, \ldots, \tau \in(0 ; \infty)\right\}
$$

It is contained in $U[\rho, \sigma]$ for $\sigma=A+M_{1}+1$ and is $T_{\bullet}$-invariant. Let $\Lambda$ be its closure in $\mathcal{D}^{\prime}$. Let us show that

$$
\begin{equation*}
g(\phi)=\inf \left\{v\left(e^{i \phi}\right): v \in \Lambda\right\} . \tag{5.4.2.12}
\end{equation*}
$$

Indeed, for every sequence $v_{j} \in \Lambda_{1}$

$$
v_{j}\left(e^{i \phi}\right) \geq \inf _{n} g_{n}(\phi)=g(\phi)
$$

Let $v \in \Lambda$. By Theorem 2.7.4.1 ( $\mathcal{D}^{\prime}$ and Quasi-everywhere Convergence)

$$
v(z):=\left(\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} v_{j}\right)(z)=\left(\limsup _{j \rightarrow \infty} v_{j}\right)^{*}(z)
$$

Hence

$$
v\left(e^{i \phi}\right) \geq g(\phi)
$$

However, the infimum is attained for every $\phi$ on the sequence $v_{\phi, n}(z)$ because of (5.4.2.11). Hence (5.4.2.12) holds and Theorem 5.4.1.4 is proved.
5.4.3 Now we describe the pair: indicator-lower indicator. Let $h$ be a $2 \pi$-periodic, $\rho$-trigonometrically convex function ( $\rho$-t.c.f) and let $g$ be a $2 \pi$-periodic upper semicontinuous function. Further they are indicator and lower indicator of an entire function, and hence must satisfy the condition

$$
\begin{equation*}
h(\phi) \geq g(\phi), \phi \in[0,2 \pi) . \tag{5.4.3.1}
\end{equation*}
$$

An interval $(a, b) \subset[0,2 \pi)$ is called a maximal interval of $\rho$-trigonometricity of the function $h$ if

$$
\begin{equation*}
h(\phi)=A \cos \rho \phi+B \sin \rho \phi, \phi \in(a, b) \tag{5.4.3.2}
\end{equation*}
$$

for some constants $A, B$, and $h$ has no such representation on any larger interval $\left(a^{\prime}, b^{\prime}\right) \supset(a, b)$.

A function $h$ is said to be strictly $\rho$-t.c.f. if it is a $\rho$-t.c.f. and is not $\rho$ trigonometrical on any interval.

If the function $h$ is a strictly $\rho$-t.c.f., then $h$ and $g$ (satisfying other previous bounds) could be an indicator and lower indicator of an entire function $f \in A(\rho(r))$. However this is not so if the function $h$ has an interval of trigonometricity.

Recall, for example, the famous M. Cartwright Theorem [Le, Ch. IV, § 2, Thm. 6]: if an indicator of an entire function is trigonometrical on an interval $(a, b)$ with $b-a>\pi / \rho$, then the function is a CRG -function on this interval, i.e.,

$$
\begin{equation*}
h(\phi)=g(\phi), \phi \in(a, b) . \tag{5.4.3.3}
\end{equation*}
$$

Let us formulate all the necessary conditions of such kind. Let $(a, b)$ be a maximal interval of $\rho$-trigonometricity of the function $h$. The M. Cartwright theorem can be formulated as the implication

$$
\begin{equation*}
(b-a>\pi / \rho) \Rightarrow(5.4 .3 .3) \tag{5.4.3.4}
\end{equation*}
$$

The following implications are also necessary:

$$
\begin{align*}
\left(\exists \phi_{0} \in(a, b): h\left(\phi_{0}\right)=g\left(\phi_{0}\right)\right) & \Rightarrow(5.4 .3 .3),  \tag{5.4.3.5}\\
\left(h(a)=g(a) \wedge h_{+}^{\prime}(a)=h_{-}^{\prime}(a)\right) & \Rightarrow(5.4 .3 .3),  \tag{5.4.3.6a}\\
\left(h(b)=g(b) \wedge h_{+}^{\prime}(b)=h_{-}^{\prime}(b)\right) & \Rightarrow(5.4 .3 .3), \tag{5.4.3.6b}
\end{align*}
$$

where $h_{ \pm}^{\prime}(a)$ and $h_{ \pm}^{\prime}(b)$ are the right and left derivatives of the function $h$ at the points $a$ and $b$.

$$
\begin{align*}
\left(b-a=\pi / \rho \wedge h_{+}^{\prime}(a)=h_{-}^{\prime}(a)\right) & \Rightarrow(5.4 .3 .3),  \tag{5.4.3.7a}\\
\left(b-a=\pi / \rho \wedge h_{+}^{\prime}(b)=h_{-}^{\prime}(b)\right. & \Rightarrow(5.4 .3 .3),  \tag{5.4.3.7b}\\
\left(\liminf _{\phi \rightarrow a+0} \frac{h(\phi)-g(\phi)}{\phi-a}=0\right) & \Rightarrow(5.4 .3 .3),  \tag{5.4.3.8a}\\
\left(\liminf _{\phi \rightarrow b-0} \frac{h(\phi)-g(\phi)}{b-\phi}=0\right) & \Rightarrow \text { (5.4.3.3). } \tag{5.4.3.8b}
\end{align*}
$$

Now we shall give an exact formulation. The functions $h$ and $g$ are said to be concordant if at least one of the following conditions holds:

1. $h$ is strictly $\rho$-t.c.;
2. for each $(a, b)$ that is a maximal interval of $\rho$-trigonometricity of the function $h$ the implications (5.4.3.4)-(5.4.3.8b) are satisfied.

Theorem 5.4.3.1 Let $0<\rho<\infty, h(\phi)$ be a $2 \pi$-periodic, $\rho$-t.c.f., $g(\phi)$ be an upper semicontinuous, $2 \pi$-periodic function, $h(\phi) \geq g(\phi)$ for all $\phi$, and $h \not \equiv g$.

A function $f \in A(\rho(r))$ which simultaneously satisfies the identity $h_{f} \equiv$ $h, \underline{h}_{f} \equiv g$ with an arbitrary proximate order $\rho(r) \rightarrow \rho$ exists if and only if the functions $h$ and $g$ are concordant.

### 5.4.4

Proof of necessity. Note that implication (5.4.3.4) is a corollary of (5.4.3.6a) or (5.4.3.6b), because every $\rho$-trigonometrical function is continuous and has continuous derivative in $(a, b)$. Recall that $(\bullet)_{[t]}$ was defined by (3.1.2.4a).

From properties of the limit set $\mathbf{F r}[f]$ (Theorem 3.1.2.2, fr2), fr3)) and the definition of indicators $((3.1 .2 .1),(3.1 .2 .2))$ we can obtain for every function $v \in$ $\operatorname{Fr}[f]$ the inequality

$$
\begin{equation*}
v\left(\tau e^{i \phi}\right) \leq \tau^{\rho} h(\phi), \phi \in[0,2 \pi), \tau>0 \tag{5.4.4.1}
\end{equation*}
$$

Since $h(\phi)$ is $\rho$-trigonometrical for $\phi \in(a, b)$, the function

$$
H\left(r e^{i \phi}\right):=r^{\rho} h(\phi)
$$

is harmonic in the angle

$$
\Gamma(a, b):=\left\{r e^{i \phi}: \phi \in(a, b), r \in(0, \infty)\right\}
$$

whence the function $v-H$ is subharmonic and nonpositive in $\Gamma(a, b)$. By virtue of the maximum principle, either $v<H$ in $\Gamma(a, b)$ or $v \equiv H$ in $\Gamma(a, b)$ for each $v \in \operatorname{Fr}[f]$. Note that the condition $v \equiv H$ in $\Gamma(a, b)$ implies $v \equiv H$ in $\Gamma[a, b]$ for the closed interval because of the upper semicontinuity of $v$.

Let us prove (5.4.3.5). For every $v \in \mathbf{F r}[f]$ we have $v\left(r e^{i \phi_{0}}\right)-H\left(r e^{i \phi_{0}}\right)=0$ whence by the maximum principle $v=H$ in $\Gamma(a, b)$. Hence (5.4.3.3) holds.

Let us prove (5.4.3.6a). Assume the contrary: $h(a)=g(a) \wedge h_{+}^{\prime}(a)=h_{-}^{\prime}(a)$ holds, but there exists $\phi_{0} \in(a, b)$ such that $h\left(\phi_{0}\right)>g\left(\phi_{0}\right)$. Then there exists $v \in \operatorname{Fr}[f]$ such that

$$
g\left(\phi_{0}\right) \leq v\left(e^{i \phi_{0}}\right)<h\left(\phi_{0}\right)
$$

whence

$$
\begin{equation*}
v\left(\tau e^{i \phi}\right)<\tau^{\rho} h(\phi) \in \Gamma(a, b) \tag{5.4.4.2}
\end{equation*}
$$

Without loss of generality, we can assume that $v(z)>-\infty$, otherwise we can replace $v$ with $\max (v,-C)$ for a large positive constant $C>0$.

We choose $0<\tau_{1}<\tau_{2}$ and to every function

$$
W_{j}\left(r e^{i \phi}\right):=v_{\left[\tau_{j}\right]}\left(r e^{i \phi+a}\right)-r^{\rho} h(\phi+a), j=1,2, \gamma=b-a, r e^{i \phi} \in \Gamma(0, \gamma)
$$

we apply the following lemma due to A.E. Eremenko and M.L. Sodin [So] (see also [PW, Ho]):

Lemma 5.4.4.1 (E.S.) Let $W$ be a subharmonic nonpositive function inside the angle $\Gamma(0, \gamma), \gamma>0$. Then the following implication is valid,

$$
\left(\limsup _{\phi \rightarrow 0} \frac{W\left(e^{i \phi}\right)}{\phi}=0\right) \Rightarrow W \equiv 0 .
$$

If the condition of this theorem is not satisfied for

$$
W^{*}\left(r e^{i \phi}\right)=\max _{\tau \in\left[\tau_{1}, \tau_{2}\right]} v_{[\tau]}\left(r e^{i \phi}\right)
$$

it would be possible to insert a $\rho$-t.c.function between $h(\phi)-\epsilon(\phi-a)$ (for a small $\epsilon)$ and $v\left(e^{i \phi}\right)$. However, such a function does not exist, because of the negative jump of the derivative. So it will be a contradiction. See further for details.

From Lemma 5.4.4.1 we get

$$
\liminf _{\phi \rightarrow a+0} \frac{h(\phi)-v_{\left[\tau_{1}\right]}\left(e^{i \phi}\right)}{\phi-a}:=\alpha_{1}>0
$$

and likewise

$$
\liminf _{\phi \rightarrow a+0} \frac{h(\phi)-v_{\left[\tau_{2}\right]}\left(e^{i \phi}\right)}{\phi-a}:=\alpha_{2}>0 .
$$

So a $\Delta>0$ can be chosen such that $a+\Delta<b$ and the inequalities

$$
\begin{equation*}
H\left(\tau_{j} e^{i \phi}\right)-v_{\left[\tau_{j}\right]}\left(e^{i \phi}\right)>\alpha \tau_{j}^{\rho}(\phi-a), j=1,2, \tag{5.4.4.3}
\end{equation*}
$$

where $\alpha:=1 / 2 \min \left(\alpha_{1}, \alpha_{2}\right)$, hold for all $\phi \in[a, a+\Delta]$.
We write

$$
\beta:=\min _{\tau \in\left[\tau_{1}, \tau_{2}\right]}\left(H\left(\tau e^{i(a+\Delta)}\right)-v\left(\tau e^{i(a+\Delta)}\right)\right)
$$

which is positive because of (5.4.4.2).
Let us choose $\epsilon>0$ small enough to

$$
\begin{equation*}
\epsilon<\min \left(\alpha, \beta\left(\tau_{2}\right)^{-\rho} \Delta^{-1}\right) \tag{5.4.4.4}
\end{equation*}
$$

and let us consider the $\rho$-trigonometrical function

$$
h_{\epsilon}(\phi):=\rho^{-1}\left(h^{\prime}(a)-\epsilon\right) \sin \rho(\phi-a)+h(a) \cos \rho(\phi-a), \phi \in(a, b)
$$

that coincides with

$$
h(\phi)=\rho^{-1} h^{\prime}(a) \sin \rho(\phi-a)+h(a) \cos \rho(\phi-a), \phi \in(a, b)
$$

in the point $\phi=a$ but has a tangent that is lower than the tangent of $h$.
Further

$$
\begin{equation*}
h(\phi)-h_{\epsilon}(\phi)=\rho^{-1} \epsilon \sin \rho(\phi-a) \leq \epsilon(\phi-a), \phi \in[a, a+\Delta] . \tag{5.4.4.5}
\end{equation*}
$$

Combining (5.4.4.3)-(5.4.4.5) we obtain

$$
\begin{align*}
v_{\left[\tau_{j}\right]}\left(e^{i \phi}\right) & <\tau_{j}^{\rho} h(\phi)-\alpha(\phi-a)<\tau_{j}{ }^{\rho} h(\phi)-\epsilon(\phi-a)  \tag{5.4.4.6}\\
& \leq \tau_{j}^{\rho} h_{\epsilon}(\phi), \phi \in[a, a+\Delta], j=1,2 \\
v\left(\tau e^{i \phi}\right) & \leq \tau^{\rho} h_{\epsilon}(a+\Delta)+\tau^{\rho} \epsilon \Delta-\beta  \tag{5.4.4.7}\\
& <\tau^{\rho} h_{\epsilon}(a+\Delta), \tau \in\left[\tau_{1}, \tau_{2}\right] .
\end{align*}
$$

We write

$$
\begin{equation*}
G:=\left\{r e^{i \phi}: \phi \in[a, a+\Delta], \tau \in\left[\tau_{1}, \tau_{2}\right]\right\} . \tag{5.4.4.8}
\end{equation*}
$$

It follows from (5.4.4.6), (5.4.4.7) that

$$
v\left(r e^{i \phi}\right)<r^{\rho} h_{\epsilon}(\phi), r e^{i \phi} \in \partial G
$$

where $\partial G$ is the boundary of the domain $G$. Since the functions $v\left(r e^{i \phi}\right)$ and $r^{\rho} h_{\epsilon}(\phi)$ are subharmonic in $G$, by virtue of the maximum principle we have

$$
\begin{equation*}
v\left(r e^{i \phi}\right)<r^{\rho} h_{\epsilon}(\phi), r e^{i \phi} \in G \tag{5.4.4.9}
\end{equation*}
$$

Let us consider the function

$$
H_{1}\left(r e^{i \phi}\right):=r^{\rho} h_{1}(\phi), r e^{i \phi} \in \Gamma(a-\Delta, a+\Delta)
$$

where

$$
h_{1}(\phi):= \begin{cases}h(\phi), & \phi \in(a-\Delta, a], \\ h_{\epsilon}(\phi), & \phi \in[a, a+\Delta) .\end{cases}
$$

The function $H_{1}$ is continuous in $\Gamma(a-\Delta, a+\Delta)$ and subharmonic in the angles $\Gamma(a-\Delta, a)$ and $\Gamma(a, a+\Delta)$. Let us prove that it is subharmonic at the point $z=e^{i a}$. Let $\mathcal{M}(z, R, v)$ be the mean value of $v$ over the circle $\{\zeta:|\zeta-z|=R\}$ (see (2.6.1.1)). Taking into consideration (5.4.4.9) and subharmonicity of $v$ (see (2.6.1.1)), for all small $R$ we have

$$
\mathcal{M}\left(e^{i a}, R, H_{1}\right) \geq M\left(e^{i a}, R, v\right) \geq v\left(e^{i a}\right)=H_{1}\left(e^{i a}\right)
$$

Hence $H_{1}$ is subharmonic for $z=e^{i a}$. Since $H_{1}$ is homogeneous, i.e., $H_{1}(k z)=$ $k^{\rho} H_{1}(z)$,

$$
\mathcal{M}\left(k e^{i a}, k R, H_{1}\right)=k^{\rho} \mathcal{M}\left(e^{i a}, R, H_{1}\right) \geq k^{\rho} H_{1}\left(e^{i a}\right)=H_{1}\left(k e^{i a}\right)
$$

So $H_{1}$ is subharmonic on the ray $\left\{z=k e^{i a}: k \in(0, \infty)\right\}$ and hence in the angle $\Gamma(a-\Delta, a+\Delta)$. Thus $h_{1}(\phi)$ is a $\rho$-t.c.f. for $\phi \in(a-\Delta, a+\Delta)$. However, by construction

$$
\left(h_{1}\right)_{-}^{\prime}(a)=h_{-}^{\prime}(a)=h_{+}^{\prime}(a)=\left(h_{\epsilon}\right)_{+}^{\prime}+\epsilon=\left(h_{1}\right)^{\prime}(a)+\epsilon
$$

and this contradicts the fact that $h_{1}$ is $\rho$-t.c.f.
Concordance of the implication (5.4.3.6a) is proved.
5.4.5 Here we continue the proof of necessity. Pass to the proof of necessity of the condition (5.4.3.7a). Assume the contrary. Then there exists $v \in \mathbf{F r}[f]$ and $\phi_{0} \in[a, b]$ such that $g\left(\phi_{0}\right) \leq v\left(e^{i \phi_{0}}\right)<h\left(\phi_{0}\right)$, whence by virtue of the maximum principle, $v\left(\tau e^{i \phi}\right)<\tau^{\rho} h(\phi)$ for $\tau e^{i \phi} \in \Gamma(a, b)$. Actually $v\left(\tau e^{i \phi}\right) \leq \tau^{\rho} h(\phi)$ everywhere and on the circle we have strict inequality. If $v\left(\tau e^{i a}\right)=H\left(\tau e^{i a}\right)$ for a $\tau>0$, then $v_{[\tau]}\left(e^{i a}\right)=h(a)$, and it will suffice to repeat the arguments used in proving (5.4.3.6a) with $v_{[\tau]}$ instead of $v$.

Exercise 5.4.5.1 Do that.
So it is sufficient to examine the case $v\left(\tau e^{i a}\right)<H\left(\tau e^{i a}\right), \tau>0$. Denote

$$
T(\phi):=h^{\prime}(a) \rho^{-1} \sin \rho(\phi-a)+h(a) \cos \rho(\phi-a) .
$$

This is a $\rho$-trigonometrical function, the graph of which is tangent to the graph of $h(\phi)$ at the point $a$.

There are two possibilities for $T(\phi)$ on some small interval $\phi \in(a-\gamma, a), \gamma>$ 0 : either $T(\phi)<h(\phi)$ or $T(\phi)=h(\phi)$.

Inequality $T(\phi)>h(\phi)$ contradicts $\rho$-t.convexity at the point $a$. The equality on the sequence of points $\phi_{j} \rightarrow a-0$ contradicts the maximum principle for $\rho$ t.c.functions.

Exercise 5.4.5.2 Why is it?
If $T(\phi)=h(\phi), \phi \in(a-\gamma, a)$, then $h$ is $\rho$-trigonometrical on the interval $(a-\gamma, b) \supset(a, b)$ that was already considered in the case (5.4.3.4) (M. Cartwright's Theorem).

So we assume $T(\phi)<h(\phi), \phi \in(a-\gamma, a)$. We set

$$
\begin{aligned}
h_{1}(\phi) & :=h(\phi)-T(\phi), \phi \in(a-\gamma, a), \\
v_{1}\left(r e^{i \phi}\right) & :=v\left(r e^{i \phi}\right)-r^{\rho} T(\phi), r e^{i \phi} \in \Gamma(a-\gamma, b) .
\end{aligned}
$$

Then $h_{1}(\phi)=0$ for $\phi \in[a, b], h_{1}(\phi)>0$ for $\phi \in(a-\gamma, a)$ and $h^{\prime}(a)=0$.
The function $v_{1}\left(e^{i \phi}\right)<0, \phi \in[a, b)$. Let us analyze the behavior of the function $v_{1}\left(e^{i \phi}\right)$ at the point $b$. Either $v_{1}\left(e^{i b}\right)<0$ or $v_{1}\left(e^{i b}\right)=0$ but

$$
\limsup _{\phi \rightarrow b-0} v_{1}\left(e^{i \phi}\right)(b-\phi)^{-1} \leq-C
$$

for some $C>0$ by Lemma 5.4.4.1 (E.S.).
From the other side $v_{1}\left(e^{i \phi}\right)$ is strictly negative also in some left (say, ( $a-$ $\Delta, a)$ ) neighborhood of $a$ because of upper semicontinuity. In any case $v_{1}\left(e^{i \phi}\right)$ can be majorated on the interval $(a-\Delta, b)$ by the function

$$
h_{\epsilon}:=-A \sin (\rho-\epsilon)(b-\phi)
$$

with sufficiently small $A$.
A point of intersection of the graph of $h_{\epsilon}$ with the axis $0, \phi$ can be regulated by $\epsilon$ and can be chosen so close to the point $a$ that the graph of $h_{\epsilon}$ also intersects the graph of $h_{1}(\phi)$, at some point $\theta_{0}<a$ because $h_{1}(a)=h_{1}^{\prime}(a)=0$.

Exercise 5.4.5.3 Make the precise proof with all the estimates.
Let the parameters $A, \epsilon, \theta_{0}$ be fixed as above. Denote

$$
S:=\left\{r e^{i \phi}: \phi \in\left(\theta_{0}, b\right), 0<r<1\right\} .
$$

Then $H_{\epsilon}\left(r e^{i \phi}\right):=r^{\rho-\epsilon} h_{\epsilon}(\phi)$ is harmonic in the sector $S$ and satisfies the inequality $H_{\epsilon}\left(r e^{i \phi}\right) \geq v_{1}\left(r e^{i \phi}\right)$ on $\partial S$. Hence $H_{\epsilon}\left(r e^{i \phi}\right) \geq v_{1}\left(r e^{i \phi}\right)$ on $S$. Thus

$$
\begin{equation*}
v\left(r e^{i \phi}\right) \leq H\left(r e^{i \phi}\right)+H_{\epsilon}\left(r e^{i \phi}\right), r e^{i \phi} \in S . \tag{5.4.5.1}
\end{equation*}
$$

Let $\mathcal{M}(r, v)$ be the mean value of the function on the circle $\{\zeta:|\zeta|=r\}$ (see 2.6.1.1). Using (5.4.5.1) we have

$$
\begin{aligned}
\mathcal{M}(r, v) & \leq \int_{\theta_{0}}^{b}\left[H\left(r e^{i \phi}\right)+H_{\epsilon}\left(r e^{i \phi}\right)\right] d \phi+\int_{[0,2 \pi) \backslash\left(\theta_{0}, b\right)} H\left(r e^{i \phi}\right) d \phi \\
& \leq d_{1} r^{\rho}-d_{2} r^{\rho-\epsilon}, \quad d_{1}, d_{2}>0 .
\end{aligned}
$$

So we get $\mathcal{M}(r, v)<0=v(0)$ for sufficiently small $r>0$ which contradicts the subharmonicity of the function $v$ at zero.
5.4.6 Now we complete proof of necessity, proving (5.4.3.8a,b). Assume the contrary: suppose

$$
\begin{equation*}
\liminf _{\phi \rightarrow a+0} \frac{h(\phi)-g(\phi)}{\phi-a}=0 \tag{5.4.6.1}
\end{equation*}
$$

but there exists a $\phi_{0} \in(a, b)$ such that $h\left(\phi_{0}\right)>g\left(\phi_{0}\right)$. Then there exists a function $v \in \mathbf{F r}[f]$ such that

$$
\begin{equation*}
g\left(\phi_{0}\right) \leq v\left(e^{i \phi_{0}}\right)<h\left(\phi_{0}\right) . \tag{5.4.6.2}
\end{equation*}
$$

Then the function $v_{1}\left(r e^{i \phi}\right):=v\left(r e^{i \phi}\right)-H\left(r e^{i \phi}\right)$ is subharmonic and nonpositive in $\Gamma(a, b)$. By virtue of the maximum principle $v_{1}\left(r e^{i \phi}\right)<0$, $r e^{i \phi} \in \Gamma(a, b)$.

From (5.4.6.1) we obtain

$$
0=\liminf _{\phi \rightarrow a+0} \frac{h(\phi)-g(\phi)}{\phi-a} \geq \liminf _{\phi \rightarrow a+0} \frac{h(\phi)-v\left(e^{i \phi}\right)}{\phi-a}=-\liminf _{\phi \rightarrow a+0} \frac{v_{1}\left(e^{i \phi}\right)}{\phi-a}
$$

whence, recollecting that $v_{1}\left(e^{i \phi}\right)<0$, we get

$$
\limsup _{\phi \rightarrow a+0} \frac{v_{1}\left(e^{i \phi}\right)}{\phi-a}=0
$$

Applying Lemma 5.4.4.1 (E.S.) to the function

$$
W\left(r e^{i \phi}\right)=v_{1}\left(r e^{i \phi+a}\right), r e^{i \phi} \in \Gamma(0, \gamma), \gamma=b-a
$$

we get $v_{1} \equiv 0$ in $\Gamma(a, b)$ which leads to a contradiction. The implication (5.4.3.8b) is proved in the same way. So the proof of necessity in Theorem 5.4.3.1 is completed.

We do not include here the proof of sufficiency and refer the readers to the original paper $[\mathrm{Po}(1992)]$.

### 5.5 Asymptotic extremal problems. Semiadditive integral

5.5.1 Suppose some class of entire functions is determined by asymptotic behavior of their zeros, and we want to know what is the restriction on asymptotic behavior of functions: for example, to estimate the indicator of such a function. The first example of such a problem was considered by B.Ya. Levin in [Le, Ch. IV, §1, Example]. A developed theory of such estimates was constructed in the papers of A.A. Gol'dberg [Go(1962)] and his pupils [Kon], [KF]. We consider this theory from the point of view of limit sets.

Let $\mathcal{M} \Subset \mathcal{M}(\rho)$ (see (3.1.3.4)) be a convex set of measures which is closed in $\mathcal{D}^{\prime}$ and is invariant with respect to the transformation $(\bullet)_{t}$ (see (3.1.3.1), (3.1.3.2)) and let $A(\mathcal{M})$ be a class of entire functions $f$ for which $\operatorname{Fr}\left[n_{f}\right] \subset \mathcal{M}$. We suppose $\rho$ is non-integer. Recall that canonical potential $\Pi(z, \nu, p)$ is defined by: (see (2.9.2.1))

$$
\Pi(z, \nu, p):=\int_{\mathbb{C}} G_{p}(z / \zeta) \nu(d \zeta)
$$

where $\nu$ is a measure and

$$
G_{p}(z):=\log |1-z|+\Re \sum_{k=1}^{p} \frac{z^{k}}{k} .
$$

Theorem 5.5.1.1 [AP] The relation

$$
\begin{equation*}
h(\phi, f)=\sup \left\{\Pi\left(e^{i \phi}, \nu, p\right): \nu \in \mathcal{M}\right\} \tag{5.5.1.1}
\end{equation*}
$$

is valid. There exists $f \in A(\mathcal{M})$ for which the equality holds in (5.5.1.1) for all $\phi$.
Proof. We should only prove that there exists an entire function with such indicator. Consider the set

$$
\Lambda:=\left\{\Pi\left(e^{i \phi}, \nu, p\right): \nu \in \mathcal{M}\right\}
$$

It is a convex set contained in $U[\rho]$. Thus there exists a subharmonic (see Corollary 4.1.4.2) and hence entire (see Corollary 5.3.1.5) function $f$ such that $\operatorname{Fr}[f]=\Lambda$. By Theorem 5.4.2.1, (5.5.1.1) holds.

For some $\mathcal{M}$ it is possible to compute the supremum in (5.5.1.1) and thus to obtain explicit precise estimates of indicators in the respective class $A(\mathcal{M})$. As an example, we shall present an estimate given by A.A. Gol'dberg.

We recall that the upper density of zeros of an entire function $f \in A(\rho)$ is defined by the equality

$$
\bar{\Delta}\left[n_{f}\right]:=\limsup _{r \rightarrow \infty} \frac{n_{f}(r)}{r^{\rho}}
$$

where $n_{f}$ is the distribution of zeros of the function $f$, and denote

$$
\begin{equation*}
K(t, \phi):=-\left[\frac{d}{d t} G_{p}^{+}\left(e^{i \phi} / t\right)\right]^{-} \tag{5.5.1.2}
\end{equation*}
$$

where $a^{+}:=\max (a, 0), a^{-}:=\min (a, 0)$. This function is piecewise continuous.
Corollary 5.5.1.2 [Go(1962)] Let the distribution of zeros $n_{f}$ of a function $f$ be concentrated on the positive ray, and let $\bar{\Delta}\left[n_{f}\right] \leq \Delta<\infty$. Then

$$
\begin{equation*}
h(\phi, f) \leq \Delta \int_{0}^{\infty} t^{\rho} K(t, \phi) d t, \phi \in[0,2 \pi) \tag{5.5.1.3}
\end{equation*}
$$

and there exists a function from the same class for which equality is attained for all $\phi$.

Proof of Corollary 5.5.1.2. We exploit Theorem 5.5.1.1. The class of functions $f$ satisfying the assumption of the corollary coincides with the class of $f$ for which

$$
\begin{equation*}
\operatorname{Fr}\left[n_{f}\right] \subset \mathcal{M}=\left\{\nu \in \mathcal{M}(\rho): \operatorname{supp} \nu \subset[0, \infty] \wedge \nu(r) \leq \Delta r^{\rho}\right\} . \tag{5.5.1.4}
\end{equation*}
$$

Exercise 5.5.1.1 Show this by using Corollary 3.3.2.6.
Thus

$$
\Pi\left(e^{i \phi}, \nu, p\right)=\int_{0}^{\infty} G_{p}\left(e^{i \phi} / t\right) \nu(d t) \leq \int_{0}^{\infty} G_{p}^{+}\left(e^{i \phi} / t\right) \nu(d t)
$$

Integrating by parts we obtain

$$
\Pi\left(e^{i \phi}, \nu, p\right) \leq-\int_{0}^{\infty} \nu(t)\left[\frac{d}{d t} G_{p}^{+}\left(e^{i \phi} / t\right)\right]^{-} d t
$$

By (5.5.1.4) we get (5.5.1.3).
We write

$$
M_{p}(r):=\max \left\{G_{p}\left(r e^{i \phi}\right): \phi \in[0,2 \pi)\right\}
$$

In the same way one can prove
Corollary 5.5.1.3 [Go(1962), Thm. 4.1] Let distribution of zeros of the function $f \in A(\rho)$ satisfy only the condition $\bar{\Delta}\left[n_{f}\right] \leq \Delta<\infty$. Then

$$
\begin{equation*}
h(\phi, f) \leq \Delta \rho \int_{0}^{\infty} t^{\rho-1} M_{p}(1 / t) d t, \phi \in[0,2 \pi) \tag{5.5.1.5}
\end{equation*}
$$

and there exists a function from the same class for which equality is attained for all $\phi$.

Exercise 5.5.1.2 Prove this corollary exploiting

$$
\mathcal{M}:=\left\{\nu \in \mathcal{M}(\rho): \nu(r) \leq \Delta r^{\rho}, \forall r>0\right\} .
$$

5.5.2 To be able to obtain explicit estimates for more diverse classes of entire functions defined by a restriction on the density of zeros, Gol'dberg introduced an integral with respect to a nonadditive measure and obtained estimates for indicators in terms of a one-dimensional integral (along a circumference) with respect to such a measure ([Go(1962)]. Gol'dberg initially constructed the integral sum of a special form. The construction presented here is based on the Levin-MatsaevOstrovskii theorem (see [Go(1962), Thm. 2.10]). Fainberg (1983) developed this approach using a two-dimensional integral. This made it possible to extend significantly the set of classes of entire functions for which the estimate expressed by a nonadditive integral is precise. We shall present these results after the necessary definitions.

Let $\delta(X)$ be a nonnegative monotonic function of $X \subset \mathbb{C}$, the function being finite on bounded sets and $\delta(\varnothing)=0$. For a given family of sets $\mathcal{X}:=\{X\}$ we denote by $N(\delta, \mathcal{X})$ the class of countable-additive measures $\mu$ defined by the relation

$$
N(\delta, \mathcal{X}):=\{\mu: \mu(X) \leq \delta(X), X \in \mathcal{X}\}
$$

For a Borel function $f \geq 0$ we define the quantity

$$
(\mathcal{X}) \int f d \delta:=\sup \left\{\int f d \mu: \mu \in N(\delta, \mathcal{X})\right\}
$$

called an $(\mathcal{X})$-integral with respect to a nonnegative measure $\delta$. For a Borel set $E \subset \mathbb{C}$ we set

$$
(\mathcal{X}) \int_{E} f d \delta:=(\mathcal{X}) \int f I_{E} d \delta
$$

where $I_{E}$ is an indicator of the set $E$, i.e.,

$$
I_{E}(z):= \begin{cases}1, & \text { if } z \in E \\ 0 & \text { if } z \notin E\end{cases}
$$

This integral possesses a number of natural properties: it is monotonic with respect to $f$ and $\delta$ and the family $\mathcal{X}$, positively homogeneous and semi-additive with respect to the function $f$ and $\delta$. If $\delta$ is a measure, if $\mathcal{X}$ is a Borel ring, and if $f$ is a measurable function, then $(\mathcal{X})$-integral coincides with the Lebesgue-Stieltjes integral.

Exercise 5.5.2.1 Check these properties.

Let $\delta(\Theta)$ be a nonadditive measure on the unit circle $\mathbb{T}$, defined initially on the family of all open sets $\Theta \subset \mathbb{T}$. It can be naturally extended to all closed sets $\Theta^{F}$ using the equality

$$
\delta\left(\Theta^{F}\right):=\inf \left\{\delta(\Theta): \Theta \supset \Theta^{F}\right\}
$$

Let $\chi_{\Theta}$ be a set of open sets containing the set $\mathbb{T}$. We write

$$
D_{r, \Theta}:=\left\{z=t e^{i \theta}: 0<t<r, e^{i \theta} \in \Theta\right\}, \chi_{z}:=\left\{D_{r, \Theta}: r>0, \Theta \in \chi_{\Theta}\right\}
$$

The subscripts $\Theta$ and $z$ at $\chi$ indicate that the families under consideration are located either on $\mathbb{T}$ or on the plane, respectively.

Let us define a nonadditive measure $\delta_{z}$ on $\chi_{z}$ by the equalities

$$
\delta_{z}\left(D_{r, \Theta}\right):=r^{\rho} \delta(\Theta), D_{r, \Theta} \in \chi_{z}
$$

Now the integral $\left(\chi_{z}\right) \int G_{p}^{+}\left(e^{i \theta} / \zeta\right) d \delta_{z}$ is defined.
Recall that the classical angular upper density of zeros of an entire function $f \in A(\rho)$ is defined by the equality (compare (3.3.2.7))

$$
\bar{\Delta}^{\mathrm{cl}}\left[n_{f}, \Theta\right]:=\limsup _{r \rightarrow \infty} n_{f}\left(D_{r, \Theta}\right) r^{-\rho}
$$

Consider the class of entire functions $A^{\mathrm{cl}}\left(\delta, \chi_{\Theta}\right)$ defined by the equality

$$
\begin{equation*}
A^{\mathrm{cl}}\left(\delta, \chi_{\Theta}\right):=\left\{f: \bar{\Delta}^{\mathrm{cl}}\left[n_{f}, \Theta\right] \leq \delta(\theta), \forall \Theta \in \chi_{\Theta}\right\} \tag{5.5.2.1}
\end{equation*}
$$

for a given non-additive measure $\delta(\Theta)$ and a family $\chi_{\Theta}$.
Theorem 5.5.2.1 [Fa] Let $\delta(\Theta)$ satisfy the condition

$$
\begin{equation*}
\delta(\Theta)=\delta(\bar{\Theta}), \forall \Theta \in \chi_{\Theta} \tag{5.5.2.2}
\end{equation*}
$$

(the dash means the closure of a set). Then

$$
\begin{equation*}
h(\phi, f) \leq\left(\chi_{z}\right) \int G_{p}^{+}\left(e^{i \theta} / \zeta\right) d \delta_{z} \tag{5.5.2.3}
\end{equation*}
$$

There exists a function $f \in A^{\mathrm{cl}}\left(\delta, \chi_{\Theta}\right)$ such that equality in (5.5.2.7) is attained for all $\phi \in[0,2 \pi)$ simultaneously.

Proof. Let us note the following: If we replace in this theorem $\bar{\Delta}^{\mathrm{cl}}\left[n_{f}, \Theta\right]$ with its $\mathcal{D}^{\prime}$ counterpart $\bar{\Delta}\left(C o_{\Theta}\left(I_{1}\right)\right)$ (see Theorem 3.3.1.2) and consider the corresponding class of entire functions $A\left(\delta, \chi_{\Theta}\right)$, the assertion of the theorem holds without conditions (5.1.5.6). You should only apply Theorem 5.5.1.1 with the corresponding $\mathcal{M}$. The condition (5.5.2.7) is exploited only for replacing " $\mathcal{D}$ '" quantities by the classic ones using results of Sections 3.3.2.

Exercise 5.5.2.2 Prove this theorem in detail.
It is also worth noting that every family $\chi_{\Theta}$ can be replaced by a family $\chi_{\Theta}^{\prime}$ that is dense in $\chi_{\Theta}$ (see 3.2.2) and such that for $\chi_{\Theta}^{\prime}$ (5.5.2.6) already holds (see Theorem 3.3.2.3).

### 5.6 Entire functions of completely regular growth. Levin-Pfluger Theorem. Balashov's theory

5.6.1 The most famous definition of a function of completely regular growth (CRGfunction) is the following:

A function $f \in A(\rho(r))$ is a function of completely regular growth, if the limit

$$
\lim _{z \rightarrow \infty} r^{-\rho(r)} \log |f(z)|, r:=|z|
$$

exists when $z \rightarrow \infty$ uniformly outside some $C_{0}^{1}$-set (see Section 5.2.1.)
Actually, it is equivalent to all other definitions of the functions of completely regular growth in the plane (compare [Le, Ch. III], $[\operatorname{Pf}(1938)],[\operatorname{Pf}(1939)])$.

By A.A. Gol'dberg ([Go(1967)]) this definition was reduced to the following:
A function $f \in A(\rho(r))$ is a function of completely regular growth, if

$$
\underline{h}_{f}(\phi)=h_{f}(\phi), \forall \phi \in[0,2 \pi) .
$$

Because of the formulae (3.2.1.1), (3.2.1.2) (see also Section 3.2.7) we have the following

Theorem 5.6.1.1 A function $f \in A(\rho(r))$ is a function of completely regular growth (CRG-function) iff $\mathbf{F r}[f]$ consists of only one subharmonic function $h(z)$.

Because of (3.2.1.11) the function $h(z)$ has the form

$$
\begin{equation*}
h(z)=r^{\rho} h\left(e^{i \phi}\right) . \tag{5.6.1.1}
\end{equation*}
$$

The function $h(\phi):=h\left(e^{i \phi}\right)$ is $\rho$-trigonometrically convex and it was studied in Sections 3.2.3, 3.2.4, 3.2.5.
5.6.2 The initial definition of regular zero distribution [Le, Ch. II, § 1] is the following:

Let $n$ be a zero distribution (divisor, or mass distribution) of convergence exponent $\rho_{1}:=\rho[n]$ (see Section 2.8.3), and let $\rho_{1}>\left[\rho_{1}\right]$. Let $\rho_{1}(r) \rightarrow \rho_{1}$ be a proper proximate order of $n(r)$ (see Theorem 2.8.1.2). It means that $n \in \mathcal{M}(\rho(r)), \rho(r) \rightarrow$ $\rho_{1}$ (see Section 3.1.3).

The initial definition of regular zero distribution for $\rho_{1}$ being non-integer is:
A zero distribution $n$ is regular if the limit

$$
\lim _{r \rightarrow \infty} \frac{n\left(\operatorname{Co}_{(\alpha, \beta)}\left(I_{t}\right)\right)}{t^{\rho_{1}(t)}}:=\Delta((\alpha, \beta))
$$

exists for all $\alpha>\beta$ except perhaps for a countable set on the circle.

[^5]By using results of Section 3.3, one can prove
Theorem 5.6.2.1 The zero distribution $n$ is regular iff $\mathbf{F r}[n]$ consists of only one measure $\nu_{\mathrm{reg}}$.

Exercise 5.6.2.1 Prove this exploiting Theorems 3.3.3.1 and 3.3.2.4.
Recall that for $f \in A(\rho(r)), \rho(r) \rightarrow \rho$ we have $n_{f} \in \mathcal{M}(\rho(r)), \rho(r) \rightarrow \rho$ (see Theorem 2.9.3.2). Now we can formulate

Theorem 5.6.2.2 (Levin-Pfluger) [Le, Ch. II, Ch. III] An entire function $f \in$ $A(\rho(r)), \rho(r) \rightarrow \rho$ of non-integer order $\rho$ is of completely regular growth function iff its zero distribution is regular.

After Theorems 5.6.1.1, 5.6.2.1 this theorem is a direct corollary of Theorem 3.1.5.1.
5.6.3 Consider now the case of integer $\rho$. In general, this case differs from the case of non-integer $\rho$. For example, Theorem 2.9.4.2 (Brelot-Lindelöf) implies that

$$
(f \in A(\rho(r)), \rho(r) \rightarrow \rho) \Longleftrightarrow n_{f} \in \mathcal{M}(\rho(r)), \rho(r) \rightarrow \rho
$$

iff the family of polynomials (2.9.4.4a) is compact.
To describe the regularity of zero distribution for the case of integer $\rho$ we assume that the limit

$$
\begin{equation*}
\lim _{R \rightarrow \infty} \delta_{R}(z, \nu, \rho):=\Re\left[\delta_{\infty} z^{\rho}\right] \tag{5.6.3.1}
\end{equation*}
$$

exists, where

$$
\delta_{\infty}:=\lim _{R \rightarrow \infty}\left[\int_{\lfloor\zeta \mid<R}|\zeta|^{-\rho} \cos \arg \zeta n(d \zeta)+i \int_{|\zeta|<R}|\zeta|^{-\rho} \sin \arg \zeta n(d \zeta)\right] .
$$

Now a zero distribution $n \in \mathcal{M}(\rho(r)), \rho(r) \rightarrow \rho$ with integer $\rho$ is called regular if $\mathbf{F r}[n]$ consists of only one measure $\nu_{\text {reg }}$ as in Theorem 5.6.2.1 and the limit (5.6.3.1) exists.
Under this definition Theorem 5.6.2.2 still holds, because the set $(\mathcal{H}, \mathbf{F r})[\log |f|]$ from Theorem 3.1.5.2 consists of only one element ( $\left.\Re\left[\delta_{\infty} z^{\rho}\right], \nu_{\text {reg }}\right)$.

Note also
Proposition 5.6.3.1 The measure $\nu_{\text {reg }}$ has the form

$$
\nu_{\mathrm{reg}}(d r d \phi)=\rho r^{\rho} d r \otimes \Delta(d \phi)
$$

where $\Delta$ is a measure of bounded variation on the unit circle.
This assertion is a corollary of invariance of $\mathbf{F r}[n]$, Theorem 3.1.3.3, frm3).
5.6.4 In the papers $[\operatorname{Bal}(1973), \operatorname{Bal}(1976)]$ functions of completely regular growth along curves of regular rotation were considered. A curve of regular rotation is a curve that is described by the equation

$$
z=t e^{i(\gamma(t) \log t+\phi)}, 0<t<\infty
$$

for a fixed $\phi$.
If $\gamma(t) \equiv \gamma$, then this curve is a logarithmic spiral. In the general case $\gamma(t)$ is a differentiable function such that

$$
\gamma(t) \rightarrow \gamma, t \gamma^{\prime}(t) \rightarrow 0, t \rightarrow \infty
$$

To describe this theory in terms of limit sets we consider the transformation

$$
\begin{aligned}
P_{t} & =t e^{i \gamma(t) \log t}, \\
u_{t}(z) & =u\left(P_{t} z\right) t^{-\rho(t)} .
\end{aligned}
$$

The following theorem is similar to Theorem 3.1.2.1.
Proposition 5.6.4.1 (Existence of spiral Limit Set) The following holds:
esls 1) $u_{t} \in S H(\rho(r))$ for all $t \in(0, \infty)$;
esls 2) the family $\left\{u_{t}\right\}$ is precompact at infinity.
The set of all limits $\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} u_{t_{j}}$ does not depend on $\gamma(t)$ but only on the constant $\gamma$ since

$$
\lim _{t \rightarrow \infty}(\gamma(t)-\gamma) \log t=\lim _{t \rightarrow \infty} t \gamma^{\prime}(t)=0
$$

So it is the same as that for

$$
P_{t}=t e^{i \gamma \log t}
$$

i.e., the case that was already considered in the general theory.

In particular (3.2.1.8) for this case has the form

$$
\begin{equation*}
z^{0}(z)=e^{i(-\gamma \log r+\phi)} \tag{5.6.4.1}
\end{equation*}
$$

Hence, from Theorem 3.2.1.2 the indicator (see (3.2.1.1)) has the form

$$
h\left(r e^{i \phi}\right)=r^{\rho} h(-\gamma \log r+\phi), z=r e^{i \phi},
$$

where $h(\phi)$ is a $\rho$-trigonometrically convex $2 \pi$-periodic function (see Section 3.2.3).
All other assertions of Levin-Pfluger theory can be obtained analogously from other general assertions as it was done in the previous sections.

Theorem 3.1.6.1 connects limit sets for every $\gamma$.
Exercise 5.6.4.1 Formulate and prove Balashov's analogy of the Levin-Pluger Theorem 5.6.2.2 and Theorem 3.1.6.1 for $m=2$.

For other generalizations of the Levin-Pfluger theory see $[\mathrm{AD}]$ and $[\mathrm{Az}(2007)]$.

### 5.7 General characteristics of growth of entire functions

5.7.1 A functional $\mathcal{F}(u)$ acting in the unit circle and defined on subharmonic functions $u \in S H(\rho(r))$ is called a growth characteristic if the following conditions are fulfilled:

1. continuity:

$$
\begin{equation*}
\mathcal{F}\left(u_{j}\right) \rightarrow \mathcal{F}(u) \tag{5.7.1.1}
\end{equation*}
$$

if $u_{j} \rightarrow u$ uniformly on compacts (of course, for continuous functions $u$ ) or if $u_{j} \downarrow u$;
2. positive homogeneity:

$$
\begin{equation*}
\mathcal{F}(c u)=c \mathcal{F}(r, u) \tag{5.7.1.2}
\end{equation*}
$$

for every constant $c>0$.
Here we shall list some widely used functionals that satisfy these conditions:

$$
\begin{align*}
H_{\phi}(u) & :=u\left(e^{i \phi}\right)  \tag{5.7.1.3}\\
T(u) & =\frac{1}{2 \pi} \int_{0}^{2 \pi} u^{+}\left(e^{i \phi}\right) d \phi  \tag{5.7.1.4}\\
M_{\alpha}(u) & :=\max \left\{u\left(e^{i \phi}\right):|\phi| \leq \alpha\right\}  \tag{5.7.1.5}\\
M(u) & :=M_{\pi}(u)  \tag{5.7.1.6}\\
I_{\alpha \beta}(u) & :=\int_{\alpha}^{\beta} u\left(e^{i \phi}\right) d \phi  \tag{5.7.1.7}\\
I(u, g) & :=\int_{0}^{2 \pi} u\left(e^{i \phi}\right) g(\phi) d \phi, g \in L^{1}[0,2 \pi] . \tag{5.7.1.8}
\end{align*}
$$

Exercise 5.7.1.1 Check properties 1 and 2 for these functionals.
Let $\alpha(t)$ and $\alpha_{\epsilon}(\zeta)$ be the "hats" defined by the equalities (2.3.1.1)-(2.3.1.3) and let $R_{\epsilon} u$ be defined by (2.3.1.4).

This averaging has the following properties.

## Proposition 5.7.1.1

1. if $u$ is subharmonic, then $R_{\epsilon} u$ is subharmonic;
2. $R_{\epsilon} u \downarrow u$ as $\epsilon \downarrow 0$ for every subharmonic function;
3. if $u_{j} \rightarrow u$ in $\mathcal{D}^{\prime}$ and $u_{j}, u$ are locally summable functions, $R_{\epsilon} u_{j} \rightarrow R_{\epsilon} u$ uniformly on every compact set.

Exercise 5.7.1.2 Prove this using Theorem 2.3.4.5, 2.6.2.3.
Now we can define the asymptotic characteristics of growth of entire function $f \in A(\rho(r)):$

$$
\begin{align*}
& \overline{\mathcal{F}}[f]:=\lim _{\epsilon \rightarrow 0} \limsup _{t \rightarrow \infty} \mathcal{F}\left(R_{\epsilon} u_{t}(\bullet)\right),  \tag{5.7.1.9}\\
& \underline{\mathcal{F}}[f]:=\lim _{\epsilon \rightarrow 0} \liminf _{t \rightarrow \infty} \mathcal{F}\left(R_{\epsilon} u_{t}(\bullet)\right), \tag{5.7.1.10}
\end{align*}
$$

where $u=\log |f|$ and $(\bullet)_{t}$ is defined by (3.1.2.1).
Proposition 5.7.1.2 For $\mathcal{F}(u)$ defined by (5.7.1.3)

$$
\overline{\mathcal{F}}[f]=h_{f}(\phi) ; \quad \underline{\mathcal{F}}[f]=\underline{h}_{f}(\phi) .
$$

For other functionals from the list (5.7.1.4)-(5.7.1.8) one may replace $R_{\epsilon} u$ by $u$ and omit $\lim _{\epsilon \rightarrow 0}$.

Exercise 5.7.1.3 Prove this.
The following assertion connects the asymptotic growth characteristics with limit sets.

Theorem 5.7.1.3 The relations

$$
\begin{aligned}
& \overline{\mathcal{F}}[f]=\sup \{\mathcal{F}(v): v \in \mathbf{F r}[f]\}, \\
& \underline{\mathcal{F}}[f]=\inf \{\mathcal{F}(v): v \in \mathbf{F r}[f]\}
\end{aligned}
$$

are true.
Proof. Let $v \in \operatorname{Fr}[f]$ and $u_{t_{j}} \rightarrow v$ in $\mathcal{D}^{\prime}$. Then $R_{\epsilon} u_{t_{j}} \rightarrow R_{\epsilon} v$ uniformly on every compact set. Hence

$$
\lim _{t_{j} \rightarrow \infty} \mathcal{F}\left(R_{\epsilon} u_{t_{j}}\right)=\mathcal{F}\left(R_{\epsilon} v\right)
$$

Passing to the limit as $\epsilon \rightarrow 0$ we obtain

$$
\lim _{\epsilon \rightarrow 0} \lim _{t_{j} \rightarrow \infty} \mathcal{F}\left(R_{\epsilon} u_{t_{j}}\right)=\mathcal{F}(v) .
$$

Choosing a sequence that corresponds to lim sup or liminf we obtain the assertion of the theorem.

Applying this theorem to the functional (5.7.1.3) we obtain the RHS's of (3.2.1.1), (3.2.1.2) and hence another definition for the indicator and lower indicator.
5.7.2 A family of growth characteristics $\chi_{A}:=\left\{\mathcal{F}_{\alpha}(r, \bullet): \alpha \in A\right\}$ is called total if the equation

$$
\begin{equation*}
\mathcal{F}_{\alpha}\left(v_{1}\right)=\mathcal{F}_{\alpha}\left(v_{2}\right), \forall r>0, \alpha \in A \tag{5.7.2.1}
\end{equation*}
$$

implies $v_{1} \equiv v_{2}$ for $v_{1}, v_{2} \in U[\rho]$ (see 3.1.2.4).

Here are some examples of the total families:

$$
\begin{align*}
\chi_{H} & :=\left\{H_{\phi}\left(u\left(e^{i \phi}\right): \phi \in[0,2 \pi)\right\}\right.  \tag{5.7.2.2}\\
\chi_{I} & :=\left\{I_{\alpha, \beta}(u): \alpha, \beta \in[0,2 \pi)\right\} ;  \tag{5.7.2.3}\\
\chi_{F o} & :=\left\{c_{k}(u)=I\left(u, g_{k}\right): k \in \mathbb{Z}\right\} ; \tag{5.7.2.4}
\end{align*}
$$

where

$$
\begin{equation*}
g_{0}:=1, g_{k}:=\cos k \phi ; g_{-k}=\sin k \phi, k \in \mathbb{N} . \tag{5.7.2.5}
\end{equation*}
$$

It is easy to deduce from Theorem 5.6.1.1
Theorem 5.7.2.1 Let a family $\left\{\mathcal{F}_{\alpha}(\bullet): \alpha \in A\right\}$ be a total family of characteristics. An entire function $f$ is a CRG-function iff

$$
\begin{equation*}
\overline{\mathcal{F}}_{\alpha}[f]=\underline{\mathcal{F}}_{\alpha}[f] . \tag{5.7.2.6}
\end{equation*}
$$

Exercise 5.7.2.1 Check this.
5.7.3 Let us consider a total family of characteristics of the form

$$
\begin{equation*}
\chi_{\Psi}:=\{I(u, \psi): \psi \in \Psi\} \tag{5.7.3.1}
\end{equation*}
$$

where $\Psi$ is a set which is complete in $L^{1}[0,2 \pi]$. For instance, such are the families $\chi_{I}$ and $\chi_{F o}$.

Theorem 5.7.3.1 $[\operatorname{Po}(1985)]$ Let $f \in A(\rho(r))$. The following assertions are equivalent:
a) $\overline{\mathcal{F}}[f g]=\overline{\mathcal{F}}[f]+\overline{\mathcal{F}}[g], \forall \mathcal{F} \in \chi_{\Psi}$,
b) $\underline{\mathcal{F}}[f g]=\underline{\mathcal{F}}[f]+\underline{\mathcal{F}}[g], \forall \mathcal{F} \in \chi_{\Psi}$, for all entire functions $g \in A(\rho(r))$.
c) $f$ is a GRG-function.

Proof of sufficiency of assertion c). Let us prove c) $\Longrightarrow$ a) and c) $\Longrightarrow b$ ). Using Theorem 5.7.1.3 we obtain for every characteristic $\mathcal{F}$

$$
\begin{equation*}
\overline{\mathcal{F}}[f g]=\sup \{\mathcal{F}(w): w \in \mathbf{F r}[f g]\} . \tag{5.7.3.2}
\end{equation*}
$$

Because of Theorem 3.1.2.4 fru1),

$$
\mathbf{F r}[f g] \subset \mathbf{F r}[f]+\mathbf{F r}[g] .
$$

Since $f$ is a CRG-function, $\mathbf{F r}[f]$ consists of only one subharmonic function $v_{\text {reg }}$ (see Theorem 5.6.1.1) and it is easy to check that in this case we have equality

$$
\mathbf{F r}[f g]=v_{\mathrm{reg}}+\mathbf{F r}[g] .
$$

Exercise 5.7.3.1 Check this.

Since $\mathcal{F}\left(v_{\text {reg }}+v_{g}\right)=\mathcal{F}\left(v_{\text {reg }}\right)+\mathcal{F}\left(v_{g}\right)$, we obtain

$$
\overline{\mathcal{F}}[f g]=\mathcal{F}\left(v_{\mathrm{reg}}\right)+\sup \left\{\mathcal{F}\left(v_{g}\right): v_{g} \in \mathbf{F r}[g]\right\}=\overline{\mathcal{F}}[f]+\overline{\mathcal{F}}[g]
$$

So c) $\Longrightarrow$ a) was proved. In the same way one can prove $c) \Longrightarrow b$ ).
Exercise 5.7.3.2 Prove this.
In the proof of sufficiency for c ) of a ) and b ) we can suppose that $\psi$ belong to the space $\mathcal{D}(\mathbb{T})$ of infinitely differentiable functions on the unit circle $\mathbb{T}$ because $\mathcal{D}(\mathbb{T})$ is complete in $L^{1}[0,2 \pi]$. We prove now sufficiency of b ) in Theorem 5.7.3.1.

We recall that (see (3.1.2.4a))

$$
v_{[t]}(z)=v(t z) t^{-\rho}, v \in U[\rho]
$$

to distinguish it from $(\bullet)_{t}$ that we define as

$$
u_{t}(z)=u(t z) t^{-\rho(t)}, u \in S H(\rho(r))
$$

The main constructive element for the proof sufficiency of b) in Theorem 5.7.3.1 is

Lemma 5.7.3.2 Let $\psi^{0} \in \mathcal{D}(S)$. There exists $v \in U[\rho]$ with the following properties:

$$
\begin{gather*}
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0} v_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty} v_{[t]}=\tilde{v},  \tag{5.7.3.3}\\
\left\langle v_{[t]}\left(e^{i \bullet}\right), \psi^{0}\right\rangle>\left\langle v\left(e^{i \bullet}\right), \psi^{0}\right\rangle \text { for } t \in(0, \infty), t \neq 1,  \tag{5.7.3.4}\\
 \tag{5.7.3.5}\\
\left\langle\tilde{v}\left(e^{i \bullet}\right), \psi^{0}\right\rangle>\left\langle v\left(e^{i \bullet}\right), \psi^{0}\right\rangle
\end{gather*}
$$

Proof of Lemma 5.7.3.2. Let $\psi^{0}$ be represented by Fourier series

$$
\psi^{0}(\phi)=\frac{a_{0}}{2}+\sum_{n=1}^{\infty}\left(a_{n} \cos n \theta+b_{n} \sin n \theta\right)
$$

Since $\psi^{0} \not \equiv 0$ there exists $a_{k} \neq 0$ or $b_{k} \neq 0$. Suppose there exists $a_{k} \neq 0$. In the proof we will consider three cases:

1. $k=0$;
2. $k \neq 0 \wedge k \leq p$;
3. $k \geq p+1$.

The number $\rho$ is supposed non-integer and $p=[\rho]$.
Consider the case $a_{0} \neq 0, a_{0}>0$. Set

$$
\begin{align*}
\psi(x) & :=\log \left(-e^{-\alpha|x|}+C\right), \alpha>0, C>1 \\
v(z) & :=|z|^{\rho} e^{\psi(\log |z|)}=\exp (\rho \log r+\psi(\log r)) \tag{5.7.3.6}
\end{align*}
$$

Applying the Laplace operator, we obtain:

$$
\begin{align*}
\Delta v & =\frac{1}{r^{2}} r \frac{\partial}{\partial r} r \frac{\partial}{\partial r} v(r)=e^{-2 x} \frac{\partial^{2}}{\partial x^{2}} e^{\rho x+\psi(x)}  \tag{5.7.3.7}\\
& =\exp ((\rho-2) x+\psi(x))\left[\left(\rho+\psi^{\prime}(x)\right)^{2}+\psi^{\prime \prime}(x)\right], x=\log r
\end{align*}
$$

Since

$$
\psi^{\prime}(x)=\alpha \operatorname{sgn} x \exp (-\alpha|x|) \rightarrow 0, \psi^{\prime \prime}(x)=-\alpha^{2} \operatorname{sgn} x \exp (-\alpha|x|) \rightarrow 0
$$

as $x \rightarrow \pm \infty$, it is possible to choose $\alpha$ such that the expression (5.7.3.7) is positive. So $v(z)$ is subharmonic.

It is easy to check that all the assertions of the lemma are satisfied and $\tilde{v}(z)=b|z|^{\rho}$ where $b(>0)$ is a constant.

Exercise 5.7.3.3 Check this.
If $a_{0}=-\left|a_{0}\right|<0$, consider the function

$$
v^{0}(z):= \begin{cases}\log |z|, & |z| \geq 1 \\ 0, & |z|<1\end{cases}
$$

it is subharmonic and

$$
\begin{equation*}
\left\langle v_{[t]}^{0}\left(e^{i \bullet}\right), \psi^{0}\right\rangle=a_{0} t^{-\rho} \log ^{+} t \tag{5.7.3.8}
\end{equation*}
$$

Since the RHS of (5.7.3.8) is minimized for $t_{0}=e^{\rho^{-1}}$, the function

$$
v(z):=v_{\left[t_{0}-1\right]}^{0}(z)
$$

satisfies the assertions of the lemma with $\tilde{v}=\lim _{t \rightarrow 0, \infty} v_{[t]}=0$
Now let $a_{0}=0, a_{k} \neq 0,0<k<p$. We will search for a function $v$ of the form

$$
\begin{equation*}
v\left(r e^{i \phi}\right):=\int_{0}^{2 \pi} G_{p}\left(r e^{i(\phi-\theta)}\right)\left(1-\operatorname{sgn} a_{k} \cos k \theta\right) d \theta \tag{5.7.3.9}
\end{equation*}
$$

This is the convolution $G_{p}\left(r e^{i \bullet}\right) * g$ of the primary kernel (see Section 2.9.1)

$$
G_{p}(z)=\log |1-z|+\Re \sum_{n=1}^{p} z^{n} / n
$$

with a positive function $g(\theta):=\left(1-\operatorname{sgn} a_{k} \cos k \theta\right)$ on the circle. So it is subharmonic. Recall that the cos-Fourier coefficients of the function $G_{p}\left(r e^{i \theta}\right)$ are (see Exercise 2.3.7.2).

$$
\hat{G}_{p}(m, r)=\left\{\begin{array}{ll}
0, & m=0,1, \ldots, p  \tag{5.7.3.10}\\
(1 / m) r^{m}, & m=p+1, \ldots
\end{array} \quad \text { if } r \leq 1\right.
$$

and

$$
\hat{G}_{p}(m, r)= \begin{cases}\log r, & m=0  \tag{5.7.3.11}\\ \frac{1}{m}\left(r^{m}-r^{-m}\right), & m=1, \ldots, p \quad \text { if } r \geq 1 \\ (1 / m) r^{m}, & m=p+1, \ldots\end{cases}
$$

All the sin-Fourier coefficients are equal to zero. The Fourier coefficients of the function $g$ are 1 and $-\operatorname{sgn} a_{k}$.

Using well-known properties of Fourier coefficients, we obtain for $0<k \leq p$,

$$
\begin{aligned}
\hat{v}_{[t]}(0) & = \begin{cases}0, & t \leq 1, \\
\frac{\log t}{t \rho}, & t \geq 1,\end{cases} \\
\hat{v}_{[t]}(k) & = \begin{cases}0, & t \leq 1, \\
-\frac{1}{k} \frac{t^{k}-t^{-k}}{t^{\rho}} \operatorname{sgn} a_{k} & t \geq 1,\end{cases} \\
\left\langle v_{[t]}\left(e^{i \bullet}\right), \psi^{0}\right\rangle & = \begin{cases}0, & t \leq 1, \\
-1 / k\left(t^{k-\rho}-t^{-k-\rho}\right)\left|a_{k}\right| & t \geq 1 .\end{cases}
\end{aligned}
$$

The function $t \mapsto\left\langle v_{t}\left(e^{i \bullet}\right), \psi^{0}\right\rangle$ tends to zero when $t \rightarrow 0, \infty$ and has its only minimum at the point

$$
t_{0}=\left(\frac{\rho+k}{\rho-k}\right)^{1 / k}
$$

Thus $v_{\left[\left(t_{0}\right)^{-1}\right]}$ satisfies the conditions of the lemma with $\tilde{v}=0$.
For $k \geq p+1$ we should take the same $g$ and then

$$
\left\langle v_{t}\left(e^{i \bullet}\right), \psi^{0}\right\rangle= \begin{cases}-(1 / k) t^{k-\rho}\left|a_{k}\right|, & t \leq 1 \\ -(1 / k) t^{-k-\rho}\left|a_{k}\right| & t \geq 1\end{cases}
$$

So the corresponding function $t \mapsto\left\langle v_{t}\left(e^{i \bullet}\right), \psi^{0}\right\rangle$ obtains its minimum at the point $t_{0}=1$ and the function $v$ satisfies the assertions of the lemma with $\tilde{v}=0$.
Exercise 5.7.3.4 Prove the lemma for the case $b_{k} \neq 0$.
Lemma 5.7.3.3 Let $v \in U[\rho]$ with the condition

$$
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0} v_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty} v_{[t]}=\tilde{v}
$$

fulfilled, and let $u \in S H(\rho(r))$ with some $v^{0} \in \operatorname{Fr}[u]$. Then there exists $w^{0} \in$ $S H(\rho(r))$ such that

$$
\begin{equation*}
\mathbf{F r}\left[w^{0}\right]=\left\{v_{[t]}: t \in(0, \infty)\right\} \cup \tilde{v} \tag{5.7.3.12}
\end{equation*}
$$

and the following condition holds:

1. If the sequence $\lim _{t_{n} \rightarrow \infty} w_{t_{n}}^{0}=v_{[t]}$ for some $t \in(0, \infty)$ and the sequence $u_{t_{n}}$ converges in $\mathcal{D}^{\prime}$ as $t_{n} \rightarrow \infty$, then $\lim _{n \rightarrow \infty} u_{t_{n}}=v_{[t]}^{0}$.
For the proof see Corollary 4.4.1.4.

Lemma 5.7.3.4 Let $w \in S H(\rho(r)), \psi \in \mathcal{D}(S)$. Then the following holds:

$$
\begin{aligned}
\liminf _{t \rightarrow \infty}\langle w, \psi\rangle & =\min _{v \in \mathbf{F r} w}\langle v, \psi\rangle \\
\limsup _{t \rightarrow \infty}\langle w, \psi\rangle & =\max _{v \in \mathbf{F r} w}\langle v, \psi\rangle .
\end{aligned}
$$

Exercise 5.7.3.5 Prove this exploiting completeness of Fr.
Proof of sufficiency of b) in Theorem 5.7.3.1. In assumption b) we should prove that $f$ is a CRG-function, i.e., by Theorem 5.6.1.1 its $\operatorname{Fr}[f]$ consists of only one function. Since $\log |f| \in S H(\rho(r))$ and because of Theorem 5.3.1.4 (Approximation) it is enough to prove the corresponding theorem for subharmonic functions. Suppose

$$
\begin{equation*}
\underline{\mathcal{F}}[u+w]=\underline{\mathcal{F}}[u]+\underline{\mathcal{F}}[w], \forall \mathcal{F} \in \chi_{\Psi} \tag{5.7.3.13}
\end{equation*}
$$

for all $w \in S H(\rho(r)$. We exploit Lemma 5.7.3.4 and write (5.7.3.13) in the form:

$$
\min _{v \in \mathbf{F r}[u+v]}\langle v, \psi\rangle=\min _{v \in \mathbf{F r} u}\langle v, \psi\rangle+\min _{v \in \mathbf{F r} w}\langle v, \psi\rangle, \forall \psi \in \Psi .
$$

Suppose the contrary, i.e., $u$ is not a CRG-function and $\operatorname{Fr} u$ does not consist of only one $v_{\min } \in U[\rho]$. Then there exists $v^{0} \neq v_{\min }$. The family $\chi_{\Psi}$ is total; therefore there exists $\psi^{0} \in \Psi$ such that $\left\langle v^{0}, \psi^{0}\right\rangle \neq\left\langle v_{\min }, \psi^{0}\right\rangle$ and hence

$$
\begin{equation*}
\left\langle v^{0}, \psi^{0}\right\rangle>\left\langle v_{\min }, \psi^{0}\right\rangle \tag{5.7.3.14}
\end{equation*}
$$

Using Lemma 5.7.3.2, construct for the function $\psi^{0}$ a function $v \in U[\rho]$ satisfying the conditions (5.7.3.3), (5.7.3.4) and (5.7.3.5). Apply Lemma 5.7.3.3 to construct a function $w^{0}$ satisfying (5.7.3.12) and the condition 1 . Under conditions of the theorem,

$$
\begin{equation*}
\min _{\omega \in \operatorname{Fr}\left(u+w^{0}\right)}\left\langle\omega, \psi^{0}\right\rangle=\min _{\omega \in \mathbf{F r}(u)}\left\langle\omega, \psi^{0}\right\rangle+\min _{\omega \in \mathbf{F r}\left(w^{0}\right)}\left\langle\omega, \psi^{0}\right\rangle . \tag{5.7.3.15}
\end{equation*}
$$

Let $\gamma \in \operatorname{Fr}\left(u+w^{0}\right)$ be the function on which the minimum of LRH in (5.7.3.15) is attained. Using (5.7.3.4), (5.7.3.5) and (5.7.3.12), we can rewrite (5.7.3.15) in the form

$$
\begin{equation*}
\left\langle\gamma, \psi^{0}\right\rangle=\min _{\omega \in \mathbf{F r} u}\left\langle\omega, \psi^{0}\right\rangle+\left\langle v, \psi^{0}\right\rangle \tag{5.7.3.16}
\end{equation*}
$$

Since $\gamma \in \operatorname{Fr}\left(u+w^{0}\right), \gamma=\mathcal{D}^{\prime}-\lim _{n \rightarrow \infty}\left(u+w^{0}\right)_{t_{n}}$. Passing to subsequences, we can suppose that the sequences $\left\{u_{t_{n}}\right\}$ and $\left\{w_{t_{n}}^{0}\right\}$ have limits. Since $\mathbf{F r} w^{0}$ has the form (5.7.3.12), there are two possible cases : $w_{t_{n}}^{0} \rightarrow v_{[t]}, t \in(0, \infty)$ and $w_{t_{n}}^{0} \rightarrow \tilde{v}$.

Consider the first case. Because of condition 1 from Lemma 5.7.3.3, $u_{t_{n}} \rightarrow v_{[t]}^{0}$ and $\gamma=v_{[t]}^{0}+v_{[t]}$. Substituting this in (5.7.3.15), we obtain

$$
\left\langle v_{[t]}^{0}, \psi^{0}\right\rangle-\min _{\omega \in \operatorname{Fr} u}\left\langle\omega, \psi^{0}\right\rangle=\left\langle v, \psi^{0}\right\rangle-\left\langle v_{[t]}, \psi^{0}\right\rangle
$$

This equality leads to a contradiction because for $t=1$ it contradicts (5.7.3.14) and for $t \neq 1$ it contradicts (5.7.3.4).

Consider the second case, when $w_{t_{n}}^{0} \rightarrow \tilde{v}$. Denote $v^{2}=\lim _{n \rightarrow \infty} u_{t_{n}}$ and rewrite (5.7.3.15) in the form

$$
\left\langle v^{2}, \psi^{0}\right\rangle-\min _{\omega \in \mathbf{F r} u}\left\langle\omega, \psi^{0}\right\rangle=\left\langle v, \psi^{0}\right\rangle-\left\langle\tilde{v}, \psi^{0}\right\rangle
$$

The last equality contradicts (5.7.3.5).
Sufficiency of condition a) of Theorem 5.7.3.1 can be proved using the Lemmas 5.7.3.3, 5.7.3.4 and a variation of Lemma 5.7.3.2.

Lemma 5.7.3.2' Let $\psi^{0} \in \mathcal{D}(S)$. There exists $v \in U[\rho]$ with the following properties:

$$
\begin{gather*}
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0} v_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty} v_{[t]}=\tilde{v},  \tag{5.7.3.3'}\\
\left\langle v_{[t]}\left(e^{i \bullet}\right), \psi^{0}\right\rangle<\left\langle v\left(e^{i \bullet}\right), \psi^{0}\right\rangle \text { for } t \in(0, \infty), t \neq 1, \\
\left\langle\tilde{v}\left(e^{i \bullet}\right), \psi^{0}\right\rangle<\left\langle v\left(e^{\bullet \bullet}\right), \psi^{0}\right\rangle . \tag{5.7.3.5'}
\end{gather*}
$$

Exercise 5.7.3.6 Prove this lemma and sufficiency of a) in Theorem 5.7.3.1.
5.7.4 In this section we consider the question of summing the asymptotic characteristics connected with the functional (5.7.1.3), i.e., indicator and lower indicator. Recall that $f \in A(\rho(r))$ is completely regular on the ray $\{\arg z=\phi\}\left(f \in A_{\text {reg, }, \phi}\right)$ if

$$
\begin{equation*}
h_{f}(\phi)=\underline{h}_{f}(\phi) \tag{5.7.4.0}
\end{equation*}
$$

We are going to prove the following assertions:
Theorem 5.7.4.1 Let $f \in A_{\mathrm{reg}, \phi}$. Then for every $g \in A(\rho(r))$,

$$
\begin{align*}
& h_{f g}(\phi)=h_{f}(\phi)+h_{g}(\phi),  \tag{5.7.4.1}\\
& \underline{h}_{f g}(\phi)=\underline{h}_{f}(\phi)+\underline{h}_{g}(\phi) . \tag{5.7.4.2}
\end{align*}
$$

Theorem 5.7.4.2 Suppose the equality (5.7.4.1) holds for every $g \in A(\rho(r))$. Then $f \in A_{\text {reg }, \phi}$.

Let us note that the assertion of Theorem 5.7.4.2 holds also if the equality (5.7.4.1) fulfilled for some sequence $\phi_{n} \rightarrow \phi$, because the indicator is a continuous function (see Section 3.2.5). So if the equality (5.7.4.1) holds for the set $\Phi$ of $\phi$ that is dense in $[0,2 \pi$ ) (or the set

$$
\begin{equation*}
e^{i \Phi}:=\left\{e^{i \phi}: \phi \in \Phi\right\} \tag{5.7.4.3}
\end{equation*}
$$

is dense on the unit circle), then $f \in A_{\text {reg, }, \phi}$ for all $\phi$, i.e., $f$ is a CRG-function.

On the other hand, the following assertion holds
Theorem 5.7.4.3 If the set $\Theta$ of $\theta$ is not dense in $[0,2 \pi)$, there exists $f \in A_{\text {reg, } \theta, \theta \in}$ $\Theta$ that is not a CRG-function.

The situation with a lower indicator is analogous, but in another topology.
A set $E$ is called non-rarefied at a point $z_{0}$ if for every function $v$ subharmonic in a neighborhood of $z_{0}$ the following holds:

$$
v\left(z_{0}\right)=\limsup _{z \in E, z \rightarrow z_{0}, z \neq z_{0}} v(z)=\limsup _{z \in E, z \rightarrow z_{0}} v(z)
$$

A set is rarefied if it is not non-rarefied.
Note that if $\underline{h}_{f}(\phi)=-\infty$, then $\underline{h}_{f g}(\phi)=-\infty$ for every $g \in A(\rho(r))$. It is obvious that $f \notin A_{\text {reg }, \phi}$.

The next theorems were proved in [GPS].
Theorem 5.7.4.4 Let (5.7.4.2) be fulfilled for $\psi \in E$ for all $g \in A(\rho(r))$ and $e^{i E}$ be non-rarefied at the point $e^{i \phi}$. Then $f \in A_{\text {reg }, \phi}$.

Theorem 5.7.4.5 Let $E_{0}$ be a set such that $e^{i E_{0}}$ is rarefied at all points of the unit circle. Then there exists $f \in A(\rho(r))$ for which (5.7.4.2) is fulfilled for all $\phi \in E_{0}$ and all $g \in A(\rho(r))$, but $f \notin A_{\mathrm{reg}, \phi}$ for all $\phi$ and $\underline{h}_{f}(\phi)>-\infty, \forall \phi$.

Let us note that $E_{0}$ can be dense in $[0,2 \pi)$ and $E$ from Theorem 5.7.4.4 can even be of zero measure.

The proof of Theorems 5.7.4.4 and 5.7.4.5 is based on the following assertion that gives a criterion for (5.7.4.2) in terms of limit sets $\operatorname{Fr}[f]$.

Theorem 5.7.4.6 Let $f \in A(\rho(r))$ and $\underline{h}_{f}(\phi)>-\infty$. The condition (5.7.4.2) holds for every $g \in A(\rho(r))$, such that $h_{g}(\phi)>-\infty$ iff

$$
\begin{equation*}
\liminf _{t \rightarrow 1} v\left(t e^{i \phi}\right)=\underline{h}_{f}(\phi) \tag{5.7.4.4}
\end{equation*}
$$

for all $v \in \mathbf{F r}[f]$.
An analogous criterion holds for (5.7.4.1).
Theorem 5.7.4.7 Let $f \in A(\rho(r))$. Then (5.7.4.1) holds for every $g \in A(\rho(r))$, iff

$$
\begin{equation*}
\limsup _{t \rightarrow 1} v\left(t e^{i \phi}\right)=h_{f}(\phi), \tag{5.7.4.5}
\end{equation*}
$$

for all $v \in \operatorname{Fr}[f]$.
Corollary 5.7.4.8 The equality (5.7.4.5) implies $f \in A_{\mathrm{reg}, \phi}$.

Actually, for every $v \in \operatorname{Fr}[f]$ we have, using semicontinuity of subharmonic functions and the definition (3.2.1.1) of the indicator,

$$
h_{f}(\phi)=\limsup _{t \rightarrow 1} v\left(t e^{i \phi}\right) \leq v\left(e^{i \phi}\right) \leq h_{f}(\phi)
$$

for all $v \in \mathbf{F r}[f]$. So $\operatorname{Fr}[f]$ consists of functions $v$ that coincide at the point $e^{i \phi}$ and hence on the ray $\left\{r e^{i \phi}: r \in(0, \infty)\right\}$.

Note also that the set $e^{i E}$ for which (5.7.4.1) holds is closed and Theorem 5.7.4.4 means that the set where (5.7.4.2) holds is thinly closed, i.e., closed in thin topology (see [Br, §6]).

Therefore if $e^{i \phi_{0}}$ is a limit point of $e^{i E}$ in the euclidian (respectively, thin) topology, then (5.7.4.1) ((5.7.4.2), respectively) is also a sufficient condition for completely regular growth at $\phi_{0}$.
5.7.5 The main constructive element for proving Theorem 5.7.4.6 is

Lemma 5.7.5.1 Let $\epsilon>0, t_{0}>0$ and $\phi_{0} \in[0,2 \pi)$ be fixed. Then there exists $v \in U[\rho]$ with the following properties:

$$
\begin{gather*}
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0} v_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty} v_{[t]}=0,  \tag{5.7.5.1}\\
v\left(e^{i \phi_{0}}\right)<v_{[t]}\left(e^{i \phi_{0}}\right), t \in(0,1) \cup(1, \infty),  \tag{5.7.5.2}\\
-\infty<v\left(e^{i \phi_{0}}\right)<-\epsilon, \tag{5.7.5.3}
\end{gather*}
$$

and the inequality

$$
\begin{equation*}
v_{[t]}\left(e^{i \phi_{0}}\right)-v\left(e^{i \phi_{0}}\right) \leq \epsilon / 2 \tag{5.7.5.4}
\end{equation*}
$$

implies

$$
\begin{equation*}
t \in\left[1 / t_{0}, t_{0}\right] \tag{5.7.5.5}
\end{equation*}
$$

The last condition means that the function $\psi(t):=v_{[t]}\left(e^{i \phi_{0}}\right)$ can be less than $\psi(1)+\epsilon / 2$ only in a neighborhood of $t=1$.

Proof. Set

$$
\begin{gather*}
w(z):=\max \left(\log \left|1-z e^{-i \phi_{0}}\right|,-N\right)+\Re \sum_{n=1}^{p} \frac{1}{n}\left(z e^{-i \phi_{0}}\right)^{n}, \\
N>0, p=[\rho] . \tag{5.7.5.6}
\end{gather*}
$$

It is obvious that $w$ is subharmonic, with masses $\nu_{w}$ concentrated in a neighborhood of the point $e^{i \phi_{0}}$. Thus $\nu_{w} \in \mathcal{M}[\rho]$ (see (3.1.3.4)) and

$$
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0}\left(\nu_{w}\right)_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty}\left(\nu_{w}\right)_{[t]}=0
$$

Hence (see Theorem 3.1.4.2) $w \in U[\rho]$, and (see (3.1.5.0))

$$
\begin{equation*}
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0} w_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty} w_{[t]}=0 \tag{5.7.5.7}
\end{equation*}
$$

Let us capitalize on the behavior of $w_{[t]}$ on the ray $\left\{\arg z=\phi_{0}\right\}$ :

$$
\begin{equation*}
w_{[t]}\left(e^{i \phi_{0}}\right):=\psi(t)=\left(\max (\log |1-t|,-N)+\Re \sum_{n=1}^{p} \frac{1}{n}\right) t^{-\rho} t^{n} \tag{5.7.5.8}
\end{equation*}
$$

It is possible to prove directly the following properties of $\psi(t)$.
i) outside interval $\left(1-e^{-N}, 1+e^{-N}\right), \psi(t)=G_{p}(t) t^{-\rho}$; where $G_{p}$ is the Primary Kernel (see Section 2.9.1) and inside this interval the first summand is $-N$;
ii) $\psi(t)>0$ for $t>t_{1}$ where $t_{1}$ is a zero of the equation $G_{p}(t)=0, \psi(t)$ decreases monotonically on the interval $\left(0,1-e^{-N}\right)$ and increases monotonically on the interval $\left(1-e^{-N}, t_{1}\right)$.

## Exercise 5.7.5.1 Prove this.

Now set $t_{2}:=1-e^{-N}$ and $v(z):=D w_{t_{2}}(z)$, where $D$ is a constant. This function satisfies the conditions (5.7.5.1) and (5.7.5.2) of the lemma and $v_{[t]}\left(e^{i \phi_{0}}\right)$ has only one negative minimum for $t=1$. Thus it is possible to take $D$ sufficiently large to satisfy the conditions (5.7.5.3) and (5.7.5.4) for fixed $\epsilon$ and $t_{0}$.

Exercise 5.7.5.2 Prove this in detail.
In the proof of Theorem 5.7.4.6 we also use Lemma 5.7.3.3. We can prove all the assertions for subharmonic functions from $S H(\rho(r))$.

Proof of Theorem 5.7.4.6. Necessity. We should prove that if the equality

$$
\begin{equation*}
\underline{h}\left(e^{i \phi_{0}}, u+w\right)=\underline{h}\left(e^{i \phi_{0}}, u\right)+\underline{h}\left(e^{i \phi_{0}}, w\right) \tag{5.7.5.9}
\end{equation*}
$$

holds for a fixed $u \in S H(\rho(r)), \phi_{0}$ and every $w \in S H(\rho(r))$, then

$$
\begin{equation*}
\liminf _{t \rightarrow 1} v\left(t e^{i \phi}\right)=\underline{h}\left(e^{i \phi}, u\right) \tag{5.7.5.10}
\end{equation*}
$$

for all $v \in \mathbf{F r} u$. Assume that $\underline{h}\left(e^{i \phi}, u\right)>-\infty$ and $\underline{h}\left(e^{i \phi_{0}}, w\right)>-\infty$. Suppose the contrary, i.e., there exists $v^{0} \in \mathbf{F r} u$ such that

$$
\begin{equation*}
\liminf _{t \rightarrow 1} v^{0}\left(t e^{i \phi_{0}}\right)>\underline{h}\left(e^{i \phi_{0}}, u\right) \tag{5.7.5.11}
\end{equation*}
$$

The inequality (5.7.5.11) implies that there exists $\epsilon>0$ and $t_{0}>0$ such that for every $t \in\left[1 / t_{0}, t_{0}\right]$ the inequality

$$
\begin{equation*}
v^{0}\left(t e^{i \phi_{0}}\right)>\underline{h}\left(e^{i \phi_{0}}, u\right)+\epsilon \tag{5.7.5.11a}
\end{equation*}
$$

holds. Let us construct by Lemma 5.7.5.1 for these $\epsilon, t_{0}, \phi_{0}$ a function $v$ and by Lemma 5.7.3.3 for the functions $u, v^{0}$ and the already found $v$ a function $w^{0}$. Let us show that for $w^{0}$ the equality (5.7.5.9) does not hold.

Compute $\underline{h}\left(e^{i \phi}, w^{0}\right)$. From (3.2.1.2)

$$
\underline{h}\left(e^{i \phi_{0}}, w^{0}\right)=\min \left\{0, \inf \left\{v_{[t]}\left(e^{i \phi_{0}}: t \in(0, \infty)\right\}\right\}\right.
$$

The inequalities (5.7.5.3) imply that 0 can be omitted and (5.7.5.2) implies that the infimum is attained at $t=1$, i.e.,

$$
\begin{equation*}
\underline{h}\left(e^{i \phi_{0}}, w^{0}\right)=v\left(e^{i \phi_{0}}\right) . \tag{5.7.5.12}
\end{equation*}
$$

Find $v^{\epsilon} \in \operatorname{Fr}\left(u+w^{0}\right)$ such that $\underline{h}\left(e^{i \phi_{0}}, u+w^{0}\right)>v^{\epsilon}\left(e^{i \phi_{0}}\right)-\epsilon / 3$. Let $t_{n} \rightarrow \infty$ and $\left(u+w^{0}\right)_{t_{n}} \rightarrow v^{\epsilon}$ in $\mathcal{D}^{\prime}$. Passing to subsequences we can assume that $u_{t_{n}}$ and $w_{t_{n}}^{0}$ also converge. Consider two cases. The first, when

$$
\begin{equation*}
\mathcal{D}^{\prime}-\lim w_{t_{n}}^{0}=v_{[t]}, t \in(0, \infty) \tag{5.7.5.13}
\end{equation*}
$$

By Lemma 5.7.3.3 $\lim u_{t_{n}}=v_{[t]}^{0}$ and hence $v^{\epsilon}=\lim \left(w^{0}+u\right)_{t_{n}}=v_{[t]}+v_{[t]}^{0}$. If $t \notin\left[1 / t_{0}, t_{0}\right]$, then by (5.7.5.4)

$$
\begin{equation*}
v_{[t]}\left(e^{i \phi_{0}}\right)>v\left(e^{i \phi_{0}}\right)+\epsilon / 2=\underline{h}\left(e^{i \phi_{0}}, w^{0}\right)+\epsilon / 2 . \tag{5.7.5.14}
\end{equation*}
$$

In this case we have

$$
\begin{equation*}
\underline{h}\left(e^{i \phi_{0}}, u+w^{0}\right) \geq v^{\epsilon}\left(e^{i \phi_{0}}\right)-\epsilon / 3 \geq v_{[t]}+v_{[t]}^{0}-\epsilon / 3 . \tag{5.7.5.14a}
\end{equation*}
$$

Using (5.7.5.14a), we obtain

$$
\begin{equation*}
\underline{h}\left(e^{i \phi_{0}}, u+w^{0}\right) \geq \underline{h}\left(e^{i \phi_{0}}, w^{0}\right)+\underline{h}\left(e^{i \phi_{0}}, u\right)+\epsilon / 6 . \tag{5.7.5.15}
\end{equation*}
$$

If $t \in\left[1 / t_{0}, t_{0}\right]$, then from (5.7.5.11) we have

$$
\begin{equation*}
\underline{h}\left(e^{i \phi_{0}}, u+w^{0}\right) \geq \underline{h}\left(e^{i \phi_{0}}, w^{0}\right)+\underline{h}\left(e^{i \phi_{0}}, u\right)+2 \epsilon / 3 \tag{5.7.5.16}
\end{equation*}
$$

So the case (5.7.5.13) is settled.
Let $\mathcal{D}^{\prime}-\lim w_{t_{n}}^{0}=0$. In this case we have
$\underline{h}\left(e^{i \phi_{0}}, u+w^{0}\right) \geq v^{\epsilon}\left(e^{i \phi_{0}}\right)-\epsilon / 3 \geq \underline{h}\left(e^{i \phi_{0}}, u\right)-\epsilon+\epsilon-\epsilon / 3=\underline{h}\left(e^{i \phi_{0}}, u\right)-\epsilon+2 \epsilon / 3$.
Using (5.7.5.12) and (5.7.5.3) we obtain

$$
\underline{h}\left(e^{i \phi_{0}}, u+w^{0}\right) \geq \underline{h}\left(e^{i \phi_{0}}, u\right)+\underline{h}\left(e^{i \phi_{0}}, w^{0}\right)+2 \epsilon / 3 .
$$

So we proved in any case that (5.7.5.9) does not hold if (5.7.5.10) does not hold.
Let us prove sufficiency in Theorem 5.7.4.6. We prove it for subharmonic functions. Let $u \in S H(\rho(r))$ and for every $v \in \operatorname{Fr} u$ (5.7.5.10) holds. Let us show that for all $w \in S H(\rho(r))(5.7 .5 .9)$ holds. It is sufficient to prove that

$$
\begin{equation*}
\underline{h}\left(e^{i \phi_{0}}, u+w\right) \leq \underline{h}\left(e^{i \phi_{0}}, u\right)+\underline{h}\left(e^{i \phi_{0}}, w\right) \tag{5.7.5.17}
\end{equation*}
$$

holds since the inverse inequality holds for every $w \in S H(\rho(r))$ (see (3.2.1.5)). Let us begin by noting that for every $v^{2} \in \mathbf{F r} w$ there exist $v \in \mathbf{F r}(u+w)$ and $v^{1} \in \mathbf{F r} u$ such that

$$
\begin{equation*}
v=v^{1}+v^{2} \tag{5.7.5.18}
\end{equation*}
$$

Indeed, let $t_{n} \rightarrow \infty$ be a sequence such that $w_{t_{n}} \rightarrow v^{2}$. We can suppose, in choosing a subsequence, that $u_{t_{n}} \rightarrow v^{1}$ and $(u+w)_{t_{n}} \rightarrow v$. Then (5.7.5.18) holds.

Let $\epsilon$ be arbitrarily small. Choose $v^{2} \in \operatorname{Fr} w$ such that $v^{2}\left(e^{i \phi}\right)<h\left(e^{i \phi}, w\right)+\epsilon$ holds. From upper semicontinuity of $v^{2}$ we have

$$
\begin{equation*}
\limsup _{t \rightarrow 1} v^{2}\left(e^{i \phi}\right) \leq \underline{h}\left(e^{i \phi_{0}}, w\right)+\epsilon \tag{5.7.5.19}
\end{equation*}
$$

Let $v^{1} \in \operatorname{Fr} u$ and $v \in \operatorname{Fr}(u+w)$ satisfy (5.7.5.18). Then we have

$$
\underline{h}\left(e^{i \phi_{0}}, u+w\right) \leq\left(v^{1}+v^{2}\right)_{[t]}\left(e^{i \phi_{0}}\right)=v_{[t]}^{1}\left(e^{i \phi_{0}}\right)+v_{[t]}^{2}\left(e^{i \phi_{0}}\right), \forall t .
$$

Hence

$$
\underline{h}\left(e^{i \phi_{0}}, u+w\right) \leq \liminf _{t \rightarrow 1} v_{[t]}^{1}\left(e^{i \phi_{0}}\right)+\limsup _{t \rightarrow 1} v_{[t]}^{2}\left(e^{i \phi_{0}}\right)
$$

Using (5.7.5.10) and (5.7.5.19) we obtain

$$
\underline{h}\left(e^{i \phi_{0}}, u+w\right) \leq \underline{h}\left(e^{i \phi_{0}}, u\right)+\underline{h}\left(e^{i \phi_{0}}, w\right)+\epsilon
$$

This proves the inverse inequality and hence the equality (5.7.5.9), because $\epsilon$ is arbitrarily small.
5.7.6 Now we are going to prove Theorem 5.7.4.4. We need the following assertion from Potential Theory.
Lemma 5.7.6.1 Let $E$ be a set that is non-rarefied at the point $e^{i \phi_{0}}$. Let $E^{\prime}$ be a set in $\mathbb{C}$, such that $\forall e^{i \phi} \in E$ and $\forall \delta>0$ there exists a point $z^{\prime} \in E^{\prime}$ on the ray $\{\arg z=\phi\}$ such that $\left|z^{\prime}-e^{i \phi}\right|<\delta$. Then $E^{\prime}$ is also non-rarefied at the point $e^{i \phi_{0}}$.

Proof. We can suppose without loss of generality that $E^{\prime}$ has no intersection with some neighborhood of zero. Denote by $P(z)$ the map $z \mapsto e^{i \arg z}$. It is easy to see that for all pairs $z_{1}^{\prime}, z_{2}^{\prime} \in E^{\prime}$ the inequality $\left|P\left(z_{1}^{\prime}\right)-P\left(z_{2}^{\prime}\right)\right|<A\left|z_{1}^{\prime}-z_{2}^{\prime}\right|$ holds for some constant $A$. Thus the logarithmic capacity (2.5.2.5) satisfies ([La, Ch. II, § 4, it. 11, 15]).

$$
\begin{equation*}
\operatorname{cap}_{l}(M)<A \operatorname{cap}_{l}\left(M^{\prime}\right) \tag{5.7.6.1}
\end{equation*}
$$

where $M^{\prime} \subset E^{\prime}, M=P\left(M^{\prime}\right)$. Now we exploit the following properties of nonrarefied sets. First, if $E$ is non-rarefied at a point $z_{0}$, then there exists a compact set that is non-rarefied at $z_{0}([\mathrm{La}, \mathrm{Ch} . \mathrm{V}, \S 1$, it. $5, \S 3$, it. 9$])$. Second, for a compact set $K$ that is non-rarefied at $z_{0}$,

$$
\begin{equation*}
\sum_{n=1}^{\infty} \frac{n}{\log \left(\mathbf{c a p}_{l} K_{n}\right)^{-1}}=\infty \tag{5.7.6.2}
\end{equation*}
$$

where $K_{n}:=K \cap\left\{z: q^{n+1} \leq\left|z-z_{0}\right| \leq q^{n}\right\}, 0<q<1$.

Using the inequality (5.7.6.1), we obtain that divergence of the series (5.7.6.2) for a compact $K \subset E$ implies divergence for $K^{\prime} \subset E^{\prime}$ where $K=P\left(K^{\prime}\right)$, i.e., $E^{\prime}$ is non-rarefied at the point $P\left(e^{i \phi_{0}}\right)=e^{i \phi_{0}}$.

Proof of Theorem 5.7.4.4. Let $\epsilon(\phi) \rightarrow 0$ as $\phi \rightarrow \phi_{0}$ and let $v \in \operatorname{Fr} u$. Suppose (5.7.5.9) holds for $e^{i \phi} \in E$. By Theorem 5.7.4.6 the equality (5.7.5.10) holds. Thus $\forall \Delta>0, \exists z^{\prime}=z^{\prime}\left(e^{i \phi}, \Delta\right)$ such that

$$
\begin{equation*}
\left|z^{\prime}-e^{i \phi}\right|<\Delta, \arg z^{\prime}=\phi, v\left(z^{\prime}\right)<\underline{h}\left(e^{i \phi}\right)+\epsilon(\phi) \tag{5.7.6.3}
\end{equation*}
$$

Set

$$
E^{\prime}:=\bigcup_{\phi \in E} \bigcup_{n=1}^{\infty} z^{\prime}\left(e^{i \phi}, 1 / n\right)
$$

By (5.7.6.3) and upper semicontinuity of $h\left(e^{i \phi}\right)$ we obtain

$$
\begin{equation*}
\limsup _{z^{\prime} \rightarrow e^{i \phi_{0}}, z^{\prime} \in E^{\prime}} v\left(z^{\prime}\right) \leq \underline{h}\left(e^{i \phi_{0}}\right) . \tag{5.7.6.4}
\end{equation*}
$$

Since $E^{\prime}$ is non-rarefied, by Lemma 5.7.6.1 the upper limit of $v$ coincides with $v\left(e^{i \phi_{0}}\right)$ and hence $v\left(e^{i \phi_{0}}\right) \leq \underline{h}\left(e^{i \phi_{0}}\right)$. The inverse inequality holds always. Thus $v\left(e^{i \phi_{0}}\right)=\underline{h}\left(e^{i \phi_{0}}\right), \forall v \in \mathbf{F r} u$. Hence $h\left(e^{i \phi_{0}}\right)=\underline{h}\left(e^{i \phi_{0}}\right)$.
5.7.7 Now we are going to prove Theorem 5.7.4.5. Before this we need to describe a construction and prove some auxiliary assertions.

Let $B_{j}:=\left\{z: T^{j}<|z|<T^{j+1}\right\}, j=0, \pm 1, \pm 2, \ldots$ where $T>1$ is a fixed number. Denote $L_{E_{0}}:=\left\{z: e^{i \arg z} \in e^{i E_{0}}\right\}$. Recall that $e^{i E_{0}}$ is a set rarefied at every point of the unit circle. Let $Q$ be the set of rational numbers on the interval $(1, T)$. Set

$$
\begin{aligned}
S_{Q} & :=\{z:|z| \in Q\}, \\
T^{j} S_{Q} & :=\left\{z T^{j}: z \in S_{Q}\right\}, \\
A_{j} & :=L_{E_{0}} \cap T^{j} S_{Q}, \quad j=0, \pm 1, \pm 2, \ldots
\end{aligned}
$$

Lemma 5.7.7.1 There exists $v \in U[\rho]$ such that

$$
\begin{equation*}
v(z)=-\infty \tag{5.7.7.1}
\end{equation*}
$$

for $z \in A_{0}$ and

$$
\begin{equation*}
\mu_{v}(e)=0, \forall e \subset \mathbb{C} \backslash B_{0} . \tag{5.7.7.2}
\end{equation*}
$$

Proof. The set $E$ is rarefied at every point, hence it is polar ([Br, Ch. 7, §4]). Thus the set $\{z:|z|=r\} \cap L_{E_{0}}$ is polar (see [Br, Ch. 3, §2]). A countable union of polar sets is polar $([\mathrm{Br}, \mathrm{Ch} .3, \S 2])$. Thus $A_{0}$ is polar.Hence there exists a positive measure $\mu$ concentrated on $B_{0}$ for which the potential $v(z):=\int G_{p}(z / \zeta) \mu(d \zeta)$ is equal to $-\infty$ on $A_{0}$ (see [Br, Ch. $4, \S 6$, Applications]). It is easy to see that $\mu \in \mathcal{M}(\rho)$ and hence $v \in U[\rho]$ (see Theorem 3.1.4.2).

Lemma 5.7.7.2 There exists $\omega \in U[\rho]$ such that the following conditions are fulfilled:

$$
\begin{equation*}
\omega(z)=-\infty, z \in A:=\cup_{j=-\infty}^{+\infty} A_{j} ; \omega(T z)=T^{\rho} \omega(z) \tag{5.7.7.3}
\end{equation*}
$$

Proof. Set for every $E \Subset \mathbb{C} \backslash 0$,

$$
\begin{equation*}
\nu(E):=\sum_{j=-\infty}^{j=+\infty} T^{j \rho} \mu_{v}\left(T^{-j} E \cup B_{0}\right) \tag{5.7.7.4}
\end{equation*}
$$

(compare Theorem 4.1.7.1). We have $\nu \in \mathcal{M}(\rho)$. Set

$$
\omega(z):=\int G_{p}(z / \zeta) \nu(d \xi d \eta), \zeta=\xi+i \eta
$$

This $\omega$ satisfies (5.7.7.3).
Exercise 5.7.7.1 Prove this using Theorem 4.1.7.1.
Lemma 5.7.7.3 Let $\omega$ be a subharmonic function in $\mathbb{C}$. Denote

$$
m(\phi):=\max \left\{\omega\left(r e^{i \phi}\right): r \in[1, T]\right\}
$$

Then there exists a constant $C>-\infty$ such that $m(\phi)>C \forall \phi$.

Proof. If not, there exists a sequence $\phi_{n}$ that we can assume to converge to $\phi_{\infty}$ such that $m\left(\phi_{n}\right) \rightarrow-\infty$. By upper semicontinuity of $\omega$ we have $\omega(z)=$ $-\infty, z e^{-i \phi_{\infty}} \in[1, T]$. Thus $\omega(z) \equiv-\infty$ because the capacity of the segment in the plane is positive and hence it is not polar for some subharmonic function.

Recall that for $v \in U[\rho]$ (see (4.1.3.1))

$$
\begin{align*}
\mathbb{C}(v) & :=\mathcal{D}^{\prime}-\operatorname{clos}\left\{v_{[t]}: 0<t<\infty\right\},  \tag{5.7.7.5}\\
\Omega(v) & :=\left\{v^{\prime} \in U[\rho]:\left(\exists t_{k} \rightarrow \infty\right)\left(v^{\prime}=\lim _{k \rightarrow \infty} v_{\left[t_{k}\right]}\right\},\right.  \tag{5.7.7.6}\\
A(v) & :=\left\{v^{\prime} \in U[\rho]:\left(\exists \tau_{k} \rightarrow 0\right)\left(v^{\prime}=\lim _{k \rightarrow \infty} v_{\left[t_{k}\right]}\right\} .\right. \tag{5.7.7.7}
\end{align*}
$$

By Theorems 4.1.3.3 and 4.2.1.2, if

$$
\begin{equation*}
A(v) \cap \Omega(v) \neq \varnothing \tag{5.7.7.8}
\end{equation*}
$$

there exists $u \in S H(\rho(r))$ such that

$$
\begin{equation*}
\operatorname{Fr} u=\mathbb{C}(v) \tag{5.7.7.9}
\end{equation*}
$$

Lemma 5.7.7.4 There exists $v^{1} \in U[\rho]$ such that the following holds:

$$
\begin{gather*}
A\left(v^{1}\right)=\Omega\left(v^{1}\right)  \tag{5.7.7.10}\\
\inf \left\{v\left(e^{i \phi}\right): v \in \mathbb{C}\left(v^{1}\right)\right\}=\liminf _{t \rightarrow 1} v\left(t e^{i \phi}\right)=0, \forall v \in \mathbb{C}\left(v^{1}\right), \forall e^{i \phi} \in e^{i E_{0}},  \tag{5.7.7.11}\\
\sup \left\{v\left(e^{i \phi}\right): v \in \mathbb{C}\left(v^{1}\right)\right\} \neq \inf \left\{v\left(e^{i \phi}\right): v \in \mathbb{C}\left(v^{1}\right)\right\} \tag{5.7.7.12.}
\end{gather*}
$$

Proof. Let $\omega(z)$ be constructed by Lemma 5.7.7.2. Set

$$
v(z):=\omega(z)+D \log ^{+} 2|z| .
$$

The condition (5.7.7.3) implies

$$
A(\omega)=\Omega(\omega)=\left\{\omega_{[t]}: t \in[1, T]\right\}
$$

because it is a Periodic Limit Set (see Theorem 4.1.7.1).
Since $\left(\log ^{+} 2|z|\right)_{[t]} \rightarrow 0, t \rightarrow 0, t \rightarrow \infty$, the function $v$ satisfies the condition

$$
A(v)=\Omega(v)=\left\{\omega_{[t]}: t \in[1, T]\right\}
$$

By Theorem 2.1.7.4 for the function $v^{1}:=v^{+}$we have

$$
A\left(v^{1}\right)=\Omega\left(v^{1}\right)=\left\{\omega_{[t]}^{+}: t \in[1, T]\right\}
$$

Note that $v^{1}(z)=0$ for $z \in A$ and since $A$ is dense in $L_{E_{0}}$ (5.7.7.11) holds. Choosing $D$ sufficiently large it is possible (using Lemma 5.7.7.3) to find on every $\operatorname{ray}\{\arg z=\phi\}$ a point $z_{\phi}$ where $v^{1}\left(z_{\phi}\right)>0$. Hence $\sup \left\{v\left(e^{i \phi}\right): v \in \mathbb{C}\left(v^{1}\right)\right\}>0$. Because of (5.7.7.11) and upper semicontinuity of $\inf \left\{v\left(e^{i \phi}\right): v \in \mathbb{C}\left(v^{1}\right)\right\}$ it is zero for every $e^{i \phi}$. Thus (5.7.7.12) holds.

Proof of Theorem 5.7.4.5. Let us construct by Theorems 4.1.3.3 and 4.2.1.2 a function $u \in S H\left(\rho(r)\right.$ such that $\operatorname{Fr} u=\mathbb{C}\left(v^{1}\right)$ where $v^{1}$ is taken from Lemma 5.7.7.4. It does not belong to $A_{\mathrm{reg}, \phi}$ for any $\phi$. The equality (5.7.5.9) holds for every $\phi \in E_{0}$ because of (5.7.7.11) by Theorem 5.7.4.6.
5.7.8 The proof of Theorem 5.7.4.1 is a copy of the proof of sufficiency of assertion c) in Theorem 5.7.3.1.

Exercise 5.7.8.1 Prove Theorem 5.7.4.1.
Now we are going to prove Theorem 5.7.4.7 which implies (as it was shown in Corollary 5.7.4.8) Theorem 5.7.4.2.

The main constructive element of the proof of necessity is
Lemma 5.7.8.1 Let $\epsilon>0, t_{0}>0$ and $\phi_{0} \in[0,2 \pi)$ be fixed. Then there exists $v \in U[\rho]$ with the following properties:

$$
\begin{gather*}
\mathcal{D}^{\prime}-\lim _{t \rightarrow 0} v_{[t]}=\mathcal{D}^{\prime}-\lim _{t \rightarrow \infty} v_{[t]}=0,  \tag{5.7.8.1}\\
v\left(e^{i \phi_{0}}\right)>v_{[t]}\left(e^{i \phi_{0}}\right), t \in(0,1) \cap(1, \infty), \tag{5.7.8.2}
\end{gather*}
$$

and the inequality

$$
\begin{equation*}
v_{[t]}\left(e^{i \phi_{0}}\right)-v\left(e^{i \phi_{0}}\right) \geq-\epsilon / 2 \tag{5.7.8.3}
\end{equation*}
$$

implies

$$
\begin{equation*}
t \in\left[t_{0}, 1 / t_{0}\right] \tag{5.7.8.4}
\end{equation*}
$$

The last condition means that the function $\psi(t):=v_{[t]}\left(e^{i \phi_{0}}\right)$ can be more than $\psi(1)-\epsilon / 2$ only in a neighborhood of $t=1$.

Proof. Consider the function

$$
\begin{equation*}
w(z):=\log ^{+}|z| . \tag{5.7.8.5}
\end{equation*}
$$

It is subharmonic and satisfies (5.7.8.1). Since the function

$$
\psi(t):=w_{[t]}\left(e^{i \phi_{0}}\right)=t^{-\rho} \log ^{+} t
$$

has its only strict maximum in the point $t_{\max }>1$, the function

$$
v(z):=w\left(z / t_{\max }\right)
$$

has all the properties (5.7.8.1)-(5.7.8.4).
After this lemma all the proof of Theorem 5.7.4.6 can be repeated with minimal changes.

Exercise 5.7.8.2 Prove Theorem 5.7.4.7.
5.7.9 Now we are going to prove Theorem 5.7.4.3. Let us prove the following

Lemma 5.7.9.1 Let $\bar{\Theta}$ be a closed subset of $[0,2 \pi)$. Then for every $\sigma>0$ there exists a $2 \pi$-periodic $\rho$-trigonometrically convex function $h(\phi)$ such that

$$
\begin{equation*}
h(\phi)=\sigma \tag{5.7.9.1}
\end{equation*}
$$

for $\phi \in \bar{\Theta}$ and

$$
\begin{equation*}
h(\phi)>\sigma \tag{5.7.9.2}
\end{equation*}
$$

for $\phi \notin \bar{\Theta}$.

Proof. We can suppose that $0 \in \bar{\Theta}$, otherwise we can shift it a little. The set $[0,2 \pi) \backslash \bar{\Theta}$ is open and it can be represented as the union of non-intersecting open intervals. If length of an interval is $\leq \pi / \rho$ we can construct a $\rho$-trigonometrical function that is equal to $\sigma$ on the ends of the interval. It is greater than $\sigma$ in all inner points of the interval because $f(\phi) \equiv \sigma$ is a strictly $\rho$-trigonometrical function. If the length of the interval is greater than $\pi / \rho$, for example $(-l / 2, l / 2)$ with $l>\pi / 2 \rho$, we cover it by intersecting intervals of length less then $\pi / \rho$, construct $h_{I}(\phi)$ as before for every interval $I$ and set $h(\phi)=\max _{I} h_{I}(\phi)$. It is obvious that $h(\phi)$ is greater than $\sigma$ and it is $\rho$-trigonometrically convex.

Theorem 5.7.4.3 is a corollary of Lemma 5.7.9.1 and the following
Theorem 5.7.9.2 Let $h_{1}$ and $h_{2}$ be two $\rho$-trigonometrically convex functions. Then there exists a function $f \in A(\rho(r))$ such that

$$
h_{f}(\phi)=\max \left(h_{1}(\phi), h_{2}(\phi)\right), \underline{h}_{f}(\phi)=\min \left(h_{1}(\phi), h_{2}(\phi)\right) .
$$

Proof. Consider the set

$$
\begin{equation*}
U:=\left\{v(z)=c r^{\rho} h_{1}(\phi)+(1-c) r^{\rho} h_{2}(\phi): 0 \leq c \leq 1\right\} . \tag{5.7.9.3}
\end{equation*}
$$

It consists of invariant subharmonic functions and is contained in $U[\rho]$ and satisfies the condition of Theorem 4.1.4.1. Hence (Theorems 4.2.1.2, Corollary 5.3.1.5) there exists a function $f \in A(\rho(r))$ such that

$$
\begin{equation*}
\operatorname{Fr} f=U . \tag{5.7.9.4}
\end{equation*}
$$

By formulae (3.2.1.1), (3.2.1.2) we obtain the assertion of the theorem, using (5.7.9.3).

Exercise 5.7.9.1 Prove Theorem 5.7.4.3.
5.7.10 The family of characteristics $\left\{\mathcal{F}_{\alpha}, \alpha \in A\right\}$ is called independent if for every subset $A^{\prime} \subset A$ (or subset in some class of subsets, for example, measurable or closed) there exists a function $f=f_{A^{\prime}} \in A(\rho(r))$ such that

$$
\begin{aligned}
& \underline{\mathcal{F}}_{\alpha}[f]=\overline{\mathcal{F}}_{\alpha}[f], \quad \alpha \in A^{\prime} \\
& \underline{\mathcal{F}}_{\alpha}[f] \neq \overline{\mathcal{F}}_{\alpha}[f], \quad \alpha \in A \backslash A^{\prime} .
\end{aligned}
$$

It means that for every pointed subset of characteristics there exists a function that has regular growth with respect to this subset of characteristics and is not of regular growth with respect to all other characteristics.

Theorem 5.7.4.3 can be considered as an assertion of independence of the family (5.7.2.2).

Theorem 5.7.10.1 The family $\chi_{F o}$ (5.7.2.4) is independent.
I.e., for every $A \subset \mathbb{Z}$ there exists $f \in A(\rho(r)$ such that

$$
\lim _{r \rightarrow \infty} r^{-\rho(r)} \int_{0}^{2 \pi} \log \mid f\left(r e^{i \phi}\right) g_{k}(\phi) d \phi
$$

exists for all $k \in A$ and does not exist for $k \in \mathbb{Z} \backslash A$. To begin we prove
Lemma 5.7.10.2 There exist two $\rho$-trigonometrically convex functions $h_{1}$ and $h_{2}$ for which

$$
\begin{align*}
& \int_{0}^{2 \pi} h_{1}(\phi) g_{k}(\phi) d \phi=\int_{0}^{2 \pi} h_{2}(\phi) g_{k}(\phi) d \phi, k \in A  \tag{5.7.10.1}\\
& \int_{0}^{2 \pi} h_{1}(\phi) g_{k}(\phi) d \phi \neq \int_{0}^{2 \pi} h_{2}(\phi) g_{k}(\phi) d \phi, k \in \mathbb{Z} \backslash A \tag{5.7.10.2}
\end{align*}
$$

Proof. Let $g(\phi) \in C^{2}$ be a function, the Fourier coefficients of which with indices $k \in A$ are equal to zero. We can represent it as a difference of $\rho$-trigonometrically convex functions in the following way. Suppose for simplicity that $\rho$ is non-integer. Then take $T_{\rho} g=g^{\prime \prime}+\rho^{2} g$ and consider

$$
\begin{aligned}
h_{1}(\phi) & :=\frac{1}{2 \rho \sin \pi \rho} \int_{0}^{2 \pi} \widetilde{\cos \rho}(\phi-\psi-\pi)\left(T_{\rho} g\right)^{+}(\phi) d \phi \\
h_{2}(\phi) & :=\frac{1}{2 \rho \sin \pi \rho} \int_{0}^{2 \pi} \widetilde{\cos \rho}(\phi-\psi-\pi)\left(T_{\rho} g\right)^{-}(\phi) d \phi .
\end{aligned}
$$

By Theorem 3.2.3.3 these functions are $\rho$ trigonometrically convex and $h_{1}-h_{2}=g$. Hence (5.7.10.11), (5.7.10.12) holds.

Proof of Theorem 5.7.10.1. We consider a function $f \in A(\rho(r))$ with the limit set $U:=\left\{v(z)=c r^{\rho} h_{1}(\phi)+(1-c) r^{\rho} h_{2}(\phi): 0 \leq c \leq 1\right\}$ with $h_{1}, h_{2}$ from the conditions of lemma, and we exploit Theorem 5.7.1.3.

Exercise 5.7.10.1 Do this in detail.

### 5.8 A generalization of the Valiron-Titchmarsh theorem

5.8.1 The point of departure on this topic is the following

Theorem VT [Va, Ti] Let $f \in A(\rho), \rho<1$ have its zeros on the negative ray. If the limit

$$
\lim _{r \rightarrow \infty} r^{-\rho} \log |f(r)|
$$

exists, then the limit

$$
\lim _{r \rightarrow \infty} r^{-\rho} n(r)
$$

exists.
The latter means that $f$ is a CRG-function.
The general problem is the following. Let $\rho$ be any non-integer number, $f \in$ $A(\rho(r))$, and suppose all zeros of $f$ lie on a finite system of rays

$$
\begin{equation*}
K_{S_{1}}:=\left\{z=r e^{i \phi}: 0<r<\infty, \phi \in S_{1}\right\} \tag{5.8.1.1}
\end{equation*}
$$

where

$$
\begin{equation*}
S_{1}:=\left\{e^{i \theta_{j}}: j=1,2, \ldots, m\right\} \tag{5.8.1.2}
\end{equation*}
$$

We write $n_{f} \in \mathcal{M}_{S_{1}}$.
Let $n_{j}$ be a zero distribution on the ray $\left\{\arg z=\theta_{j}\right\}$ and all the limits

$$
\begin{equation*}
\lim _{r \rightarrow \infty} r^{-\rho} n_{j}(r):=\Delta_{j} \tag{5.8.1.3}
\end{equation*}
$$

exist. In such a case we write $n_{f} \in \mathcal{M}_{\text {reg, } S_{1}}$.
Let $K_{S}$ be one more system of rays

$$
\begin{equation*}
S=\left\{e^{i \psi_{k}}: k=1,2, \ldots, n\right\} . \tag{5.8.1.4}
\end{equation*}
$$

Some $\psi_{k}$ can coincide with some $\theta_{j}$. Suppose that $f$ has regular growth on this system, i.e.,

$$
\begin{equation*}
h_{f}(\phi)=\underline{h}_{f}(\phi), e^{i \phi} \in S \tag{5.8.1.5}
\end{equation*}
$$

In such a case we write $f \in A_{\text {reg }, S}$.
The problem is, what is the connection between $S$ and $S_{1}$ so that the implication $\left(f \in A_{\text {reg }, S}\right) \Longrightarrow\left(n_{f} \in \mathcal{M}_{\text {reg }, S_{1}}\right)$ holds.

This problem can be reformulated in another way. For $n_{f} \in \mathcal{M}_{\text {reg, } S_{1}}$ if $n_{f} \in$ $\mathcal{M}_{S_{1}}$ it is necessary and sufficient that $f$ is a CRG-function, because existence of an angle density is equivalent to existence of all the limits. So the problem can be reformulated in the form: what is the connection between $S$ and $S_{1}$, so that the implication $\left(f \in A_{\mathrm{reg}, S}\right) \Longrightarrow(f$ is CRG-function $)$ holds.

We write

$$
\begin{equation*}
G(t, \gamma, \rho):=G_{p}\left(e^{t-i \gamma}\right) e^{-\rho t}, p=[\rho] \tag{5.8.1.6}
\end{equation*}
$$

where $G_{p}$ is the Primary Kernel:

$$
G_{p}(z)=\log |1-z|+\Re \sum_{k=1}^{p} z^{k} / k
$$

Set

$$
\hat{G}(s, \gamma, \rho):=\int_{-\infty}^{\infty} G(t, \gamma, \rho) e^{-i s t} d t
$$

This is the Fourier transformation of $G(t, \gamma, \rho)$. It can be computed (see, e.g, $[\mathrm{Oz}$, Lem. 3]);

$$
\hat{G}(s, \gamma, \rho)=\frac{\pi \cos (\pi+\gamma)(\rho+i s)}{(\rho+i s) \sin \pi(\rho+i s)}
$$

Consider the matrix

$$
\begin{equation*}
\hat{\mathbb{G}}\left(s, S_{1}-S\right):=\left\|\hat{G}\left(s, \theta_{j}-\psi_{k}, \rho\right)\right\| . \tag{5.8.1.7}
\end{equation*}
$$

We are going to prove (see $[\mathrm{Az}(1998)])$
Theorem 5.8.1.1 The implication

$$
\left\{f \in A_{\text {reg }, S}\right\} \wedge\left\{n_{f} \in \mathcal{M}_{S_{1}}\right\} \Longrightarrow\{f \text { is a } C R G \text {-function }\}
$$

holds iff

$$
\begin{equation*}
\operatorname{rank} \hat{\mathbb{G}}\left(s, S_{1}-S\right)=m, \forall s \in(-\infty, \infty) \tag{5.8.1.8}
\end{equation*}
$$

As a corollary we obtain the following ([De])
Theorem 5.8.1.2 (Delange) Suppose that $S_{1}$ and $S$ consist of one ray, i.e.,

$$
S_{1}=\left\{e^{i \theta_{1}}\right\}, S=\left\{e^{i \psi_{1}}\right\}
$$

The implication (5.8.1.5) holds iff

$$
\begin{equation*}
\theta_{1}-\psi_{1} \neq(1-(2 k+1) / 2 \rho) \pi, \quad k=1,2, \ldots \tag{5.8.1.9}
\end{equation*}
$$

5.8.2 A Fourier transformation for distribution $\nu$ on the real axes is a distribution in the standard space $\mathcal{S}^{\prime}$ (see [Hö, vol. 1, Ch. 7, §7.1]). For a locally bounded measure whose variation is "not very quickly" growing, it can be defined by

$$
(\mathcal{F} \nu)(s):=\lim _{\epsilon \rightarrow 0} \int_{-\infty}^{\infty} e^{i t s} e^{-\frac{\epsilon t^{2}}{2}} \nu(d t)
$$

where the right side is understood in the sense of distributions.
For example, if $\nu(d t):=e^{i s_{0} t} d t$, we have $\mathcal{F} \nu(s)=\delta\left(s-s_{0}\right)$ where $\delta$ is the Dirac function.

Exercise 5.8.2.1 Check this.
For distribution and a summable function one can define a convolution for which the property $\mathcal{F}(f * \nu)(s)=\mathcal{F} f(s) \mathcal{F} \nu(s)$ holds.

Proof of Theorem 5.8.1.1. Since $f \in \mathcal{M}_{S_{1}}$ the limit set $\operatorname{Fr} n_{f}$ is concentrated on $K_{S_{1}}$. So every $v \in \mathbf{F r}[f]$ can be represented in the form (see Theorem 3.1.5.1):

$$
\begin{equation*}
v(z)=\sum_{j=1}^{j=m} \int_{0}^{\infty} G_{p}\left(z / r e^{i \theta_{j}}\right) \mu_{j}(d r) \tag{5.8.2.1}
\end{equation*}
$$

where $\mu_{j}$ is concentrated on the ray $\left\{\arg \zeta=\theta_{j}\right\}$ and belongs to $U[\rho]$. After changing variables,

$$
r=e^{\tau},|z|=e^{t}
$$

we obtain from (5.8.2.1)

$$
\begin{equation*}
v^{1}\left(t e^{i \phi}\right)=\sum_{j=1}^{j=m} \int_{-\infty}^{\infty} G\left(t-\tau, \phi-\theta_{j}, \rho\right) \mu_{j}^{1}(d \tau) \tag{5.8.2.1a}
\end{equation*}
$$

where

$$
\begin{equation*}
e^{\rho \tau} \mu_{j}^{1}(d \tau):=\mu_{j}(d r), v^{1}\left(t e^{i \phi}\right):=v\left(|z| e^{i \phi}\right) e^{-\rho|z|} . \tag{5.8.2.2}
\end{equation*}
$$

The equality (5.8.2.1a) can be written as

$$
\begin{equation*}
v^{1}\left(t e^{i \phi}\right)=\sum_{j=1}^{j=m}\left[G\left(\bullet, \phi-\theta_{j}, \rho\right) * \mu_{j}^{1}\right](t) \tag{5.8.2.3}
\end{equation*}
$$

where * stands for convolution. Then $f \in A_{\text {reg,S }}$ with $n_{f} \in \mathcal{M}_{S_{1}}$, iff every pair $v_{1}, v_{2} \in \mathbf{F r}[f]$ satisfies the condition

$$
\begin{equation*}
v_{1}(z)=v_{2}(z), z \in K_{S} \tag{5.8.2.4}
\end{equation*}
$$

Denote by $\mu_{1, j}, \mu_{2, j}$ the restriction of $\mu_{v_{1}} \mu_{v_{1}}$ to the ray $\left\{\arg z=\theta_{j}\right\}$. Set $\nu_{j}:=$ $\mu_{1, j}-\mu_{2, j}$ Using (5.8.2.3) we can rewrite (5.8.2.4) in the form

$$
\begin{equation*}
\sum_{j=1}^{j=m}\left[G\left(\bullet, \phi_{k}-\theta_{j}, \rho\right) * \nu_{j}^{1}\right](t) \equiv 0, k=1,2, \ldots, n \tag{5.8.2.5}
\end{equation*}
$$

Applying Fourier transforms we obtain a system of linear equations:

$$
\begin{equation*}
\sum_{j=1}^{j=m}\left[\hat{G}\left(\bullet, \psi_{k}-\theta_{j}, \rho\right) \cdot \hat{\nu}_{j}^{1}\right](t) \equiv 0, k=1,2, \ldots, n \tag{5.8.2.6}
\end{equation*}
$$

Suppose now that $\operatorname{rank} \hat{\mathbb{G}}\left(s, S-S_{1}\right)=m$ for every $s \in \mathbb{R}$. The system (5.8.2.6) has only the trivial solution for every $s$. Thus $\hat{\nu}_{j}^{1}(s) \equiv 0$, for $j=1,2, \ldots, m$. This implies $\nu_{j}^{1}(t) \equiv 0$ for $j=1,2, \ldots, m$ and $\nu_{j} \equiv 0$ for $j=1,2, \ldots, m$. Thus $\mu_{v_{1}}=\mu_{v_{2}}$, i.e., (by (5.8.2.3)) $\mathbf{F r}[f]$ consists of one function $v \in U[\rho]$. Thus $f$ is a CRG-function.

Conversely, suppose that $\operatorname{rank} \hat{\mathbb{G}}\left(s, S-S_{1}\right)<m$ for some $s_{0}$.
Then there exists a nontrivial solution $\left(b_{1}, \ldots, b_{m}\right)$ that satisfies the corresponding system. We obtain that $\left\{\hat{\nu}_{j}^{1} b_{j} \delta\left(s-s_{0}\right), j=1,2 \ldots, m\right\}$ is a solution of (5.8.2.6) for all $s \in \mathbb{R}$ and hence

$$
\nu_{j}^{1}(d t)=b_{j} e^{i t s_{0}} d t, j=1,2, \ldots, m
$$

Since $\nu_{j}^{1}$ have bounded densities $d \nu_{j}^{1} / d t$, we can find a constant $C$ such that $\sup \left\{\left|d \nu_{j}^{1} / d t\right|: 0<t<\infty, j=1,2 \ldots, m\right\} \leq C$.

Set

$$
\begin{equation*}
\mu_{1, j}^{1}(d t)=C d t+\nu_{j}^{1}(d t) ; \mu_{2, j}^{1}=C d t . \tag{5.8.2.7}
\end{equation*}
$$

Both of these are measures. Now we pass to $m_{1, j}, m_{2, j}$ via (5.8.2.2). It is easy to check that $m_{1, j}, m_{2, j} \in \mathcal{M}(\rho)$.

Exercise 5.8.2.2 Check this.
Consider $\mu_{1}, \mu_{2} \in \mathcal{M}(\rho)$ which are defined uniquely by their restrictions $\mu_{1, j}, \mu_{2, j}$ respectively on $K_{S_{1}}$. Set

$$
v_{1}(z):=\int_{\mathbb{C}} G_{p}(z / \zeta) \mu_{1}(d \xi d \eta) ; v_{2}(z):=\int_{\mathbb{C}} G_{p}(z / \zeta) \mu_{2}(d \xi d \eta) ; \zeta=\xi+i \eta
$$

It is easy to check that the equality

$$
\begin{equation*}
v_{1}(z)=v_{2}(z), z \in K_{S} \tag{5.8.2.8}
\end{equation*}
$$

holds.
Exercise 5.8.2.3 Check this.
Since $\mu_{1}$ and $\mu_{2}$ are finite sums of trigonometrical functions, for $v_{1}$ and $v_{2}$ the condition (4.1.3.3) is satisfied. Thus by Theorem 4.3.6.1 there exists a function $f \in A(\rho(r))$ for which

$$
\operatorname{Fr}[f]=\bigcup_{0 \leq c \leq 1} \mathbb{C}\left(c v_{1}+(1-c) v_{2}\right) .
$$

Since for $v \in \mathbb{C}\left(c v_{1}+(1-c) v_{2}\right)$ (5.8.2.8) also holds, the same holds for $v \in \operatorname{Fr}[f]$ and this function is not a CRG-function.

## Chapter 6

## Application to the Completeness of Exponential Systems in Convex Domains and the Multiplicator Problem

The completeness of exponential systems in convex domains is intimately connected to the multiplicator problem. Considering a special form of exponent system is related to the study of special subharmonic functions that determine the periodic limit set, the so-called automorphic subharmonic functions. The next Sections 6.1, 6.2 are devoted to these problems.

### 6.1 The multiplicator problem

6.1.1 Let $\Phi \in A(\rho(r))$ and let $H(\phi)$ be a $\rho$-trigonometrically convex function. A function $g \in A(\rho(r))$ is called an $H$-multiplicator of $\Phi$ if the indicator $h_{g \Phi}$ of the product $g \Phi$ satisfies the inequality

$$
h_{g \Phi}(\phi) \leq H(\phi), \forall \phi .
$$

In some questions (see Section 6.3) we need to determine whether a given function $\Phi$ has a multiplicator. We shall study this problem in terms of the limit set of $\Phi$. Define $H(z):=r^{\rho} H(\phi), z=r e^{i \phi}$. Let $v \in U[\rho]$ (see (3.1.2.4)). Consider the function

$$
m(z, v, H):=H(z)-v(z) .
$$

As will be proved in Corollary 6.1.9.3, the maximal subharmonic minorant of $m(z, v, H)$ exists and is continuous. The maximal subharmonic minorant of $m$
(m.s.m.) belonging to $U[\rho]$ will be denoted by $\mathcal{G}_{H} v$, while the domain of definition of the operator $\mathcal{G}_{H}$ will be denoted by $D_{H}$. Though $m(0, \bullet, \bullet)=0$, the m.s.m. of $m$ can differ from zero (as was remarked by A.E. Eremenko and M.L. Sodin), but if the m.s.m. of $m$ equals zero at zero, then it belongs to $U[\rho]$.

Exercise 6.1.1.0 Prove this.
Exercise 6.1.1.1 Consider the function

$$
w(z)= \begin{cases}|z| \log |z|, & \text { if }|z| \leq 1 \\ |z|-1, & \text { if }|z| \geq 1\end{cases}
$$

It is subharmonic and belongs to $U[1]$. Show that the maximal subharmonic minorant of $K|z|-w(z)$ is different from zero in 0 for every $K>0$.

Theorem 6.1.1.1 ([AG(1992)]) $\Phi \in A(\rho(r))$ has an $H$-multiplicator iff

$$
\begin{equation*}
\operatorname{Fr}[\Phi] \subset D_{H} \tag{6.1.1.1}
\end{equation*}
$$

Proof of necessity. Let $g$ be a multiplicator of $\Phi$, i.e.,

$$
\begin{equation*}
h_{g \Phi}(\phi) \leq H(\phi) \tag{6.1.1.2}
\end{equation*}
$$

and let $v \in \operatorname{Fr}[\Phi]$. We can choose $v_{g \Phi} \in \operatorname{Fr}[g \Phi]$ and $v_{g} \in \operatorname{Fr}[g]$ such that $v_{g \Phi}=$ $v+v_{g}$ (see Theorem 3.1.2.4, fru1)).

Exercise 6.1.1.2 Prove this directly.
By definition of indicator (3.2.1.1) and (6.1.1.2) we have $v_{g \Phi}(z) \leq H(z)$ or $v_{g}(z) \leq m(z, v, H)$. Since $v_{g} \in U[\rho], v \in D_{H}$.

For proving sufficiency we need the following
Theorem 6.1.1.2 The operator $\mathcal{G}_{H}$ is

1. upper semicontinuous in the $\mathcal{D}^{\prime}$-topology, 6.1.1.5, i.e.,

$$
\left(v_{j} \rightarrow v\right) \wedge\left(\mathcal{G}_{H} v_{j} \rightarrow w\right) \Longrightarrow(w \in U[\rho]) \wedge\left(w(z) \leq \mathcal{G}_{H}(z), z \in \mathbb{C}\right)
$$

2. invariant: $\left(\mathcal{G}_{H} v\right)_{[t]}=\mathcal{G}_{H} v_{[t]} ;\left(\right.$ see (3.1.2.4a) for $\left.P_{t} \equiv t I\right)$;
3. concave:

$$
\left(\forall v_{1}, v_{2} \in D_{H}, c \in[0 ; 1]\right) \Longrightarrow\left(v_{c}:=c v_{1}+(1-c) v_{2} \in D_{H}\right)
$$

and

$$
\mathcal{G}_{H}\left(v_{c}\right) \geq c \mathcal{G}_{H}\left(v_{1}\right)+(1-c) \mathcal{G}_{H}\left(v_{2}\right)
$$

Proof. Let us prove 1). Suppose $v_{j} \in U[\rho] \rightarrow v$ and $\mathcal{G}_{H} v_{j} \rightarrow w$. Then

$$
\begin{equation*}
\mathcal{G}_{H} v_{j} \leq H(z)-v_{j}(z), z \in \mathbb{C} . \tag{6.1.1.3}
\end{equation*}
$$

Applying $(\bullet)_{\epsilon}$ from (2.6.2.2) and Theorem 2.3.4.5, reg 3), we obtain

$$
w_{\epsilon} \leq(H)_{\epsilon}(z)-(v)_{\epsilon}(z), z \in \mathbb{C} w_{\epsilon}(0) \geq 0
$$

Passing to the limit as $\epsilon \downarrow 0$ we obtain by Theorem 2.6.2.3, ap2)

$$
w(z) \leq H(z)-v(z)=m(z, H, v), z \in \mathbb{C} .
$$

Since $0 \leq w(0) \leq m(0, H, v)=0$ we have $w(0)=0$ and, hence, $w \in U[\rho]$. Thus $v \in D_{H}$ and $w(z) \leq \mathcal{G}_{H} v(z)$.

Let us prove 2). Since $H(z)$ is invariant with respect to $(\bullet)_{[t]}$,

$$
(\mathcal{G} v)_{[t]} \leq H(z)-v_{[t]} .
$$

Hence,

$$
\begin{equation*}
(\mathcal{G} v)_{[t]}(z) \leq\left(\mathcal{G}\left(v_{[t]}\right)\right)(z), \tag{6.1.1.4}
\end{equation*}
$$

because $\mathcal{G}\left(v_{[t]}\right)$ is the maximal subharmonic minorant. We can replace $v$ with $v_{[1 / t]}$ and obtain $\left(\mathcal{G} v_{[1 / t]}\right)_{[t]}(z) \leq \mathcal{G} v(z)$. Applying $(\bullet)_{[t]}$ to the two sides of the inequality, we obtain $\mathcal{G} v_{[1 / t]}(z) \leq(\mathcal{G} v(z))_{[1 / t]}$. Now we can replace $1 / t$ with $t$ and obtain the reverse inequality to (6.1.1.4), which, together with (6.1.1.4), proves 2 ).
$3)$. Let $v_{1}, v_{2} \in D_{H}$ and $c \in[0 ; 1]$. One has

$$
\mathcal{G} v_{i}(z) \leq H(z)-v_{i}(z), i=1,2, \quad \forall z .
$$

Then

$$
\left[c \mathcal{G} v_{1}+(1-c) \mathcal{G} v_{2}\right](z) \leq H(z)-\left[c v_{1}+(1-c) v_{2}\right](z) .
$$

Thus

$$
\left[c \mathcal{G} v_{1}+(1-c) \mathcal{G} v_{2}\right](z) \leq \mathcal{G}\left[c v_{1}+(1-c) v_{2}\right](z)
$$

Proof of sufficiency in Theorem 6.1.1.1. Assume that $\operatorname{Fr}[\Phi] \subset D_{H}$ and consider the set

$$
\begin{equation*}
\boldsymbol{U}:=\left\{\left(v^{\prime}, v^{\prime \prime}\right): v^{\prime \prime} \leq \mathcal{G} v^{\prime}, v^{\prime} \in \mathbf{F r}[\Phi]\right\} \tag{6.1.1.5}
\end{equation*}
$$

Then $\boldsymbol{U}$ is nonempty, because of (6.1.1.1), closed, because of Theorem 6.1.1.2, 1), and invariant, because of Theorem 6.1.1.2, 2).

Every fiber $\boldsymbol{U}^{\prime \prime}=\left\{v^{\prime \prime}: v^{\prime \prime} \leq \mathcal{G} v^{\prime}\right\}$ is convex because of Theorem 6.1.1.2, 3). By Theorem 4.4.1.2 there exists $u^{\prime \prime} \in U(\rho(r))$ such that for the curve $\boldsymbol{u}:=\left(u^{\prime}, u^{\prime \prime}\right)$,

$$
\begin{equation*}
\operatorname{Fr}[\boldsymbol{u}]=\boldsymbol{U} \tag{6.1.1.6}
\end{equation*}
$$

By Theorem 5.3.1.4 (Approximation Theorem) the function $u^{\prime \prime}$ can be replaced with $\log |g|$, where $g \in A(\rho(r))$, retaining the property (6.1.1.6).

Let us prove that $g$ is an $H$-multiplicator of $\Phi$. Indeed, set $\Pi:=g \Phi$. It is enough to prove that for every $v_{\Pi} \in \mathbf{F r}[\Pi]$,

$$
\begin{equation*}
v_{\Pi}(z) \leq H(z) \tag{6.1.1.7}
\end{equation*}
$$

Note that every $v_{\Pi}$ has the form $v_{\Pi}=v_{g}+v$, where $\left(v, v_{g}\right) \in \boldsymbol{U}$. Thus, because of definitions (6.1.1.5) and (6.1.1.6), $v_{\Pi}$ satisfies (6.1.1.7).

Let us note that the pair $\left(v, \mathcal{G}_{H} v\right) \in \boldsymbol{U}$ because of closeness of $\boldsymbol{U}$. Hence the following assertion holds.

Proposition 6.1.1.3 Every $\Phi \in A(\rho)$ that satisfies (6.1.1.1) has a multiplicator $g \in A(\rho)$ such that

$$
\begin{equation*}
v+\mathcal{G}_{H} v \in \mathbf{F r}[g \Phi] . \tag{6.1.1.8}
\end{equation*}
$$

Exercise 6.1.1.3 Check this in detail.
Although $v \in U[\rho]$ is in general an upper semicontinuous function, we need
Theorem 6.1.1.4 The function $\mathcal{G}_{H} v(z), v \in U[\rho]$, is a continuous function that is harmonic outside the set $E=\left\{z: \mathcal{G}_{H} v(z)=m(z, v, H)\right\}$.

Proof. $\mathcal{G}_{G} v(z)$ is continuous because of Corollary 6.1.9.3. If $\mathcal{G}_{H} v\left(z_{0}\right)<v\left(z_{0}\right)$ and if $\mathcal{G}_{H} v(z)$ is not harmonic in a neighborhood of $z_{0}$, we can make sweeping of masses from a small disc $\left\{\left|z-z_{0}\right|<\epsilon\right\}$ (see Theorem 2.7.2.1). The obtained subharmonic function will be greater than $\mathcal{G}_{H} v(z)$, contradicting maximality.
6.1.2 Suppose that some $H$-multiplicator $g=g(z, \Phi, H)$ of the function $\Phi$ is found. We examine the function $\Pi=g \Phi$. The structure of its limit set is described by the following statement:

Proposition 6.1.2.1 Every $v_{\Pi} \in \operatorname{Fr}[g \Phi]$ can be written as $v_{\Pi}=v+w_{1}$, where $v \in \operatorname{Fr}[\Phi]$ and $w_{1} \in U[\rho]$ with the condition

$$
\begin{equation*}
w_{1}(z) \leq \mathcal{G}_{H}(z), \forall z \in \mathbb{C} \tag{6.1.2.1}
\end{equation*}
$$

and, conversely, for every $v \in \mathbf{F r}[\Phi]$ there exists $a v_{g}, v_{g}(z) \leq \mathcal{G}_{H} v(z)$, such that

$$
v+v_{g} \in \operatorname{Fr}[g \Phi] .
$$

Exercise 6.1.2.1 Prove this the same way as in Exercise 6.1.1.2.
An $H$-multiplicator $G$ of the function $\Phi$ will be called ideally complementing if it satisfies the condition

$$
\operatorname{Fr}[G \Phi]=\left\{v_{\Pi}=v+\mathcal{G}_{H} v: v \in \operatorname{Fr}[\Phi]\right\} .
$$

If a multiplicator is ideally complementing, then equality is achieved in (6.1.2.1) for all $v \in \operatorname{Fr}[\Phi]$. This make the multiplicator optimal in another respect. Recall that an entire function $f$ is of minimal type with respect to a proximate order $\rho(r), \rho(r) \rightarrow \rho$ if (see (2.8.1.6))

$$
\sigma_{f}:=\limsup _{r \rightarrow \infty} \log M(r, f) r^{-\rho(r)}=0
$$

Proposition 6.1.2.2 Let $G=G(\bullet, \Phi, H)$ be an ideally complementing $H$-multiplicator of a function $\Phi$. Then each $H$-multiplicator of the function $\Pi=G \Phi$ is of minimal type.

This proposition is proved in Section 6.1.3.
A function $\Phi$ is said to be ideally complementable if for each $H$ the condition (6.1.1.1) implies that $\Phi$ has an ideally complementing multiplicator. For instance, if $\Phi$ is a function of completely regular growth (see Section 5.6) then it is ideally complementable.

Exercise 6.1.2.2 Prove this.

Theorem 6.1.2.3 Every function with periodic limit set is ideally complementable.
This theorem is proved in Section 6.1.6.
Let $C \subset \mathbb{R}^{l}$ be an $l$-dimensional connected compact and let $\{h(\phi, c): c \in C\}$ be a set of $\rho$-t.c. functions that is continuous with respect to $c \in C$. For example, $c \in[0,1]$ and $h(\phi, c)=c h_{1}(\phi)+(1-c) h_{2}(\phi)$. The set

$$
\begin{equation*}
U_{\mathrm{ind}}:=\left\{v\left(r e^{i \phi} s\right)=r^{\rho} h(\phi, c): c \in C\right\} \tag{6.1.2.2}
\end{equation*}
$$

is the limit set of an entire function.
Exercise 6.1.2.3 Prove this using Theorem 4.3.6.1.
Such a set is called a set of indicators. Entire functions with such limit sets can be also considered as a generalization of CRG-functions.

Theorem 6.1.2.4 Every function $\Phi$ whose limit set is a set of indicators is ideally complementable.

This theorem is proved in Section 6.1.7.
The existence of an ideally complementing $H$-multiplicator depends, of course, both on $\Phi \in A(\rho(r))$ (or, more precisely, on its limit set $\operatorname{Fr}[\Phi]$ ) and on $H$.

Theorem 6.1.2.5 Let $\Phi$ and $H$ be such that the condition (6.1.1.1) is satisfied. The function $\Phi$ has an ideally complementing $H$-multiplicator if and only if the operator $\mathcal{G}_{H}$ is continuous on $\operatorname{Fr} \Phi$.

This theorem is proved in Section 6.1.6.

Now we formulate a sufficient condition for continuity of the operator $\mathcal{G}_{H}$. We shall say that the maximum principle for $U[\rho]$ is valid in the domain $G$, (which is, generally speaking, unbounded), if the conditions $w \in U[\rho], w(z)=0$ for $z \notin G$ imply $w(z) \equiv 0$.

Let us denote by $\mathcal{H}_{w}$ a region of harmonicity of $w \in U[\rho]$, i.e., a region where the conditions " $w$ is harmonic in $G$ " and " $G \supset \mathcal{H}_{w}$ " imply $G=\mathcal{H}_{w}$.

We remark that $\mathcal{H}_{w}$ is a connected component of the open set on which $w$ is harmonic. Generally it is not unique.

The image of $U \in U[\rho]$ will be denoted by $\mathcal{G}_{H} U$, while its closure in the $\mathcal{D}^{\prime}$-topology will be denoted by $\operatorname{clos} \mathcal{G}_{H} U$.

Theorem 6.1.2.6 Suppose for every $w \in \operatorname{clos} \mathcal{G}_{H} U$ and every $\mathcal{H}_{w}$ the maximum principle for $U[\rho]$ holds. Then $\mathcal{G}_{H}$ is continuous on $U$.

This theorem is proved in Section 6.1.5.
In Section 6.1 .8 we will construct an example of $\Phi$ and $H$ such that the operator $\mathcal{G}_{H}$ is not continuous on $\operatorname{Fr}[\Phi]$. This is also an example of the function that has no ideally complementing multiplicator.

### 6.1.3

Proof of Proposition 6.1.2.2. Let $g$ be an ideally complementing multiplicator of the function $\Pi=G \Phi$. We write

$$
\begin{equation*}
\theta(z):=(g G \Phi)(z) \tag{6.1.3.1}
\end{equation*}
$$

Let $v_{g} \in \mathbf{F r}[g]$. Let us choose $t_{j} \rightarrow \infty$ such that:

$$
(\log |g|)_{t_{j}} \rightarrow v_{g} ; \quad(\log |\Pi|)_{t_{j}} \rightarrow v_{\Pi} \in \operatorname{Fr}[\Pi] ; \quad(\log |\theta|)_{t_{j}} \rightarrow v_{\theta} \in \operatorname{Fr}[\theta] .
$$

It follows from (6.1.3.1) that $v_{\theta}=v_{g}+v_{\Pi}$. Since $g$ is a multiplicator of $\Pi$, we have

$$
\begin{equation*}
v_{\theta}(z)=v_{g}(z)+v_{\Pi}(z) \leq H(z) \tag{6.1.3.2}
\end{equation*}
$$

Since $G$ is an ideally complementing multiplicator, $v_{\Pi}=v+\mathcal{G}_{H} v$. So for all $z \in \mathbb{C}$ (6.1.3.2) implies

$$
\left(v_{g}+\mathcal{G}_{H} v\right)(z) \leq(H-v)(z) .
$$

Since $\mathcal{G}_{H} v$ is the maximal subharmonic minorant, $v_{g}(z) \leq 0$ and hence $v_{g}(z) \equiv 0$. Thus (see (3.2.1.1)) we have $h_{g}(\phi) \equiv 0$ and therefore

$$
\sigma_{g}=\max _{0 \leq \phi \leq 2 \pi} h_{g}(\phi)=0
$$

6.1.4 In order to prove Theorem 6.1.2.6 we need a number of auxiliary statements.

Lemma 6.1.4.1 Let the maximum principle be valid in $G$ for $U[\rho]$ and for some continuous functions $w_{1}, w \in U[\rho]$ satisfy:
a) $w$ is harmonic in $G$;
b) $w_{1}(z)=w(z)$ outside of $G$.

Then

$$
\begin{equation*}
w_{1}(z) \leq w(z), \quad z \in G \tag{6.1.4.1}
\end{equation*}
$$

Proof. We set

$$
w_{0}(z):= \begin{cases}\left(w_{1}-w\right)^{+}(z), & z \in G \\ 0, & z \notin G\end{cases}
$$

This function is continuous in $\mathbb{C}$ and, evidently, subharmonic both in $G$ and in $\mathbb{C} \backslash \bar{G}$. Since $w_{0}(z) \geq 0$, the inequality for the mean values

$$
0=w_{0}(z) \leq \frac{1}{2 \pi} \int_{0}^{2 \pi} w_{0}\left(z+\epsilon r e^{i \phi}\right) d \phi, z \in \partial G
$$

implies the subharmonicity on $\partial G$. Since $G$ satisfies the maximum principle for $U[\rho]$, we have $w_{0} \equiv 0$, which is equivalent to (6.1.4.1).

Now we shall dwell on some properties of maximal subharmonic minorants and, in particular, of $w=\mathcal{G}_{H} v$. Let

$$
\begin{equation*}
E_{v}:=\left\{z \in \mathbb{C}: \mathcal{G}_{H} v(z)=m(z, v, H)\right\} . \tag{6.1.4.2}
\end{equation*}
$$

We remark that $m(z, v, H)$ is a $\delta$-subharmonic function in $\mathbb{C}$ whose charge will be denoted by $\nu(\bullet, v)$, its positive and negative parts will be denoted by $\nu^{+}$and $\nu^{-}$.

Let us denote by $\mu_{H}$ the measure of $H(z)$. It is decomposed into the product of measures (see Section 3.2 and Proposition 5.6.3.1)

$$
\begin{equation*}
\mu_{H}=\Delta_{H} \otimes \rho r^{\rho-1} d r \tag{6.1.4.3}
\end{equation*}
$$

where $\Delta_{H}$ is the measure on the unit circle and $\rho r^{\rho-1} d r$ is the measure on the ray. It is obvious that

$$
\begin{equation*}
\nu^{+}(\bullet, v) \leq \mu_{H}(\bullet) . \tag{6.1.4.4}
\end{equation*}
$$

We shall denote the mass distribution of $w \in U[\rho]$ by $\mu_{w}$.
The modulus of continuity of $w$ (if $w$ is continuous) will be denoted by $\omega_{w}(z, h), z \in \mathbb{C}, h>0$.

The following lemma lists various properties of $w \in \mathcal{G}_{H} U, U \subset U[\rho]$ which will be useful in the sequel:

Lemma 6.1.4.2 Let $w \in \mathcal{G}_{H} U$. Then

1. $w \in U[\rho, \sigma]$ where

$$
\begin{aligned}
\sigma & =4 \cdot 2^{\rho}\left[\max \left\{H\left(e^{i \phi}\right): \phi \in[0,2 \pi]\right\}+2 \sigma_{1}\right], \\
\sigma_{1} & =\max \left\{v(z)|z|^{-\rho}: z \in \mathbb{C}, v \in U\right\}
\end{aligned}
$$

2. the charge restriction $\left.\nu(\bullet, v)\right|_{E_{v}}$ to $E_{v}$ is nonnegative, i.e.,

$$
\left.\nu(\bullet, v)\right|_{E_{v}}=\left.\nu^{+}(\bullet, v)\right|_{E_{v}}
$$

3. outside $E_{v}$ the function $w$ is harmonic, i.e.,

$$
\left.\mu_{w}\right|_{\mathbb{C} \backslash E_{v}}=0
$$

4. the measure $\mu_{w}$ is bounded from above by $\nu^{+}(\bullet, v)$, i.e.,

$$
\mu_{w} \leq \nu^{+}(\bullet, v)
$$

5. $\mathcal{G}_{H} U$ is equicontinuous on each compact set, i.e.,

$$
\omega_{w}(z, h) \leq C(R, \sigma, \rho) \sqrt{h} \log (1 / h),|z| \leq R
$$

where $C(R, \sigma, \rho)$ is independent of $w \in \mathcal{G}_{H} U$.
Proof. Let us prove property 1. We have

$$
\begin{aligned}
T(r, w) & :=\frac{1}{2 \pi} \int_{0}^{2 \pi} w^{+}\left(r e^{i \phi}\right) d \phi \\
& \leq \frac{1}{2 \pi}\left[r^{\rho} \int_{0}^{2 \pi} H^{+}\left(e^{i \phi}\right) d \phi+\int_{0}^{2 \pi} v^{+}\left(r e^{i \phi}\right) d \phi+\int_{0}^{2 \pi} v^{-}\left(r e^{i \phi}\right) d \phi\right]
\end{aligned}
$$

Since $v(0)=0$, we have

$$
\int_{0}^{2 \pi} v^{-}\left(r e^{i \phi}\right) d \phi \leq \int_{0}^{2 \pi} v^{+}\left(r e^{i \phi}\right) d \phi
$$

Therefore

$$
T(r, w) \leq\left[\max \left\{H\left(e^{i \phi}\right): \phi \in[0,2 \phi]\right\}+2 \sigma_{1}\right] r^{\rho}
$$

It is known (see Theorem 2.8.2.3, (2.8.2.5)) that $M(r) \leq 4 T(2 r)$. So we conclude that

$$
w(z) \leq 4 \cdot 2^{\rho}\left[\max \left\{H\left(e^{i \phi} ; \phi \in[0,2 \pi]\right\}+2 \sigma_{1}\right]|z|^{\rho}=\sigma|z|^{\rho} .\right.
$$

Let us prove property 2. To this end we shall use the following theorem (Grishin's Lemma) [Gr].

Theorem A.F.G Let $g$ be a nonnegative $\delta$-subharmonic function, and let $\nu_{g}$ be its charge. Then the restriction $\left.\nu_{g}\right|_{E}$ to the set $E=\{z: g(z)=0\}$ is a measure.

Applying this theorem to the function $g:=m(z, v, H)-\mathcal{G}_{H} v(z)$, we get

$$
\begin{equation*}
\left.\nu(\bullet, v)\right|_{E_{v}} \geq\left.\mu_{w}\right|_{E_{v}}, \tag{6.1.4.5}
\end{equation*}
$$

hence, we obtain property 2.

Let us prove property 3 . Since $w$ and $v$ are upper semicontinuous, and $H$ is continuous (see Theorem 3.2.5.5), the set $\{z:(w+v)(z)-H(z)<0\}$ is open.

Let us take a neighborhood of an arbitrary point of this set and replace the function $w$ within it with the Poisson integral constructed using this function, i.e., let us sweep out the mass from this neighborhood. The subharmonic function obtained would be strictly greater than the initial one, if the latter were not harmonic. This means that the initial $w$ was not the maximal minorant. We have arrived at a contradiction, which proves property 3 .

Property 4 immediately follows from property 3 and (6.1.4.5).
In order to prove property 5 we shall need an auxiliary statement which will be stated as a number of lemmas. Let

$$
P(z, \phi, R):=\frac{1}{2 \pi} \frac{R^{2}-|z|^{2}}{\left|z-R e^{i \phi}\right|}
$$

be the Poisson kernel in the disc $K_{R}=\{z:|z|<R\}$.
Below, $C^{\prime}$ s with indices will denote constants.
Lemma 6.1.4.3 In the disc $K_{R / 2}$, we have

$$
\left|\operatorname{grad}_{z} P(z, \phi, R)\right| \leq C_{1}(R)
$$

where $C_{1}(R)$ depends only on $R$.

Exercise 6.1.4.1 Prove this.
We shall introduce the notation for the Green function for the Laplace operator in the disc $K_{R}$ :

$$
G(z, \zeta, R):=\log \left|\frac{R^{2}-\zeta \bar{z}}{R(z-\zeta)}\right|
$$

The disc $\{\zeta:|\zeta-z|<t\}$ will be denoted by $K_{z, t}$.
Lemma 6.1.4.4 Let $z \in K_{R / 2} \backslash K_{\zeta, \sqrt{h}}$. Then for a small $h$,

$$
\left|\operatorname{grad}_{z} G(z, \zeta, r)\right| \leq C_{2}(R) / \sqrt{h}
$$

Exercise 6.1.4.2 Prove this.
Let us write $\mu(z, t):=\mu\left(K_{z, t}\right)$.
Lemma 6.1.4.5 For $z \in K_{R / 2}, 0<t<R / 10$, we have

$$
\mu_{H}(z, t) \leq C_{3}(\sigma, R) t
$$

Proof. Applying Theorem 2.6.5.1 (Jensen-Privalov) to the function $H(z)$, we obtain

$$
M_{H}=\max \left\{H\left(e^{i \phi}: \phi \in[0 ; 2 \pi]\right\}=\Delta_{H}(\mathbb{T}) / \rho\right.
$$

where $\mathbb{T}$ is the unit circle.
Now

$$
\mu(z, t) \leq \Delta_{H}(\mathbb{T}) \int_{|z|-t}^{|z|+t} r^{\rho-1} d r \leq \rho^{2} M_{H} R^{\rho-1} t \leq \sigma C(\rho) R^{\rho-1} t
$$

where $C(\rho)$ is a constant depending only on $\rho$. This proves the lemma.
Lemma 6.1.4.6 Let $h<1$ and suppose that a monotonic function $\mu(t)$ satisfies the condition

$$
\begin{equation*}
\mu(t)<c t . \tag{6.1.4.6}
\end{equation*}
$$

Then

$$
\int_{0}^{\sqrt{h}} \log (1 / t) \mu(d t) \leq(3 / 2) c \sqrt{h} \log h .
$$

Exercise 6.1.4.3 Prove this by integrating by parts and using (6.1.4.6).
Lemma 6.1.4.7 Let $z \in K_{R / 2}$ and $\zeta \in K_{R}$. Then

$$
|\log |\left(R^{2}-\zeta \bar{z} / R \| \leq C_{4}(R)\right.
$$

Exercise 6.1.4.4 Prove this.
Now we pass to the proof of assertion 5 from Lemma 6.1.4.2. According to the F. Riesz theorem (Theorem 2.6.4.3) we represent $w$ in the circle as

$$
\begin{equation*}
w(z)=H(z, w)-\int_{K_{R}} G(z, \zeta, R) \mu_{w}(d \xi d \eta), \zeta=\xi+i \eta \tag{6.1.4.7}
\end{equation*}
$$

where

$$
H(z, w)=\frac{1}{2 \pi} \int_{0}^{2 \pi} P(z, \phi, R) w\left(R e^{i \phi}\right) d \phi
$$

It follows from Lemma 6.1.4.3 and 1 of Lemma 6.1.4.2 that

$$
\begin{equation*}
\left|\operatorname{grad}_{z} H(z, w)\right| \leq C_{1}(R) \frac{1}{2 \pi} \int_{0}^{2 \pi}|w|\left(\operatorname{Re}^{i \phi}\right) d \phi \leq C_{1}(R) 2 \sigma R^{\rho} \tag{6.1.4.8}
\end{equation*}
$$

We split the integral in (6.1.4.7) into three terms:

$$
\begin{aligned}
& \psi_{1}(z, h):=\int_{K_{r} \backslash K_{z_{0}}, \sqrt{h}} G(z, \zeta, R) \mu_{w}(d \xi d \eta) \\
& \psi_{2}(z, h)=\int_{K_{z_{0}}, \sqrt{h}} \log \left|\left(R^{2}-\bar{z} \zeta\right) / R\right| \mu(d \xi d \eta) \\
& \psi_{3}(z, h)=\int_{K_{z_{0}}, \sqrt{h}} \log |\zeta-z| \mu(d \xi d \eta)
\end{aligned}
$$

where $z_{0}$ is an arbitrary fixed point in $K_{R / 2}$.
Combining property 4 and inequality (6.1.4.4) we have

$$
\begin{equation*}
\mu_{w}(E) \leq \mu_{H}(E), \forall E \subset K_{R} \tag{6.1.4.9}
\end{equation*}
$$

For all $z \in K_{z_{0}, \sqrt{h} / 2}$ Lemma 6.1.4.4 yields

$$
\begin{equation*}
\left|\operatorname{grad} \psi_{1}(z, h)\right| \leq C_{2}(r) \sigma R^{\rho} / \sqrt{h} . \tag{6.1.4.10}
\end{equation*}
$$

Combining Lemmas 6.1.4.5 and 6.1.4.7 with inequality (6.1.4.9), we get

$$
\begin{equation*}
\left|\psi_{2}(z, h)\right| \leq C_{4}(R) C_{3}(\sigma, R) \sqrt{h} . \tag{6.1.4.11}
\end{equation*}
$$

Further, from Lemmas 6.1.4.5 and 6.1.4.6, taking into account the fact that $\log \mid \zeta-$ $z \mid<0$, we obtain

$$
\begin{equation*}
\left|\psi_{3}(z, h)\right| \leq(3 / 2) C_{3}(\sigma, R) \sqrt{h} \log h \tag{6.1.4.12}
\end{equation*}
$$

Now consider the difference

$$
\Delta w:=w\left(z_{0}+\Delta z\right)-w\left(z_{0}\right),|\Delta z|<h<\sqrt{h} / 2 .
$$

It can be represented as

$$
\begin{equation*}
\Delta w=\Delta \psi_{1}+\Delta \psi_{2}+\Delta \psi_{3}+\Delta H(z, w) . \tag{6.1.4.13}
\end{equation*}
$$

Choosing $h$ small enough, one may assume that $z_{0}+\Delta z \in K_{\sqrt{h} / 2, z_{0}}$. Thus, according to (6.1.4.11),

$$
\begin{equation*}
\left|\Delta \psi_{2}\left(z_{0}, h\right)\right| \leq\left|\psi_{2}\left(z_{0}, h\right)\right|+\left|\psi_{2}\left(z_{0}+h, h\right)\right| \leq C_{6}(\sigma, R) \sqrt{h} . \tag{6.1.4.14}
\end{equation*}
$$

Likewise (6.1.4.12) yields

$$
\begin{equation*}
\left|\Delta \psi_{3}\left(z_{0}, h\right)\right| \leq\left|\psi_{3}\left(z_{0}, h\right)\right|+\left|\psi_{3}\left(z_{0}+h, h\right)\right| \leq C_{7}(\sigma, R) \sqrt{h} \log h . \tag{6.1.4.15}
\end{equation*}
$$

Finally, from (6.1.4.10) and (6.1.4.8), respectively, we obtain

$$
\begin{equation*}
\left|\Delta \psi_{1}\right| \leq C_{3}(\sigma, R) \sqrt{h},\left|\Delta H\left(z_{0}, w\right)\right| \leq C_{8}(\sigma, R) h \tag{6.1.4.16}
\end{equation*}
$$

Substituting (6.1.4.14)-(6.1.4.16) into (6.1.4.13), we obtain relation 5 of Lemma 6.1.4.2.

Thus we have completed the proof of Lemma 6.1.4.2.
6.1.5 In this item we are going to prove Theorem 6.1.2.6. However, before that, we prove
Lemma 6.1.5.1 Let $w_{n}=\mathcal{G}_{H} v_{n}, v_{n} \xrightarrow{\mathcal{D}^{\prime}} v$ and $w_{n} \xrightarrow{\mathcal{D}^{\prime}} w_{\infty}$. Set

$$
E_{\infty}:=\left\{z: w_{\infty}(z)=H(z)-v(z)\right\}
$$

Then $w_{\infty}$ is harmonic in $\mathbb{C} \backslash E_{\infty}$.
Let us note that $w_{\infty}$, in general, is not the maximal subharmonic minorant because the operator $\mathcal{G}_{H}$ can be only upper semicontinuous, as will be demonstrated by example in Section 6.1.8. However, $w_{\infty}$ is a minorant of $H-v$ because of Theorem 6.1.1.2, 1 .

Proof. Let $z_{0} \notin E_{\infty}$. Then there exists a $\delta>0$ such that

$$
w_{\infty}\left(z_{0}\right)+v\left(z_{0}\right) \leq H\left(z_{0}\right)-2 \delta
$$

Since the function $b(z):=w_{\infty}+v(z)-H(z)$ is upper semicontinuous, there exists an $\epsilon=\epsilon(\delta)$ such that $b(z)<-\delta$ for all $z \in\left\{\left|z-z_{0}\right|<2 \epsilon\right\}$.

Let $(\bullet)_{\epsilon}$ be a smoothing operator from (2.6.2.3). If $w_{n} \xrightarrow{\mathcal{D}^{\prime}} w$ then $\left(w_{n}\right)_{\epsilon} \rightarrow w_{\epsilon}$ uniformly on every compact set (Theorem 2.3.4.5, reg3) and for every subharmonic function $v$ the sequence $v_{\epsilon}(z) \downarrow v(z)$, when $\epsilon \downarrow 0$ (Theorem 2.6.2.3, ap2).

Then $(b)_{\epsilon}(z)<-\delta$, for $\left|z-z_{0}\right|<\epsilon$ or $\left(w_{\infty}\right)_{\epsilon}(z)+(v)_{\epsilon}(z) \leq(H)_{\epsilon}(z)-\delta$. The function $H$ is continuous, hence uniformly continuous on the circle $\left\{z:\left|z-z_{0}\right| \leq\right.$ $\epsilon\}$. Thus we can replace $(H)_{\epsilon}$ in the last inequality with $H$ and $\delta$ with $\delta / 2$. So, we have

$$
\begin{equation*}
\left(w_{\infty}\right)_{\epsilon}(z)+v_{\epsilon}(z) \leq H(z)-\delta / 2,\left|z-z_{0}\right|<\epsilon \tag{6.1.5.1}
\end{equation*}
$$

Since $(\bullet)_{\epsilon}$ is monotonic on subharmonic functions, we can replace $\epsilon$ in (6.1.5.1) with any $\epsilon_{1}<\epsilon$. So we obtain

$$
\begin{equation*}
\left(w_{\infty}\right)_{\epsilon_{1}}(z)+v_{\epsilon_{1}}(z) \leq H(z)-\delta / 2,\left|z-z_{0}\right|<\epsilon \tag{6.1.5.2}
\end{equation*}
$$

Since $\left(w_{n}\right)_{\epsilon_{1}} \rightarrow\left(w_{\infty}\right)_{\epsilon_{1}}$ uniformly in the disc $\left|z-z_{0}\right| \leq \epsilon$ we can replace in (6.1.5.2) $w_{\infty}$ with $w_{n}$ and respectively $v$ with $v_{n}$, changing $\delta / 2$ with $\delta / 4$. After that we can pass to the limit as $\epsilon_{1} \downarrow 0$ for every sufficiently large $n$. So we obtain

$$
w_{n}(z)+v_{n}(z) \leq H(z)-\delta / 4,\left|z-z_{0}\right|<\epsilon
$$

It means that the disc $\left\{\left|z-z_{0}\right|<\epsilon\right\} \subset \mathbb{C} \backslash E_{v_{n}}$. Because of Lemma 6.1.4.2, 3, $w_{n}$ is harmonic in this disc for all large $n$. Thus $w_{\infty}$ is also harmonic, as the $\mathcal{D}^{\prime}$-limit of $w_{n}$.

Proof of Theorem 6.1.2.6. Let $v_{n} \xrightarrow{\mathcal{D}^{\prime}} v$. Then the set $w_{n}=\mathcal{G}_{H} v_{n}$ is equicontinuous by Lemma 6.1.4.2, 5 , and we can choose from it a subsequence uniformly converging to a continuous function $w_{\infty}$. Let $w=\mathcal{G}_{H} v, E=E_{w}, E_{\infty}$ being defined in Lemma 6.1.5.1.

Since

$$
\left(w_{\infty}+v\right)(z) \leq H(z),(w+v)(z) \leq H(z), \forall z \in \mathbb{C}
$$

and $v$ is upper semicontinuous, whereas $w$ and $H$ are continuous, the sets $E$ and $E_{\infty}$ are closed.

Since $w_{\infty}(z) \leq H(z)-v(z)$ we have

$$
\begin{equation*}
w_{\infty}(z) \leq w(z), \forall z \in \mathbb{C}, \tag{6.1.5.3}
\end{equation*}
$$

and therefore $E_{\infty} \subset E$.
The function $w$ is subharmonic in $\mathbb{C} \backslash E_{\infty}$, and $w_{\infty}$ is harmonic in $\mathbb{C} \backslash E_{\infty}$ by Lemma 6.1.5.1. They take the same values on $E_{\infty}$. As the maximum principle holds in $\mathcal{H}_{w_{\infty}}$ by assumption we have, according to Lemma 6.1.4.1 the inequality

$$
\begin{equation*}
w(z) \leq w_{\infty}(z), \forall z \in \mathbb{C} \tag{6.1.5.4}
\end{equation*}
$$

The inequalities (6.1.5.4) and (6.1.5.3) imply that $w(z)=w_{\infty}(z)$, i.e., $\mathcal{G}_{H}$ is continuous.

### 6.1.6

Proof of Theorem 6.1.2.5. Sufficiency. We exploit the following criterion for existence of a limit set that follows from Theorems 4.2.1.1, 4.2.1.2, 4.3.1.2 and Corollary 5.3.1.5:

Proposition 6.1.6.1 In order that $U \subset U[\rho]$ be a limit set of an entire function $f \in A(\rho(r))$ it is necessary and sufficient that there exists a piecewise continuous, $\omega$-dense in $U$ asymptotically dynamical pseudo-trajectory (a.d.p.t) $v(\bullet \mid t)$.

Exercise 6.1.6.1 Check this.
Let $v_{\Phi}(\bullet \mid t)$ be an a.d.p.t. corresponding to $\operatorname{Fr} \Phi$. Consider the pseudo-trajectory $v_{g}(\bullet \mid t):=\mathcal{G}_{H} v_{\Phi}(\bullet \mid t)$. It exists because of (6.1.1.1). Prove that this pseudotrajectory is asymptotically dynamical, i.e., (4.3.1.1) is fulfilled. Recall that $T_{\tau} \bullet=(\bullet)_{\left[e^{\tau}\right]}$.

Using the property of invariance of $\mathcal{G}_{H}$ (Theorem 6.1.1.2, 2) we have

$$
T_{\tau} v_{g}\left(\bullet \mid e^{t}\right)-v_{g}\left(\bullet \mid e^{t+\tau}\right)=\mathcal{G}_{H}\left[T_{\tau} v_{\Phi}\left(\bullet \mid e^{t}\right)-v_{\Phi}\left(\bullet \mid e^{t+\tau}\right)\right] .
$$

Thus (4.3.1.1) is fulfilled because of continuity of $\mathcal{G}_{H}$. Also the condition of $\omega$ denseness (4.3.1.4) is fulfilled and

$$
\left\{w \in U[\rho]:\left(\exists t_{j} \rightarrow \infty\right) w=\mathcal{D}^{\prime}-\lim v_{g}\left(\bullet \mid e^{t_{j}}\right)\right\}=\mathcal{G}_{H}(\mathbf{F r} \Phi) .
$$

The corresponding entire function $g \in A(\rho(r))$ with the limit set $U_{g}=\mathcal{G}_{H}(\mathbf{F r} \Phi)$ is an ideally complementing multiplicator, because

$$
\operatorname{Fr}[g \Phi]=\left\{v+\mathcal{G}_{H} v: v \in \mathbf{F r}[\Phi]\right\} .
$$

## Exercise 6.1.6.2 Check this.

Necessity. Let $G$ be an ideally complementing multiplicator of $\Phi$. Let us show that $\mathcal{G}_{H}$ is continuous on $\operatorname{Fr}[\Phi]$. Assume this is not the case, i.e., there exists a sequence $v_{j} \rightarrow v$ such that $\mathcal{G}_{H} v_{j} \rightarrow W$ and $W \neq \mathcal{G}_{H} v$. Since the limit set $\operatorname{Fr}[G \Phi]$ is closed, we have $v_{j}+\mathcal{G}_{H} v_{j} \rightarrow v+\mathcal{G}_{H} v, v_{j} \in \operatorname{Fr}[\Phi]$. On the other hand, $v_{j}+\mathcal{G}_{H} v_{j} \rightarrow v+W$. Thus, $W=\mathcal{G}_{H} v$, which is a contradiction.

Proof of Theorem 6.1.2.3. Let $\operatorname{Fr}[\Phi]$ be a periodic limit set, that is

$$
\operatorname{Fr}[\Phi]=\mathbb{C}(v)=\left\{v_{[t]}: 1 \leq t \leq e^{P}\right\}
$$

where $v \in U[\rho]$. We shall show that $\mathcal{G}_{H}$ is continuous on $U[\rho]$. By Theorem 6.1.1.2, 2) the equality $\left(\mathcal{G}_{H} v\right)_{[t]}=\mathcal{G}_{H} v_{[t]}$ holds. Since the operation $(\bullet)_{[t]}$ is continuous for all $t, \mathcal{G}_{H}$ is continuous on $\mathbb{C}(v)$.
6.1.7 Now we are going to prove Theorem 6.1.2.4. However, we need some preparation.

Let $h(\phi), \phi \in[0,2 \pi)$ be a $2 \pi$-periodic $\rho$-t.c.function, satisfying the condition

$$
\max _{\phi \in[0,2 \pi]} h(\phi)=\sigma
$$

We denote this class as $T C[\rho, \sigma]$ and write

$$
T C[\rho]:=\bigcup_{\sigma>0} T C[\rho, \sigma]
$$

The class of functions $w=h_{1}-h_{2}$ where $h_{1}, h_{2} \in T C[\rho, \sigma]$ will be denoted as $\delta T C[\rho, \sigma]$ and we will also write

$$
\delta T C[\rho]:=\bigcup_{\sigma>0} \delta T C[\rho, \sigma] .
$$

From properties of a $\rho$-t.c.function (see Sections 3.2.3-3.2.5) we can obtain the following properties of $\delta-\rho$-t.c.functions:

Proposition 6.1.7.1 For $w \in \delta T C[\rho]$ the following holds:

1. $w^{\prime}(\phi-0)$ and $w^{\prime}(\phi+0)$ exist at each point and are bounded in $[0 ; 2 \pi]$;
2. $w^{\prime}(\phi-0)=w^{\prime}(\phi+0)$ for all $\phi \in[0 ; 2 \pi]$, except, perhaps, a countable set;
3. the charge $\Delta_{w}$ generated by the function

$$
\Delta_{w}:=w^{\prime}(\phi)+\rho^{2} \int^{\phi} w(\theta) d \theta
$$

has bounded variation $\left|\Delta_{w}\right| ;$ the variation $\left|\Delta_{w}\right|(\alpha, \beta)$ of the charge on the interval $(\alpha ; \beta)$ and the variation of the charge generated by derivative $\left|w^{\prime}\right|(\alpha ; \beta)$ on the same interval satisfy the relation

$$
\left|\Delta_{w}\right|(\alpha, \beta) \geq\left|w^{\prime}\right|(\alpha ; \beta)+\rho^{2}(\beta-\alpha)
$$

4. For all $w \in \delta T C[\rho, \sigma]$,

$$
\max \left(\left|w^{\prime}(\phi-0)\right|,\left|w^{\prime}(\phi+0)\right|\right) \leq C(\rho, \sigma), \phi \in[0 ; 2 \pi] ;
$$

5. if $r^{\rho} w_{n} \xrightarrow{\mathcal{D}^{\prime}} r^{\rho} w$ and $w_{n} \in \delta T C[\rho, \sigma]$, then $w_{n} \rightarrow w$ uniformly on $[0 ; 2 \pi]$.

Exercise 6.1.7.1 Prove this using properties of $\rho$-t.c.functions.
We also need a technical
Lemma 6.1.7.2 Let $M_{n}(\phi)$ be a sequence of functions that satisfy the conditions:

1. $M_{n} \geq 0 ; M_{n}(0)=0$;
2. $M_{n}$ converges uniformly to $M_{\infty}(\phi) \geq A \sin \rho \phi, A>0$;
3. $M_{n}^{\prime}(\phi-0), M_{n}^{\prime}(\phi+0)$ exist at every point, and they coincide almost everywhere;
4. there exists a sequence $\phi_{n} \downarrow 0$ such that for each arbitrarily small $\epsilon>0$ and arbitrarily large $n_{0} \in \mathbb{N}$ there exists $n>n_{0}$ for which the inequality $M_{n}\left(\phi_{n}\right)<\epsilon \phi_{n}$ holds.
Then there exists a sequence $\left(\zeta_{n}, \eta_{n}\right)$ of disjoint intervals and a subsequence $M_{k_{n}}$ such that

$$
\begin{equation*}
M_{k_{n}}^{\prime}\left(\zeta_{n}\right)-M_{k_{n}}^{\prime}\left(\eta_{n}\right) \geq A \rho / 2 \tag{6.1.7.1}
\end{equation*}
$$

Proof. Set $\epsilon_{0}=1 / 2, \eta_{0}=\pi / 4$ and choose the required sequence recurrently. Let $\epsilon_{n}, \eta_{n}, \zeta_{n}$ be already chosen. Set $\epsilon_{n+1}=\epsilon_{n} / 2$, find $\phi_{n+1}<\eta_{n}$ and choose $k_{0}=k_{0}(n)$ so that for $k>k_{0}$,

$$
M_{k}\left(\phi_{n+1}\right)-A \rho \phi_{n+1}>-\epsilon_{n+1} \phi_{n+1} .
$$

This is possible because of condition 2 and $\sin \rho \phi \sim \rho \phi, \phi \rightarrow 0$. So we have

$$
\begin{equation*}
\frac{M_{k}\left(\phi_{n+1}\right)}{\phi_{n+1}}>A \rho-\epsilon_{n+1} \tag{6.1.7.2}
\end{equation*}
$$

Now, choose $\psi_{n+1}<\phi_{n+1}$ and $k_{n+1}>k_{0}$ so that

$$
\begin{equation*}
M_{k_{n+1}}\left(\psi_{n+1}\right)<\epsilon_{n+1} \psi_{n+1} \tag{6.1.7.3}
\end{equation*}
$$

This is possible by condition 4 . Thus for small $\epsilon_{n+1}$ from (6.1.7.2) and (6.1.7.3) we obtain

$$
\begin{equation*}
\frac{M_{k_{n+1}}\left(\phi_{n+1}\right)-M_{k_{n+1}}\left(\psi_{n+1}\right)}{\phi_{n+1}-\psi_{n+1}}>(2 / 3) A \rho . \tag{6.1.7.4}
\end{equation*}
$$

On the other hand

$$
\begin{equation*}
\frac{M_{k_{n+1}}\left(\psi_{n+1}\right)-M_{k_{n+1}}(0)}{\psi_{n+1}-0}<\epsilon_{n+1} \tag{6.1.7.5}
\end{equation*}
$$

On the interval $\left(\psi_{n+1}, \phi_{n+1}\right)$ there is a point $\eta_{n+1}$ where the derivative exists and the inequality

$$
\begin{equation*}
M_{k_{n+1}}^{\prime}\left(\eta_{n+1}\right) \geq \frac{M_{k_{n+1}}\left(\phi_{n+1}\right)-M_{k_{n+1}}\left(\psi_{n+1}\right)}{\phi_{n+1}-\psi_{n+1}} \tag{6.1.7.6}
\end{equation*}
$$

is valid. Also there is a point $\zeta_{n+1} \in\left(0, \psi_{n+1}\right)$ where the derivative exists and the inequality

$$
\begin{equation*}
M_{k_{n+1}}^{\prime}\left(\zeta_{n+1}\right) \leq \frac{M_{k_{n+1}}\left(\psi_{n+1}\right)-M_{k_{n+1}}(0)}{\psi_{n+1}-0} \tag{6.1.7.7}
\end{equation*}
$$

is valid.
From the inequalities (6.1.7.4)-(6.1.7.7) we obtain (6.1.7.1).
Proof of Theorem 6.1.2.4. Denote by $\hat{\mathcal{G}}_{H} h$ the maximal $\rho$-t.c.minorant of $H\left(e^{i \phi}\right)-h(\phi)$. It follows from Theorem 6.1.1.2, 2 that

$$
\mathcal{G}_{H}\left(r^{\rho} h(\phi)\right)\left(r e^{i \phi}\right)=r^{\rho} \hat{\mathcal{G}}_{H} h(\phi) .
$$

Exercise 6.1.7.2 Prove this.
So taking in consideration Proposition 6.1.7.1, 5 , one must prove
Proposition 6.1.7.3 The operator $\hat{\mathcal{G}}_{H}$ is continuous on the set

$$
\hat{U}_{\mathrm{ind}}:=\{h(\phi, c): c \in C\}
$$

in the uniform topology.
Proof. Let $h_{n} \rightarrow h, h_{n}, h \in \hat{U}_{\text {ind }}$. Set $\hat{w}_{n}=\hat{\mathcal{G}}_{H} h_{n}, \hat{w}=\hat{\mathcal{G}}_{H} h, \hat{w}_{\infty}=\lim _{n \rightarrow \infty} \hat{w}_{n}$. We set also $\hat{M}_{n}=H-h_{n}-\hat{w}_{n}, \hat{M} \infty=H-h-\hat{w}_{\infty}, \hat{M}=H-h-\hat{w}$. Let $\left(\alpha_{n} ; \beta_{n}\right)$ be a maximum interval where $\hat{M}_{n}(\phi)>0$. We shall show that $\beta_{n}-\alpha_{n} \leq \pi / \rho$. Indeed, for a fixed $n$ let us consider the function

$$
\hat{W}_{n}:=\hat{w}_{n}+\epsilon_{n} L\left(\phi-\left(\alpha_{n}+\beta_{n}\right) / 2\right)
$$

where

$$
L(\phi)= \begin{cases}\cos |\phi|, & \phi \in(-\pi / 2 \rho ; \pi / 2 \rho) \\ 0, & \phi \in[-\pi ; \pi] \backslash(-\pi / 2 \rho ; \pi / 2 \rho)\end{cases}
$$

and $\epsilon_{n}$ is small enough. If $\beta_{n}-\alpha_{n}>\pi / \rho$, then $\hat{W}_{n}$ is also a $\rho$-t.c.minorant of $H-h_{n}$, i.e., $\hat{w}_{n}$ is not maximal.

If $\beta_{n}-\alpha_{n}=\pi / \rho$, then, to ensure that $\hat{w}_{n}$ is a maximal minorant, at least one of the conditions

$$
\begin{equation*}
\liminf _{\phi \rightarrow \alpha_{n}+0} \frac{\hat{M}_{n}(\phi)}{\phi-\alpha_{n}}=0, \liminf _{\phi \rightarrow \beta_{n}-0} \frac{\hat{M}_{n}(\phi)}{\beta_{n}-\phi}=0 \tag{6.1.7.8}
\end{equation*}
$$

must be satisfied.
Let us choose (and preserve the previous notation) a subsequence $\hat{M}_{n}(\phi)$ for which $\alpha_{n} \rightarrow \alpha$, and $\beta_{n} \rightarrow \beta$.

If $\beta-\alpha<\pi / \rho$, then the maximum principle for $\rho$-t.c.functions is valid. Repeating arguments of proof of Theorem 6.1.2.6, we obtain $\hat{w}_{\infty}=\hat{w}$ for all $\phi$, which proves Proposition 6.1.7.3 for the case considered.

Exercise 6.1.7.3 Repeat them.
Consider the case when $\beta-\alpha=\pi / \rho$. Set $q(\phi)=\left(w-w_{\infty}\right)(\phi)$. The function $q$ is $\rho$-trigonometric on the interval $(\alpha ; \beta)$ since $\hat{w}$ and $\hat{w}_{\infty}$ are $\rho$-trigonometric, i.e., have the form $A \sin \rho \phi+B \cos \rho \phi$.

Exercise 6.1.7.4 Explain this.
Besides, we have $q \geq 0$ and $q(\alpha)=q(\beta)=0$. It is easy to see that $q$ has the form

$$
\begin{equation*}
q(\phi)=A \sin \rho(\phi-\alpha), A>0 \tag{6.1.7.9}
\end{equation*}
$$

Exercise 6.1.7.5 Prove this.
Since $\hat{M}(\phi) \geq 0$, we have $\hat{M}_{\infty}(\phi)=\hat{M}(\phi)+\left(\hat{w}-\hat{w}_{\infty}\right)(\phi) \geq\left(\hat{w}-\hat{w}_{\infty}\right)(\phi), \forall \phi$, whence

$$
\begin{equation*}
\hat{M}_{\infty}(\phi) \geq A \sin \rho(\phi-\alpha), A>0 \tag{6.1.7.10}
\end{equation*}
$$

Since $\beta_{n}-\alpha_{n} \leq \pi / \rho$, the segment $[\alpha, \beta]$ contains the infinite sequence $\alpha_{n}$ or $\beta_{n}$. Let us single out a subsequence, let it be, for example, $\alpha_{n} \rightarrow \alpha+0, \alpha_{n} \in[\alpha ; \beta]$.

Consider the sequence $M_{n}(\phi)=\hat{M}_{n}(\phi-a)$. From the definition of $M_{n}$ and from relation (6.1.7.10) it follows that conditions 1 and 2 of Lemma 6.1.7.2 are fulfilled. Condition 3 is fulfilled because of property 1 of Lemma 6.1.7.2. Further, if $\alpha_{n} \not \equiv \alpha$, then condition 4 of Lemma 6.1.7.2 is trivially true,since $M_{n}\left(\alpha_{n}-\alpha\right)=0$; otherwise, if $\alpha_{n} \equiv \alpha$, condition 4 follows from (6.1.7.8).

Applying Lemma 6.1.7.2, we obtain the union of intervals satisfying (6.1.7.7). The equality $H(\phi)=\hat{M}_{n}-h_{n}-\hat{w}_{n}$ yields the following inequality for the measure $\Delta_{H}: \Delta_{H}\left(\left(\eta_{n} ; \zeta_{n}\right)\right) \geq A \rho / 2$. Summing this inequality and taking into account the fact that the intervals do not intersect, we obtain $\Delta_{H}\left(\cup_{n}\left(\eta_{n}, \zeta_{n}\right)\right)=\infty$, which is impossible. So, Proposition 6.1.7.3 is proved.

Hence, Theorem 6.1.2.4 is proved.
6.1.8 In this item we show an example of $H$ and an entire function without an ideally complementing $H$-multiplicator.

According to Theorem 6.1.2.5, to construct such an example it is sufficient to construct a limit set on which $\mathcal{G}_{H}$ is not continuous.

We set

$$
L(\eta)= \begin{cases}\cos |\eta|, & \eta \in(-\pi / 2 \rho ; \pi / 2 \rho) \\ 0, & \eta \in[-\pi ; \pi] \backslash(-\pi / 2 \rho ; \pi / 2 \rho)\end{cases}
$$

Let us define $X \in C^{\infty}$ so that $X(\xi)=1$ for $\xi<0$ and $X=0$ for $\xi>\alpha$.
We set

$$
\begin{equation*}
\kappa:=\left(1 / \rho^{2}\right) \max _{(-\infty ;+\infty)}\left[2 \rho X^{\prime}+X^{\prime \prime}\right](\xi), \quad H_{0}(\eta):=L(\eta)+\kappa \tag{6.1.8.1}
\end{equation*}
$$

We also set

$$
v(\zeta, c):=\left[H_{0}-X(\xi-c) L(\eta)\right] e^{\rho \xi}, \quad \zeta=\xi+i \eta
$$

where $H_{0}$ and $L$ have been periodically extended from the interval $[-\pi ; \pi]$ to $(-\infty,+\infty)$.

As $H(z)$ we take

$$
H(z):=H_{0}(\phi) r^{\rho}
$$

Lemma 6.1.8.1 We have

$$
\begin{gather*}
v(\log z, c) \in U[\rho, \sigma], \quad \sigma=1+\kappa  \tag{6.1.8.2}\\
\mathcal{G}_{H} v(\bullet, c) \equiv 0  \tag{6.1.8.3}\\
\lim _{c \rightarrow \infty} v(\log z, c)=\kappa r^{\rho} \tag{6.1.8.4}
\end{gather*}
$$

uniformly with respect to $z \in K \Subset \mathbb{C}$, and

$$
\begin{equation*}
\mathcal{G}_{H}\left(\kappa r^{\rho}\right)=L(\phi) r^{\rho} . \tag{6.1.8.5}
\end{equation*}
$$

Proof. For the Laplace operators in $\zeta$ and $z$ it is true that $\Delta_{\zeta}=\Delta_{z} /|\zeta|^{2}$. Let us check that $v(\zeta, c)$ is subharmonic in $\zeta$. We have

$$
\Delta_{\zeta} v(\zeta, c)=\left\{[1-X(\xi-c)]\left(L^{\prime \prime}+\rho^{2} L\right)(\eta)+\left[\rho^{2} \kappa-L(\xi)\left[X^{\prime \prime}(x-c)+2 \rho X^{\prime}(x-c)\right]\right\} e^{\rho \xi} .\right.
$$

Exercise 6.1.8.1 Check this computation.
Since $X(\xi) \leq 1$ and $L(\eta)$ is $\rho$-t.c.,

$$
[1-X(\xi-c)]\left(L^{\prime \prime}+\rho^{2} L\right)(\eta) \geq 0
$$

Since $L(\xi) \leq 1$ and $\left[X^{\prime \prime}(x-c)+2 \rho X^{\prime}(x-c)\right] \leq \kappa \rho^{2}$ we have

$$
\left[\rho^{2} \kappa-L(\xi)\left[X^{\prime \prime}(x-c)+2 \rho X^{\prime}(x-c)\right] \geq 0\right.
$$

Thus $v(\log z, c)$ is subharmonic.

Exercise 6.1.8.2 Prove that $v(\log z, c) \in U[\rho, \sigma]$ for $\sigma=1+\kappa$.
Let us prove (6.1.8.3). We have

$$
H(z)-v(\log z, c)=X(\log r-c) L(\phi) r^{\rho}
$$

Since $X=0$ for $r>e^{c+\alpha}$, the maximal subharmonic minorant of $H-v$ is zero by the maximum principle.

Relation (6.1.8.4) is obvious, since $X(\log r-c)$ converges to 1 uniformly on every disc $\{|z| \leq R\}$. Relation 6.1.8.5 follows from the equality $H(z)-\kappa r^{\rho}=$ $L(\phi) r^{\rho}$, since $L(\phi) r^{\rho} \in U[\rho]$.

Now we pass to the construction of the example. Examine the set

$$
U_{1}:=\operatorname{clos}\{v(\log z, c): c \in[0 ; \infty)\} .
$$

It contains the function

$$
\mathcal{D}^{\prime}-\lim _{c \rightarrow \infty} v(\log z, c)=\kappa r^{\rho} .
$$

Let us consider the minimal convex $(\bullet)_{[t]}$-invariant set $U$ containing $U_{1}$. The set is contained in $U[\rho, 1+\kappa]$. It is a limit set for a certain entire function $\Phi$. Let us show that $\mathcal{G}_{H}$ is not continuous on $\operatorname{Fr}[\Phi]$. We take an arbitrary sequence $c_{j} \rightarrow \infty$ and set $v_{j}(z):=v\left(\log z, c_{j}\right) \in U$. Now $\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} v_{j}=\kappa r^{\rho}$ by (6.1.8.4) and $\mathcal{G}_{H} v_{j}(z)=0$, so $\mathcal{D}^{\prime}-\lim _{j \rightarrow \infty} \mathcal{G}_{H} v_{j}=0$ but

$$
\mathcal{G}_{H}\left(\lim v_{j}\right)=\mathcal{G}_{H}\left(\kappa r^{\rho}\right)=L(\phi) r^{\rho} \not \equiv 0
$$

which shows the lack of continuity.
By virtue of Theorem 6.1.2.5, $\Phi$ is not ideally complementable.
6.1.9 Here we prove existence and continuity of the maximal subharmonic minorant for some classes of functions $m(z)$.

Theorem 6.1.9.1 Let $m(z)$ be a continuous function such that the set of subharmonic minorants is nonempty. Then the maximal subharmonic minorant of $m$ exists and is continuous.

Proof. The set of subharmonic minorants is nonempty and partially ordered. Indeed, for every subset $\left\{u_{\alpha}, \alpha \in A\right\}$ of subharmonic minorants there exists $u_{A}=$ $\left(\sup \left\{u_{\alpha}: \alpha \in A\right\}\right)^{*}$ which is subharmonic and is a minorant of $m$, because $m$ is continuous.

Exercise 6.1.9.1 Explain this in detail.

Thus there exists a uniquely maximal element m.s.m. $(z, m)$, which is a subharmonic minorant of $m$.

Let us prove that it is continuous at every point $z_{0}$. Since m.s.m. $(z, m)$ is upper semicontinuous,

$$
\text { m.s.m. }(z, m)>\text { m.s.m. }\left(z_{0}, m\right)-\epsilon
$$

for $\left|z-z_{0}\right|<\delta$ for arbitrarily small $\epsilon$ and corresponding $\delta=\delta(\epsilon)$. So we need to prove the inequality

$$
\text { m.s.m. }(z, m)<\text { m.s.m. }\left(z_{0}, m\right)+\epsilon
$$

for arbitrarily small $\epsilon$ and corresponding $\delta=\delta(\epsilon)$.
Perform sweeping m.s.m. $(z, m)$ from the disc $\left|z-z_{0}\right|<\delta$ such that the result $u(z, \delta)$ satisfies the inequality

$$
\text { m.s.m. }(z, m)<u(z, \delta)<\text { m.s.m. }(z, m)+\epsilon<m(z)+\epsilon .
$$

Thus $u(z, \delta)-\epsilon<m(z)$. Hence m.s.m. $(z, m)>u(z, \delta)-\epsilon$ for all $z$. Since $u(z, \delta)$ is continuous, $u(z, \delta)>u\left(z_{0}, \delta\right)-\epsilon$ in the disc $\left\{\left|z-z_{0}\right|<\delta_{1}\right\}$. So m.s.m. $(z, m)>$ $u\left(z_{0}, \delta\right)-\epsilon>$ m.s.m. $\left(z_{0}, m\right)-\epsilon$.

Theorem 6.1.9.2 Let $m=m_{1}-m_{2}$, where $m_{1}, m_{2}$ are subharmonic functions. Then the maximal subharmonic minorant of $m$ exists. If $m_{1}$ is continuous, then the maximal subharmonic minorant is continuous.

Proof. Set $\mathcal{M}_{\epsilon}(z, m):=\mathcal{M}_{\epsilon}\left(z, m_{1}\right)-\mathcal{M}_{\epsilon}\left(z, m_{2}\right)$, where $\mathcal{M}_{\epsilon}\left(z, m_{i}\right), i=1,2$ is defined by (2.6.1.1). Since $\mathcal{M}_{\epsilon}(z, m)$ is continuous (see Theorem 2.6.2.3 (Smooth approximation)), there exists m.s.m. $\left(z, \mathcal{M}_{\epsilon}(z, m)\right)$. We have

$$
u(z, m):=\limsup _{\epsilon \rightarrow 0} \text { m.s.m. }\left(z, \mathcal{M}_{\epsilon}(\bullet, m)\right) \leq \lim _{\epsilon \rightarrow 0} \mathcal{M}_{\epsilon}(z, m)=m_{1}-m_{2}(z)=m(z)
$$

Now we prove that the upper semicontinuous regularization $u^{*}(z, m)$ also satisfies the inequality $u^{*}(z, m) \leq m(z)$. Indeed, $m_{2}+u(z, m) \leq m_{1}(z)$. Hence,

$$
\mathcal{M}_{\epsilon}\left(z, m_{2}\right)+\mathcal{M}_{\epsilon}(z, u(\bullet, m)) \leq \mathcal{M}_{\epsilon}\left(z, m_{1}\right)
$$

Passing to the limit we obtain three subharmonic functions and inequality

$$
m_{2}(z)+u^{*}(z, m) \leq m_{1}(z)
$$

We prove that $u^{*}(z, m)$ is the m.s.m. $(z, m)$. If not, there would exist a subharmonic function $u_{1}$ which exceeds $u^{*}(z, m)$ on a set of positive measure (otherwise they coincide); thus we would have for some $z$ and $\epsilon$,

$$
u^{*}(z, m)<\mathcal{M}_{\epsilon}\left(z, u_{1}\right) \leq m \cdot \operatorname{s.m} \cdot\left(z, \mathcal{M}_{\epsilon}(z, m)\right)
$$

This contradicts the definition of $u^{*}(z, m)$.

Now suppose that $m_{1}$ is continuous at a point $z_{0}$. From Theorem 2.6.5.1 (Jensen-Privalov) we obtain that it is equivalent to

$$
\int_{0}^{\epsilon} \frac{\mu_{m_{1}}\left(\left\{z:\left|z-z_{0}\right|<t\right\}\right)}{t} d t=o(1), \epsilon \rightarrow 0
$$

Similarly to the proof of Theorem 6.1.4.2, 5, we obtain

$$
\mu_{\text {m.s.m. }(z, m)} \leq \mu_{m_{1}} .
$$

Hence m.s.m. $(z, m)$ is also continuous.

Exercise 6.1.9.2 Prove continuity in detail.
Corollary 6.1.9.3 For $m=m(z, v, H)$ the function $\mathcal{G}_{H} v(z):=$ m.s.m. $(z, m)$ exists and is continuous.

Exercise 6.1.9.3 Prove Corollary 6.1.9.3.

### 6.2 A generalization of $\rho$-trigonometric convexity

6.2.1 One of the important and useful kinds of limit sets is periodic limit sets. They are determined by one subharmonic function $v \in U[\rho]$ that satisfies the condition

$$
\begin{equation*}
v(T z)=T^{\rho} v(z), z \in \mathbb{C} . \tag{6.2.1.1}
\end{equation*}
$$

Such a function is called automorphic. They generate the class of so-called $L_{\rho^{-}}$ subfunctions, that is a generalization of $\rho$-trigonometrically convex functions. In this part we are going to review properties of such functions from different points of view that will be useful for applications (see [ADP]).

In connection with property (6.2.1.1) it is natural to consider so-called $T$ homogeneous domains in $\mathbb{C}$, i.e., such domains $G$ that satisfy the condition $\{T z$ : $z \in G\}=G$ or shortly $T G=G$. As we can see they are invariant with respect to dilation by $T$. For example, every component of an open set of harmonicity of an automorphic function is a $T$-homogeneous domain.

Let $v$ satisfy (6.2.1.1). Then the function

$$
\begin{equation*}
q(z):=v\left(e^{z}\right) e^{-\rho x} \tag{6.2.1.2}
\end{equation*}
$$

is a $2 \pi$ periodic function in $y$ and $P$-periodic in $x$, where $P=\log T$.
The function $q$ can be considered as a function on a torus $\mathbb{T}_{P}^{2}$, obtained by identifying the opposite sides of the rectangle $\Pi=(0, T) \times(-\pi, \pi)$.

The homology group of $\mathbb{T}_{P}^{2}$ is nontrivial, and generated by the cycles $\gamma_{x}, \gamma_{y}$, where $\gamma_{x}=\mathbb{T}_{P}^{2} \cap\{y=0\}, \gamma_{y}=\mathbb{T}_{P}^{2} \cap\{x=0\}$.

Let $\pi$ be the covering map of $\mathbb{C}$ onto $\mathbb{T}_{P}^{2}$, then $\phi=\pi \circ \log$ is a well-defined covering map of $\mathbb{C} \backslash\{0\}$ onto $\mathbb{T}_{P}^{2}$, where the group of deck transformations is given by the dilations by $T^{m}$ for $m \in \mathbb{Z}$. So if $G$ is a given $T$-homogeneous domain, then

$$
\begin{equation*}
D=\pi \circ \log G=\phi(G) \tag{6.2.1.3}
\end{equation*}
$$

is a domain in $\mathbb{T}_{P}^{2}$. On the other hand, not every domain in $\mathbb{T}_{P}^{2}$ has a $T$-homogeneous domain as its preimage under $\phi$. The preimage $\phi^{-1}(D)$ under $\phi$ is a possibly disconnected set which is invariant under dilations by $T^{m}$ for $m \in \mathbb{Z}$. An intrinsic description is given by the next proposition.
Proposition 6.2.1.1 Let $\gamma$ be a closed curve in a domain $D \subset \mathbb{T}_{P}^{2}$ that is homologous in $\mathbb{T}_{P}^{2}$ to a cycle $\gamma=n_{x} \gamma_{x}+n_{y} \gamma_{y}, n_{x}, n_{y} \in \mathbb{Z}$.

1. If $n_{x}=0$ for every such $\gamma$ in $D$, then

$$
\phi^{-1}(D)=\cup_{j=-\infty}^{\infty} G_{j}
$$

where $G_{j}=T^{j} G_{0}, G_{0}$ is an arbitrary connected component of $\phi^{-1}(D)$, and $G_{j} \cap G_{l}=\varnothing$ for $j \neq l$.
2. If there exists a curve $\gamma$ as above with $n_{x} \neq 0$, then

$$
\phi^{-1}(D)=\cup_{q=0}^{k-1} G_{q}
$$

where $k=\min \left|n_{x}\right|$ with the minimum taken over all such curves $\gamma ; G_{0}$ is an arbitrary component of $\phi^{-1}(D) ; G_{j}, j=0,1, \ldots, k-1$, are disjoint $T^{k}$-homogeneous domains, and for every $m \in \mathbb{Z}, T^{m} G_{0}=G_{q}$, provided $m=l k+q$, for some $q \in \mathbb{Z}, 0 \leq q \leq k-1, l \in \mathbb{Z}$.
We call domains as in part 2 of Proposition 6.2.1.1 connected on spirals. In particular, this proposition shows that for every $D$ connected on spirals, we can find a connected $T^{k}$-homogeneous domain that relates to $D$ by (6.2.1.3).

Let us give some examples. The domain $D^{\prime}=\mathbb{T}_{P}^{2} \cap\{|x-P / 2|<P / 4\}$ is not connected on spirals, whereas $D^{\prime \prime}=\mathbb{T}_{P}^{2} \cap\{|y|<\pi / 4\}$ is. It follows that $D^{\prime} \cap D^{\prime \prime}$ is not connected on spirals whereas $D^{\prime} \cup D^{\prime \prime}$ is.

The situation can be more complicated. Set

$$
\begin{aligned}
x^{\prime}(x, y, \alpha) & :=x \cos \alpha+y \sin \alpha ; \\
y^{\prime}(x, y, \alpha) & :=-x \sin \alpha+y \cos \alpha ; 0 \leq \alpha<\pi / 4 ; \\
P_{1} & :=(1 / 2)\left|x^{\prime}(P, 2 \pi,-\alpha)\right| ; \\
P_{2} & :=(1 / 2)\left|y^{\prime}(P, 2 \pi,-\alpha)\right| .
\end{aligned}
$$

Then $R^{\prime}=\left\{z^{\prime}=x^{\prime}+i^{\prime} y^{\prime}:-P_{1}<x^{\prime}<P_{1} ;-P_{2}<y^{\prime}<P_{2}\right\}$ is a fundamental rectangle for $\mathbb{T}_{P}^{2}$ in the corresponding coordinates. Set $f\left(y^{\prime}\right):=\left(P_{2}-y^{\prime}\right)^{-1}-$ $\left(y^{\prime}+P_{2}\right)^{-1}$ and $D_{0,0}:=\left\{z^{\prime}:-P_{2}<y^{\prime}<P_{2} ; f\left(y^{\prime}\right)<x^{\prime}<f\left(y^{\prime}\right)+d\right\}$ where $0<d<P_{1}$. Then the domains $D_{l, m}:=D_{0,0}+2 P_{1} l+2 P_{2} m i^{\prime}, l, m \in \mathbb{Z}$ are disjoint, and their union $D$ determines a domain $\hat{D} \subset \mathbb{T}_{P}^{2}$. This $\hat{D}$ is determined completely by the intersection of $D$ with the rectangle $R=(0, P) \times(-\pi, \pi)$. The domain $\hat{D}$ is not connected on spirals.

One more example. Consider the family of lines $L_{l}:=\{z=x+i y: y=$ $\pi /(k P) x+l \pi / k, x \in \mathbb{R}\}, l \in \mathbb{Z}$. It determines a closed curve (spiral) $\gamma$ on $\mathbb{T}_{P}^{2}$ with $n_{1}=k$. The open set $D_{k}=\left\{z:|z-\zeta|<\epsilon, \zeta \in L_{l}, l \in \mathbb{Z}\right\}, 0<\epsilon<P / 2 \sqrt{\pi^{2}+k^{2}}$, determines a domain $\hat{D}_{k}$ on $\mathbb{T}_{P}^{2}$ that is connected on spirals, and such that $\phi^{-1}\left(\hat{D}_{k}\right)$ consists of $k$ components, every one of them $T^{k}$-homogeneous.

Since the function $v$ in (6.2.1.1) is subharmonic, the function $q$ of (6.2.1.2) is upper semicontinuous and in the $D^{\prime}$ topology on $\mathbb{T}_{P}^{2}$ satisfies the inequality $L_{\rho} q \geq 0$, where

$$
\begin{equation*}
L_{\rho}:=\Delta+2 \rho \frac{\partial}{\partial x}+\rho^{2} . \tag{6.2.1.4}
\end{equation*}
$$

Such functions $q$ are called subfunctions with respect to $L_{\rho}$, or $L_{\rho}$-subfunctions. $L_{\rho} q$ is a positive measure on $\mathbb{T}_{P}^{2}$.
The operator $L_{\rho}$ arises naturally by changing variables $z \mapsto \log z$ in the Laplace operator $\Delta_{\zeta}$.

Exercise 6.2.1.1 Check this. Set $\zeta=e^{z}$.
Let us note that if $q$ depends only on the variable $y$, it is a $2 \pi$-periodic $\rho$ trigonometric convex function because $L_{\rho}$ turns into $T_{\rho}=(\bullet)^{\prime \prime}+\rho^{2}(\bullet)$ (cf. Section 3.2.3).
6.2.2 Consider the solution of the homogeneous boundary problem

$$
\begin{align*}
L_{\rho} q & =0 \quad \text { in } D \\
\left.q\right|_{\partial D} & =0, \tag{6.2.2.1}
\end{align*}
$$

where $D$ is a domain in $\mathbb{T}_{P}^{2}$ and $q$ is bounded in a neighborhood of $\partial D$ with boundary value zero quasi-everywhere. This is a spectral problem for a pencil of differential operators ([Ma]).

A solution of this problem can be defined for an arbitrary domain $D \subset \mathbb{T}_{P}^{2}$ with a boundary of positive capacity.

The spectrum of the problem (6.2.2.1) consists of those (complex) $\rho$ for which (6.2.2.1) holds for some function $q \not \equiv 0$. The minimal positive point of the spectrum $\rho(D)$ exists iff the spectrum exists. The spectrum exists iff the domain $D$ is connected on spirals. In this case $\rho(D)$ is the order of the minimal harmonic function in every one of the domains $G_{i}$ that corresponds to $D$ by Proposition 6.2.1.1.

The quantity $\rho(D)$ is strictly monotonic. It means that if two domains $D_{1}$, $D_{2} \in \mathbb{T}_{P}^{2}$ are such that $D_{1} \subset D_{2}$ and the capacity of $D_{2} \backslash D_{1}$ is positive, then $\rho\left(D_{2}\right)<\rho\left(D_{1}\right)$. For example, this is the case of $D_{2}=\{|y|<d, d<2 \pi\}$ and $D_{1}$ is the same strip without the segment $\{$ it : $0 \leq t \leq d\}$.

In connection with the multiplicator problem we considered the maximal subharmonic minorant of a function $m=H-v$ where $v$ is a $T$-automorphic function. From Theorem 6.1.1.2, 2 we can obtain that if $v$ is a $T$-automorphic function, then $\mathcal{G}_{H} v$ is also $T$-automorphic.

Exercise 6.2.2.1 Check this.
Thus for this case, finding $D_{H}$ in Theorem 6.1.1.1 is reduced to finding a maximal $L_{\rho}$-subfunction $q$ that satisfies the inequality

$$
\begin{equation*}
q(z) \leq m(z):=\left[H\left(e^{z}\right)-v\left(e^{z}\right)\right] e^{-\rho x}, z \in \mathbb{T}_{P}^{2} \tag{6.2.2.2}
\end{equation*}
$$

We say that $m(z)$ has an $L_{\rho}$-subminorant.
The idea of $\rho(D)$ gives a possibility for
Theorem 6.2.2.1 If $m$ has a non-zero $L_{\rho}$-subminorant, then $\rho(D) \leq \rho$ for some component $D$ of the open set $\mathcal{M}_{+}:=\{z: m(z)>0\}$.

Conversely, if $\rho(D)<\rho$ (strict inequality) for some component $D$ of the set $\mathcal{M}_{+}$, and $m(z) \geq 0$ for all $z \in \mathbb{T}_{P}^{2}$, then $m$ has a non-zero $L_{\rho}$-subminorant.

Exercise 6.2.2.2 Prove that $\mathcal{M}_{+}$is open.
6.2.3 If $\rho \notin \mathbb{Z}$, the operator $L_{\rho}$ has a fundamental solution $E_{\rho}(\bullet-\zeta)$ in $\mathbb{T}_{P}^{2}$, where $\zeta$ is a shift by the torus, i.e., by the modulus $P+i 2 \pi$. It means that

$$
L_{\rho} E_{\rho}(\bullet-\zeta)=\delta_{\zeta}
$$

in $\mathcal{D}^{\prime}\left(\mathbb{T}_{P}^{2}\right)$, where $\delta_{\zeta}$ is the Dirac function, concentrated at $\zeta$.
If $\rho \in \mathbb{Z}$, there exists, as for operator $T_{\rho}$ and a spherical operator (see Theorem 3.2.4.2, Theorem 3.2.6.3), a generalized fundamental solution $E_{\rho}^{\prime}$ that satisfies the equation

$$
L_{\rho} E_{\rho}^{\prime}(\bullet-\zeta)=\delta_{\zeta}-\cos \rho(y-\eta), \zeta=\xi+i \eta
$$

in $\mathcal{D}^{\prime}\left(\mathbb{T}_{P}^{2}\right)$.
Theorem 6.2.3.1 Let $\rho>0, \rho \notin \mathbb{Z}$. Then every $L_{\rho}$-subfunction on $\mathbb{T}_{P}^{2}$ can be represented in the form

$$
\begin{equation*}
q(z)=\int_{\mathbb{T}_{P}^{2}} E_{\rho}(z-\zeta) \nu(d \zeta) \tag{6.2.3.1}
\end{equation*}
$$

where $\nu=L_{\rho} q$.
This theorem is the counterpart of Theorems 3.2.3.3, 3.2.6.2.
Theorem 6.2.3.2 Let $\rho>0, \rho \in \mathbb{Z}$. Then the mass distribution $\nu=L_{\rho} v$ satisfies the condition

$$
\begin{equation*}
\int_{\mathbb{T}_{P}^{2}} e^{ \pm i \rho y} \nu(d z)=0 \tag{6.2.3.2}
\end{equation*}
$$

and the representation

$$
\begin{equation*}
q(z)=\Re\left(C e^{i \rho y}\right)+\int_{\mathbb{T}_{P}^{2}} E_{\rho}^{\prime}(z-\zeta) \nu(d \zeta) \tag{6.2.3.3}
\end{equation*}
$$

holds with $C$ that is a complex scalar.
This theorem is the counterpart of Theorems 3.2.4.2, 3.2.6.2.

Let $D \subset \mathbb{T}_{P}^{2}$ and $\rho(D)>\rho$. Then the operator $L_{\rho}$ has in $D$ the Green function $-G_{\rho}(z, \zeta, D)$. Thus for every $q$ that is an $L_{\rho}$-subfunction in $D$ and bounded from above in $\bar{D}$, we have the representation

$$
\begin{equation*}
q(z)=g(z)-\int_{D} G_{\rho}(z, \zeta, D) \nu(d \zeta) \tag{6.2.3.4}
\end{equation*}
$$

in which $\nu=L_{\rho} q$ and $g$ is the minimal majorant on $\partial D$ of the function $q$, satisfying $L_{\rho} g=0$ in $D$.

This is the counterpart of Theorem 2.6.4.3 (F. Riesz representation) and Theorem 3.2.5.1.

From (6.2.3.4) one can easily obtain
Theorem 6.2.3.3. (Maximum principle) If $\rho(D)>\rho$ and $q(z)$ is an $L_{\rho}$-subfunction such that $q(z) \leq 0, z \in \partial D$, then $q(z) \leq 0, z \in D$.

## Exercise 6.2.3.1 Prove this.

Theorem 6.2.3.4 An $L_{\rho}$-subfunction in $\mathbb{T}_{P}^{2}$ can not attain zero maximum if is not zero identically.

Exercise 6.2.3.1 Prove this by exploiting (6.2.1.2) and properties of subharmonic functions.

Theorem 6.2.3.5 Let $q$ be an $L_{\rho}$-subfunction in $\mathbb{T}_{P}^{2}$. If $q(z) \leq 0$ for $z \in \mathbb{T}_{P}^{2}$ then $q(z) \equiv 0$.

Exercise 6.2.3.2 Prove this using Theorem 3.1.4.7 (**Liouville).
Proposition 6.2.3.6 Let $q_{D}$ be the solution of the problem (6.2.2.1) in a domain $D$ with a smooth boundary, corresponding to $\rho=\rho(D)$. Suppose that $q_{D}\left(z_{0}\right)=1$ for some $z_{0} \in D$. Then

$$
\frac{\partial q_{D}}{\partial n}>0, \forall z \in \partial D
$$

Exercise 6.2.3.3 Prove this, using properties of positive harmonic functions.
6.2.4 In the part devoted to completeness of an exponential system (Section 6.3) we will need the notion of minimality of a subharmonic function from $U[\rho]$. A function $v \in U[\rho]$ is called minimal if the function $v-\epsilon r^{\rho}$ has no subharmonic minorant for arbitrarily small $\epsilon>0$. If $v$ is $T$-automorphic, the corresponding $L_{\rho^{-}}$ subfunction $q$ is called minimal if the function $q-\epsilon$ has no $L_{\rho}$-subminorant in $\mathbb{T}_{P}^{2}$. We formulate one sufficient condition for minimality and one sufficient condition for nonminimality.

Theorem 6.2.4.1 Let $\mathcal{H}_{\rho}(q)$ be the maximal open set on which $L_{\rho} q=0$. If there exists a connected component $M \subset \mathcal{H}_{\rho}(q)$ such that $\rho(M)<\rho$, then $q$ is a minimal $L_{\rho}$-subfunction.

For example, $q \equiv 0$ is minimal.
Proposition 6.2.4.2 The function $q$ is nonminimal if $q(z) \geq c$ or $L_{\rho} q-c>0$ for some positive $c$ for all $z \in \mathbb{T}_{P}^{2}$.

For example, $q \equiv c>0$ is nonminimal.

### 6.3 Completeness of exponential systems in convex domains

6.3.1 Let $\Lambda:=\left\{\lambda_{k}\right\}, k=1,2, \ldots$ be a set of points in the complex plane $\mathbb{C}$, satisfying the condition $\lambda_{k} \neq 0$ and $\lambda_{j} \neq \lambda_{k}$, if $k \neq j$.

Consider the canonical product

$$
\begin{equation*}
\Phi_{\Lambda}(\lambda):=\prod_{k}\left(1-\lambda / \lambda_{k}\right) \exp \lambda / \lambda_{k} \tag{6.3.1.1}
\end{equation*}
$$

We suppose in this section that $\Phi_{\Lambda}(\lambda)$ is an entire function of order 1 and normal type, i.e., a function of exponential type (see [Le, Ch. 1, § 20].

This fact can be expressed in terms of $\Lambda$ by using the Brelot-Lindelöf Theorem 2.9.4.2.

Exercise 6.3.1.1 Formulate this theorem for entire functions of order 1 and normal type under assumption that $\rho(r) \equiv 1$.

We will suppose that the upper density of zeros (see Section 2.8, Section 5.1). $\bar{\Delta}_{\Lambda}>0$.
6.3.2 Let $G \subset \mathbb{C}$ be a convex bounded domain containing zero. This last requirement does not restrict any of the further considerations connected to completeness, because $\exp \Lambda:=\left\{e^{\lambda_{j} z}: \lambda_{j} \in \Lambda\right\}$ can be replaced by the system $\left\{e^{\lambda_{j}\left(z-z_{0}\right)}: \lambda_{j} \in \Lambda\right\}$ and $e^{\lambda_{j}\left(z-z_{0}\right)}=C_{j} e^{\lambda_{j} z}$. Let $A(G)$ be the space of holomorphic functions in $G$ with the topology of uniform convergence on compact sets. We will study the completeness of the exponential systems

$$
\begin{equation*}
\exp \Lambda:=\left\{e^{\lambda_{j} z}: \lambda_{j} \in \Lambda\right\} \tag{6.3.2.1}
\end{equation*}
$$

in $A(G)$.
We will be interested in the following questions:

1. completeness of $\exp \Lambda$ in $A(G)$;
2. maximality of $G$ for $\exp \Lambda$, which is complete in $A(G)$;
3. extremal overcompleteness of $\exp \Lambda$ in $A(G)$ for a maximal $G$.

Let us give precise definitions of maximality and extremal overcompleteness. The completeness means that every function $f \in A(G)$ can be approximated on every compact set $K \Subset G$ with arbitrary precision by linear combinations of functions from $\exp \Lambda$.

A convex domain $G$ is called maximal for a system $\exp \Lambda$, which is complete in $A(G)$ if for every domain $G_{1}$ such that $G \Subset G_{1} \exp \Lambda$ is not complete in $A\left(G_{1}\right)$.

A system $\exp \Lambda$ is called extremely overcomplete in $A(G)$ for a maximal $G$, if for every sequence $\Lambda_{1}:=\left\{\lambda_{j}^{1}\right\}$ such that $\Lambda_{1} \cap \Lambda=\varnothing$ and $\bar{\Delta}_{\Lambda_{1}}>0$ the domain $G$ is not maximal for the system $\exp \Lambda \cup \Lambda_{1}$.

In other words, every essential enlargement of an extremely overcomplete system enlarges also the maximal domain of completeness.
6.3.3 Let

$$
h_{\Lambda}(\phi):=\limsup _{r \rightarrow \infty} \log \left|\Phi_{\Lambda}\left(r e^{i \phi}\right)\right| r^{-1}
$$

be the indicator of $\Phi_{\Lambda}$. It is a 1-trigonometrically convex function or simply a trigonometrically convex function (t.c.f). Let $G_{\Lambda}$ be the conjugate indicator diagram of $\Phi_{\Lambda}$, i.e., a convex domain of the form

$$
G_{\Lambda}:=\left\{z: \max _{z \in G_{\Lambda}} \Re\left(z e^{i \phi}\right) \leq h_{\Lambda}(\phi)\right\} .
$$

Let us describe conditions for completeness, maximality and extremal overcompleteness when $\Lambda$ is a regular set (see Section 5.6) and $\Phi_{\Lambda}$ is a CRG-function (see Section 5.6).

We say that $G_{\Lambda}$ is enclosed in $G$ if it can be enclosed in $G$ by parallel translation, enclosed with sliding, if it can be moved after enclosing only in one direction, enclosed rigidly if it is impossible to move after enclosing, freely enclosed in every other case of enclosing.
Theorem 6.3.3.1 Let $\Lambda$ be a regular set. Then the following holds:

1. $\{\exp \Lambda$ is not complete in $A(G)\} \Longleftrightarrow\left\{G_{\Lambda}\right.$ is freely enclosed in $\left.G\right\}$;
2. $\{G$ is maximal for $\exp \Lambda\} \Longleftrightarrow\left\{G_{\Lambda}\right.$ is not freely enclosed in $\left.G\right\}$;
3. $\{\exp \Lambda$ is extremely overcomplete in $A(G)\} \Longleftrightarrow\left\{G_{\Lambda}\right.$ is enclosed rigidly in $G\}$.

Let us note that $G$ is maximal for $\exp \Lambda$ but not extremely overcomplete if and only if $G_{\Lambda}$ is enclosed with sliding in $G$.

This theorem is a corollary of the more general Theorem 6.3.4.1, but will be proved independently in Section 6.3.10.
6.3.4 If $\Lambda$ is not regular, it is natural to exploit the notion of a limit set (see Section 3.1) to characterize $\exp \Lambda$.

Suppose the limit set of $\Phi_{\Lambda}$ has the form

$$
\operatorname{Fr}\left[\Phi_{\Lambda}\right]:=\left\{v(\lambda)=|\lambda|\left(c h_{1}+(1-c) h_{2}\right)(\arg \lambda): c \in[0 ; 1]\right\}
$$

where $h_{1}, h_{2}$ are t.c.f.

Such a limit set is a particular case of $U_{\text {ind }}$ (6.1.2.2). It is called an indicator limit set and it is indeed a limit set of an entire function (see Exercise 6.1.2.4).

The asymptotic behavior of the set $\Lambda$ (i.e., the limit set of the corresponding mass distribution) can be described completely using Theorem 3.1.5.2.

Exercise 6.3.4.2 Do that.
We will call such $\Lambda$ an indicator set. Denote by $G_{1}, G_{2}$ the conjugate diagram of $h_{1}, h_{2}$. Since $G_{1}, G_{2}$ are convex, the set

$$
\alpha G_{1}+\beta G_{2}:=\left\{\alpha z_{1}+\beta z_{2}: z_{1} \in G_{1}, z_{2} \in G_{2}\right\}, \alpha, \beta>0
$$

is also convex and is a conjugate diagram of the t.c.f. $h:=\alpha h_{1}+\beta h_{2}$.
Theorem 6.3.4.1 Let a set $\Lambda$ be an indicator set. Then the following holds:

1. $\{\exp \Lambda$ is not complete in $A(G)\} \Longleftrightarrow\left\{G_{1}\right.$ and $G_{2}$ are freely enclosed in $\left.G\right\}$;
2. $\{G$ is maximal for $\exp \Lambda\} \Longleftrightarrow\left\{G_{1}\right.$ and $G_{2}$ are enclosed in $G$ and at least one of them is not freely enclosed in $G\}$;
3. $\{\exp \Lambda$ is extremely overcomplete in $A(G)\} \Longleftrightarrow\left\{c G_{1}+(1-c) G_{2}\right.$ is enclosed rigidly in $G \forall c \in[0 ; 1]\}$.

This theorem is proved in Section 6.3.10.
The equality holds:

$$
\begin{equation*}
h_{\Lambda}=\max \left(h_{1}, h_{2}\right) \tag{6.3.4.1}
\end{equation*}
$$

Thus the conjugate diagram $G_{\Lambda}$ of the function $h_{\Lambda}$ is the convex hull of $G_{1}$ and $G_{2}$.
Let us note that the indicator $h_{\Lambda}$ does not determine the completeness of the system $\exp \Lambda$ if $\Lambda$ is not a regular set, as the following example shows.

Example 6.3.4.1 Let

$$
\begin{aligned}
G_{1} & :=\{z=x+i y: x=1 ;-1 \leq y \leq 1\} \\
G_{2} & :=\{z=x+i y: x=-1 ;-1 \leq y \leq 1\} \\
\text { and } \quad G & =\{z:|z|<1+\epsilon\}
\end{aligned}
$$

with a small $\epsilon$.

Exercise 6.3.4.3 Prove that $G_{1}$ and $G_{2}$ are freely enclosed in $G$ and their convex hull is not enclosed.

Let $\Lambda$ be a set such that the interior of $G_{\Lambda}$ coincides with $G$. If $\Lambda$ is a regular set, then $\exp \Lambda$ is complete in $A(G), G$ is maximal for $\exp \Lambda$ and $\exp \Lambda$ is extremely overcomplete in $A(G)$.

If $\Lambda$ is an indicator set, then the first two assertions hold but $\exp \Lambda$ can be not extremely overcomplete:

Example 6.3.4.2 Set

$$
\begin{aligned}
G_{1} & :=\{z=x+i y:-1 \leq x \leq 0 ; y=0\} \\
G_{2} & :=\{z=x+i y: x=1 ;-1 \leq y \leq 1\}
\end{aligned}
$$

Here $G_{\Lambda}$ is a triangle in which $G_{1}$ is freely enclosed and $G_{2}$ is rigidly enclosed, but $c G_{1}+(1-c) G_{2}$ is free enclosed for all $c: 0<c<1$.

Exercise 6.3.4.4 Check this.
Example 6.3.4.3 Set

$$
\begin{aligned}
G_{1} & :=\{z=x+i y: x=-1 ; y \in[-1 ; 1]\} \\
G_{2} & :=\{z=x+i y: x=1 ; y \in[-1 ; 1]\} .
\end{aligned}
$$

Exercise 6.3.4.5 Check that $G_{1}$ and $G_{2}$ are enclosed with sliding in $G_{\Lambda}$.
If $G_{1}$ and $G_{2}$ are rigidly enclosed in $G$ it does not imply in general that $c G_{1}+(1-c) G_{2}$ are rigidly enclosed for all $c \in[0 ; 1]$.

Example 6.3.4.4 Let $G_{1}$ be an equilateral triangle inscribed in the circle $|z|=1$, let $G_{2}$ be the same triangle rotated by the angle $\pi / 6$, and let $G$ be the unit disc.

Exercise 6.3.4.6 Show that $\frac{1}{2}\left(G_{1}+G_{2}\right)$ is freely enclosed in $G$.
If $G_{1}, G_{2} \subset G$ and $G_{1} \cap G_{2}$ is rigidly enclosed in $G$, then $c G_{1}+(1-c) G_{2}$ is rigidly enclosed for $c \in[0 ; 1]$.

Exercise 6.3.4.7 Check this.
However this is not a necessary condition.

## Example 6.3.4.5

$$
\begin{aligned}
G & :=\{z=x+i y:|x|<1 ;|y|<1\} ; \\
G_{1} & :=\{z=x+i y: x \in(-1,1) ;-x>y>-1\} ; \\
G_{2} & :=\{z=x+i y: x \in(-1,1) ;-1<y<x\} .
\end{aligned}
$$

Exercise 6.3.4.8 Check that every triangle $c G_{1}+(1-c) G_{2}$ is rigidly enclosed in $G$ and $G_{1} \cap G_{2}$ is freely enclosed.
6.3.5 Consider in more detail the conditions for extremal overcompleteness in the case when $\Lambda$ is an indicator set and $G_{\Lambda}=G$ or, in other words, if

$$
\begin{equation*}
h_{\Lambda}=h_{G} . \tag{6.3.5.1}
\end{equation*}
$$

We can suppose that $h_{1}$ and $h_{2}$ are linearly independent, otherwise we exploit Theorem 6.3.3.1. If, for example, the inequality $h_{1}(\phi) \leq h_{2}(\phi), \forall \phi$, holds, the extremal overcompleteness is in the case when $G_{1}$ is rigidly enclosed in $G_{2}$ because $G_{1} \cap G_{2}=G_{1}$, and this case was mentioned above (Exercise 6.3.4.7).

Consider the general case. Denote $g(\phi):=\left|h_{1}-h_{2}\right|(\phi)$, and set $\Theta_{\Lambda}:=$ $\{\phi: g(\phi)>0\}$. This is an open set on the unit circle. Denote as $I_{\Lambda}:=\left(\alpha_{1}, \alpha_{2}\right)$ the maximal interval contained in $\Theta_{\Lambda}$ and denote by $d_{\Lambda}$ its length. Since $g(\phi)$ is continuous,

$$
\begin{equation*}
g\left(\alpha_{j}\right)=0, j=1,2 . \tag{6.3.5.2}
\end{equation*}
$$

If also at least one of the conditions

$$
\liminf _{\phi \in I_{\Lambda}, \phi \rightarrow \alpha_{j}} \frac{g(\phi)}{\phi-\alpha_{j}}=0, j=1,2,
$$

is fulfilled, we say $g$ is zero with tangency on $\partial I_{\Lambda}$.
Theorem 6.3.5.1 Suppose $\Lambda$ is an indicator set that satisfies (6.3.5.1). In order that $\exp \Lambda$ be extremely overcomplete in $A(G)$ it is necessary and sufficient that at least one of the following conditions holds:

1. $d_{\Lambda}<\pi$;
2. $d_{\Lambda}=\pi$ and $g$ is zero with tangency on $\partial I_{\Lambda}$.

This theorem is proved in Section 6.3.11.
6.3.6 We call $\Lambda$ periodic if $\operatorname{Fr}\left[\Phi_{\Lambda}\right]$ is a periodic limit set (see Theorem 4.1.7.1). In such a case all the limit set is determined by one subharmonic function $v \in U[1]$ (see (4.1.3.1)). Let us characterize the system $\exp \Lambda$ for periodic $\Lambda$.

Set

$$
\begin{align*}
h_{G}(\phi) & :=\max \left\{\Re\left(z e^{i \phi}\right): z \in G\right\},  \tag{6.3.6.1}\\
m(\lambda, G, v) & :=|\lambda| h_{G}(\arg \lambda)-v(\lambda) . \tag{6.3.6.2}
\end{align*}
$$

Denote by $\mathcal{G}_{G} v$ the maximal subharmonic minorant of the function $m(\lambda, G, v)$. A function $w \in U[1]$ is called minimal if the function $w-\epsilon|\lambda|$ has no subharmonic minorant in $U[1]$ for every small $\epsilon>0$. The harmonic function of the form

$$
\begin{equation*}
H(\lambda):=|\lambda|(A \cos (\arg \lambda)+B \sin (\arg \lambda)), \tag{6.3.6.3}
\end{equation*}
$$

for example, is minimal.
We will denote as HARM the set of functions of the form (6.3.6.3).
Theorem 6.3.6.1 Let $\Lambda$ be a periodic set. The following holds:

1. $\{\exp \Lambda$ is not complete in $A(G)\} \Longleftrightarrow\left\{\mathcal{G}_{G} v\right.$ exists and is non-minimal\};
2. $\{G$ is maximal for $\exp \Lambda\} \Longleftrightarrow\left\{\mathcal{G}_{G} v\right.$ exists and is minimal $\}$;
3. $\{\exp \Lambda$ is extremely overcomplete in $A(G)\} \Longleftrightarrow\left\{\mathcal{G}_{G} v \in\right.$ HARM $\}$.

This theorem is proved in Section 6.3.12.
6.3.7 Let us characterize the completeness of $\exp _{\Lambda}$ for periodic $\Lambda$ in other terms. For this we need the information that was presented in Section 6.2. We will take $\rho=1$. Denote

$$
\begin{equation*}
q_{\Lambda}(z):=v_{\Lambda}\left(e^{z}\right) e^{-x} \tag{6.3.7.1}
\end{equation*}
$$

(compare with (6.2.1.2)). As it was explained in Section 6.2 this function is an $L_{1}$-subfunction on the torus $\mathbb{T}_{P}^{2}$. Set

$$
m\left(z, G, q_{\Lambda}\right)=h_{G}(y)-q_{\Lambda}, D(G, \Lambda):=\left\{z: m\left(z, G, q_{\Lambda}\right)>0\right\} \subset \mathbb{T}_{P}^{2}
$$

The set $D(G, \Lambda)$ is open because $-m$ is an upper semicontinuous function (see Theorem 2.1.2.4), denote

$$
\rho(\Lambda, G):=\min \rho(M)
$$

where the minimum is taken over all components $M$ of $D(G, \Lambda)$, and it is attained on one of the components because they are not intersecting and $\mathbb{T}_{P}^{2}$ is compact.

Exercise 6.3.7.1 Explain this in detail, using properties of $\rho(D)$ (Section 6.2).
Theorem 6.3.7.1 If

$$
\begin{equation*}
\rho(\Lambda, G) \geq 1 \tag{6.3.7.2}
\end{equation*}
$$

then $\exp \Lambda$ is complete in $G$.
This theorem is proved in Section 6.3.12.
Let $w:=g_{G} q_{\Lambda}(z)$ be the maximal $L_{1}$-subminorant of $m\left(z, G, q_{\Lambda}\right)$. Denote by $\mathcal{H}_{\Lambda}$ the open set in $\mathbb{T}_{P}^{2}$ where $L_{1} w=0$.

Theorem 6.3.7.2 If there exists a component $M$ of $\mathcal{H}_{\Lambda}$ such that $\rho(M)<1$, then $w$ is minimal, and, hence, $G$ is maximal for $\exp \Lambda$.

This theorem follows directly from Theorem 6.2.4.1. It is not known if the condition (6.3.7.2) is necessary.

Consider in detail the situation in which the domain $G$ coincides with $G_{\Lambda}$, the conjugated indicator diagram of $h_{\Lambda}$, i.e., we suppose that

$$
\begin{equation*}
h_{G}(\phi)=h_{\Lambda}(\phi), \forall \phi . \tag{6.3.7.3}
\end{equation*}
$$

In this case $m\left(z, G, q_{\Lambda}\right) \geq 0$ and we obtain the following criterion:
Theorem 6.3.7.3 In order that $\exp \Lambda$ be complete in $G_{\Lambda}$ it is necessary and sufficient that

$$
\begin{equation*}
\rho\left(\Lambda, G_{\Lambda}\right) \geq 1 \tag{6.3.7.4}
\end{equation*}
$$

This theorem is proved in Section 6.3.12. The condition (6.3.7.3) automatically implies maximality if there is completeness.

Since

$$
\begin{equation*}
h_{\Lambda}(y)=\max \left\{q_{\Lambda}(x+i y): x \in[0 ; P]\right\} \tag{6.3.7.5}
\end{equation*}
$$

the function $m\left(z, G_{\Lambda}, q_{\Lambda}\right)$ has a zero in $x$ for every fixed $y$.
Thus the set $D(G, \Lambda)$ does not contain any curve $y=$ const on the torus.

Theorem 6.3.7.4 Let $G_{0}$ be a strictly convex domain and let $D_{0} \subset \mathbb{T}_{P}^{2}$ be such that $\mathbb{T}_{P}^{2} \backslash D_{0}$ intersect every line $\left\{y=y_{0}\right\}, y_{0} \in[0,2 \pi]$.

Then there exists a periodic $\Lambda$ such that

$$
\begin{equation*}
G_{\Lambda}=G_{0}, D\left(G_{\Lambda}, \Lambda\right)=D_{0} \tag{6.3.7.6}
\end{equation*}
$$

This theorem is proved in Section 6.3.13.
Example 6.3.7.1 Let $D_{0}$ be the complement in $\mathbb{T}_{P}^{2}$ to the set

$$
\begin{equation*}
M:=\{z=x+i y: x=f(y), y \in[0 ; 2 \pi]\} \tag{6.3.7.7}
\end{equation*}
$$

where $f(y)$ is a continuous $2 \pi$-periodic function satisfying the condition

$$
0<f(y)<P
$$

Then $\rho\left(D_{0}\right)=\infty$, because this domain is not connected on spirals (see Section 6.2.). It means that for every strictly convex $G_{0}$ there exists a periodic $\Lambda$ such that $G_{\Lambda}=G_{0}$ and $\exp \Lambda$ is extremely overcomplete in $G_{0}$.

Example 6.3.7.2 Let $D_{0}$ be the complement to the set

$$
M:=\left\{z=x+i y: x=\frac{P}{2 \pi} y, 0 \leq y \leq 2 \pi\right\} .
$$

Then

$$
\begin{equation*}
\rho\left(D_{0}\right)=\frac{1}{2}\left(1+(2 \pi / P)^{2}\right) \tag{6.3.7.8}
\end{equation*}
$$

(see Section 6.3.13).
Thus, choosing $P$, and using Theorem 6.3.7.4, it is possible make exp $\Lambda$ complete or non-complete in $G_{0}\left(=G_{\Lambda}\right)$ for every strictly convex domain $G_{0}$.
6.3.8 Now pass to generalizations. Denote by $D_{G}$ the natural domain of definition of the operation $\mathcal{G}_{G}$, i.e., the set of $v \in U[1]$ for which $m(\lambda, G, v)$ (see (6.3.6.2)) has a subharmonic minorant belonging to $U[1]$.

Let $\Phi_{\Lambda}$ be defined by the equality (6.3.1.1). The condition that for every $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ the function $m(\lambda, G, v)$ has a subharmonic minorant belonging to $U[1]$ is possible to express by the relation

$$
\begin{equation*}
\operatorname{Fr}\left[\Phi_{\Lambda}\right] \subset D_{G} \tag{6.3.8.1}
\end{equation*}
$$

(compare with (6.1.1.1)).
We call the set $U \subset U[1]$ minimal $(U \in \operatorname{MIN})$ if for arbitrarily small $\epsilon>0$ there exists $w=w_{\epsilon} \in U$ such that the function $w_{\epsilon}-\epsilon|\lambda|$ has no subharmonic minorant, belonging to $U[1]$.

Let us note that if $U$ contains a minimal function (in the sense of Section 6.3.6), then $U \in$ MIN. We denote the image of $\operatorname{Fr}\left[\Phi_{\Lambda}\right]$ under the mapping by the operator $\mathcal{G}_{G}$ as $J_{G}(\Lambda)$.

Theorem 6.3.8.1 The following holds:

1. $\left\{\exp _{\Lambda}\right.$ is not complete in $\left.A(G)\right\} \Longleftrightarrow\left\{((6.3 .8 .1)\right.$ holds $) \wedge\left(J_{G}(\Lambda) \notin\right.$ MIN $\left.)\right\}$;
2. $\left\{G\right.$ is maximal for $\left.\exp _{\Lambda}\right\} \Longleftrightarrow\left\{((6.3 .8 .1)\right.$ holds $\left.) \wedge\left(J_{G}(\Lambda) \in \mathrm{MIN}\right)\right\}$;
3. $\{\exp \Lambda$ is extremely overcomplete for maximal $G\} \Longleftrightarrow\{((6.3 .8 .1)$ holds $) \wedge$ $\left.\left(J_{G}(\Lambda) \in \mathrm{HARM}\right)\right\} ;$
6.3.9 In the proof of Theorem 6.3.8.1 that we are going to prove now we exploit

Theorem 6.3.9.1. (A.I. Markushevich) see [Le, Ch. 4, § 7] Let $A(\mathbb{C} \backslash \bar{G})$ be a class of functions $\psi$ which are holomorphic in $\mathbb{C} \backslash \bar{G}$ and equal to zero in infinity. In order that the system $\exp \Lambda$ be complete in $A(G)$, it is necessary and sufficient that the function

$$
\begin{equation*}
\Phi(\lambda):=\int_{L_{\psi}} e^{\lambda z} \psi(z) d z, \tag{6.3.9.1}
\end{equation*}
$$

where $\psi \in A(\mathbb{C} \backslash \bar{G})$, and $L_{\psi} \Subset G$ is a rectifiable closed curve, has the following property: the condition

$$
\begin{equation*}
\Phi\left(\lambda_{k}\right)=0, \forall \lambda_{k} \in \Lambda \tag{6.3.9.2}
\end{equation*}
$$

implies $\Phi(\lambda) \equiv 0$.
Proof of Theorem 6.3.8.1, 1. Necessity. Let $\exp \Lambda$ be not complete. By Theorem 6.3.9.1 $\Phi\left(\lambda_{k}\right)=0$, but $\Phi(\lambda) \not \equiv 0$. The function $g(\lambda):=\Phi(\lambda) / \Phi_{\Lambda}(\lambda)$, where $\Phi_{\Lambda}$ is from (6.3.1.1), is an entire function and it has order one and normal or minimal type by Theorem 2.9.3.1. Set

$$
u^{g}:=\log |g(\lambda)| ; u^{\Phi}(\lambda):=\log |\Phi(\lambda)| ; u^{\Lambda}(\lambda):=\log |\Phi(\lambda)| .
$$

We have from (6.3.9.1) $\left.u^{\Phi}(\lambda) \leq \max \left\{\Re(\lambda z): z \in L_{\psi}\right\}+C_{\psi}\right\}$, where $C_{\psi}$ is a constant, depending possibly on $\psi$.

This implies that

$$
\begin{equation*}
u^{\Phi}(\lambda) \leq h_{G_{1}}(\phi) r+C_{\psi}, \lambda=r e^{i \phi} \tag{6.3.9.3}
\end{equation*}
$$

for some convex domain $G_{1} \Subset G$.
Let $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$. Choose a sequence $t_{j} \rightarrow \infty$ for which $\left(u^{\Lambda}\right)_{t_{j}} \rightarrow v$, and the sequences $\left(u^{\Phi}\right)_{t_{j}}$ and $\left(u^{g}\right)_{t_{j}}$ also converge to $v^{\Phi}$ and $v^{g}$ respectively. From the equality $u^{g}(\lambda)=u^{\Phi}(\lambda)-u^{\Lambda}(\lambda)$ we obtain $v^{g}(\lambda)=v^{\Phi}(\lambda)-v(\lambda)$ where $v^{g} \in$ $\mathbf{F r}[g], v^{\Phi} \in \mathbf{F r}[\Phi]$.

Since (6.3.9.3) implies $v^{\Phi}(\lambda) \leq h_{G_{1}}(\phi) r$,

$$
\begin{equation*}
v^{g}(\lambda) \leq h_{G_{1}}(\phi) r-v(\lambda) \tag{6.3.9.4}
\end{equation*}
$$

and it means that for every $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right] \mathcal{G}_{G_{1}} v$ and hence $\mathcal{G}_{G} v$ exist, i.e., the condition (6.3.8.1) holds.

Let us show that the condition $J_{G}(\Lambda) \notin \mathrm{MIN}$ is satisfied. We have for some $\delta>0$ the relation

$$
h_{G_{1}}(\phi)-h_{G}(\phi) \leq-\delta .
$$

From (6.3.9.4) we obtain

$$
\begin{equation*}
v^{g}(\lambda)+\delta r \leq m(\lambda, G, v) \tag{6.3.9.5}
\end{equation*}
$$

The left-hand side of the inequality (6.3.9.5) belongs to $U[1]$. Thus $w_{v}:=\mathcal{G}_{G_{1}} v$ satisfies the condition $v^{g}(\lambda)+\delta r \leq w_{v}(\lambda)$ for every $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$. It means that $J_{G}(\Lambda) \notin \mathrm{MIN}$.

Necessity is proved.
For proving sufficiency we exploit the following assertion.
Theorem 6.3.9.2 (I.F. Krasichkov-Ternovskii) Suppose there exists an entire function $g$ such that

$$
\begin{equation*}
h_{g \Phi_{\Lambda}}(\phi)<h_{G}(\phi), \forall \phi . \tag{6.3.9.6}
\end{equation*}
$$

Then the system $\exp \Lambda$ is not complete for some convex domain $G_{1} \Subset G$.
This theorem connects the problem of completeness to the multiplicator problem.

Proof of Theorem 6.3.9.2. Let $g(\lambda)$ satisfy (6.3.9.6). Denote by $\psi(z)$ the Borel transformation for $\Phi(\lambda):=g(\lambda) \Phi_{\Lambda}(\lambda)$. By the Pólya Theorem (see, for example, [Le, Ch. 1, §20]), all the singularities of $\psi$ are contained in a convex domain $G_{\Phi}$ which is the conjugate diagram of the indicator $h_{\Phi}(\phi)$. Thus the representation (6.3.9.1) holds with $L_{\psi}$ that embraces $G_{\Phi}$. It follows from (6.3.9.6) that $G_{\Phi} \Subset G$. Thus it is possible to choose $L_{\psi}$ between $\partial G_{\Phi}$ and $\partial G$. Since (6.3.9.2) for $\Phi$ is fulfilled and $\Phi(\lambda) \not \equiv 0, \exp \Lambda$ is non-complete in some convex $G_{1} \Subset G$ such that $L_{\psi} \Subset G_{1}$ by Theorem 6.3.9.1.

Now we can prove sufficiency in Theorem 6.3.8.1, 1. From the condition $J_{G}(\Lambda) \notin$ MIN it follows that one can choose $\delta>0$ such that $\forall v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ the functions $w_{v}-\delta r$ where $w_{v}:=\mathcal{G}_{G}$, have subharmonic minorants. As we already said in Section 6.3.2, completeness does not depend on shift by any fixed $z_{0}$. Thus we can suppose that $0 \in G$ and, hence, $h_{G}(\phi)>0$ for all $\phi$. Let $\gamma<2 \delta$ be such that $h_{G}(\phi)-\gamma>0$ and $G_{1} \Subset G$ satisfy

$$
\begin{equation*}
h_{G_{1}}(\phi)-\gamma / 3>0, h_{G}(\phi)-h_{G_{1}}(\phi) \leq \gamma / 2 . \tag{6.3.9.7}
\end{equation*}
$$

Let us check that

$$
\begin{equation*}
D_{G_{1}} \supset \operatorname{Fr}\left[\Phi_{\Lambda}\right], \tag{6.3.9.8}
\end{equation*}
$$

Indeed, for $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ we have

$$
\begin{align*}
m\left(\lambda, G_{1}, v\right):=h_{G_{1}}(\phi) r-v(\lambda) & \geq h_{G}(\phi) r-(\gamma / 2) r-v(\lambda) \\
& \geq h_{G}(\phi) r-v(\lambda)-\delta r \\
& \geq w_{v}-\delta r . \tag{6.3.9.9}
\end{align*}
$$

Since the right-hand side of (6.3.9.9) has a subharmonic minorant from $U[1]$, then (6.3.9.8) is proved. By Theorem 6.1.1.1 there exists a multiplicator $g(z) \in A(1)$ such that

$$
\begin{equation*}
h_{g \Phi}(\phi) \leq h_{G_{1}}(\phi)<h_{G}(\phi) . \tag{6.3.9.10}
\end{equation*}
$$

From Theorem 6.3.9.2 we obtain that $\exp \Lambda$ is non-complete in $G$.
Proof of Theorem 6.3.8.1, 2. Necessity. Let $G_{j}, j=1,2, \ldots$ be a sequence of convex domains, satisfying the conditions $G_{j} \supseteq G, G_{j} \downarrow G$. Since $\exp \Lambda$ is noncomplete in every $A\left(G_{j}\right), D_{G_{j}} \supset \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ by Theorem 6.3.8.1, 1 .

The sequence $w_{j}:=\mathcal{G}_{G_{j}} v$ satisfies

$$
w_{j}(\lambda) \leq h_{G_{j}}(\phi) r-v(\lambda), \lambda \in \mathbb{C} .
$$

Since $\left\{w_{j}\right\}$ is compact and $h_{G_{j}} \rightarrow h_{G}$, one can find a subsequence with the limit $w \in U[1]$. Then $w(\lambda) \leq h_{G}(\phi) r-v(\lambda)$. Hence $\mathcal{G}_{G} v$ exists.

If $J_{G} \in$ MIN would not hold, then, by Theorem 6.3.8.1, $1, \exp \Lambda$ is noncomplete in $A(G)$, which contradicts maximality.

Necessity is proved. Let us prove sufficiency.
Completeness of $\exp \Lambda$ in $A(G)$ follows from Theorem 6.3.8.1, 1. We will prove that $\exp _{\Lambda}$ is non-complete in $A\left(G_{1}\right)$ for every $G_{1} \ni G$ under the condition $D_{G} \supset \operatorname{Fr}\left[\Phi_{\Lambda}\right]$. Set

$$
\delta:=\min _{\phi}\left[h_{G_{1}}(\phi)-h_{G}(\phi)\right]>0 .
$$

Then $\forall v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$,

$$
\mathcal{G}_{G} v+\delta r \leq h_{G_{1}}(\phi) r-v(\lambda), \lambda \in \mathbb{C} .
$$

This means that $\mathcal{G}_{G_{1}} v \geq \mathcal{G}_{G} v+\delta r$. Hence $J_{G_{1}}(\Lambda) \notin$ MIN and, by Theorem 6.3.8.1, $1, \exp \Lambda$ is non-complete in $A\left(G_{1}\right)$.

Proof of Theorem 6.3.8.1, 3. Necessity. By Theorem 6.3.8.1, 2 from maximality $G$ (6.3.8.1) follows. We will prove that $\mathcal{G}_{G} v \in \operatorname{HARM} \forall v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$. Suppose it is not fulfilled, i.e., there exists $v_{0} \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ such that the mass distribution $\nu_{0}$ of the function $w_{0}=\mathcal{G}_{G} v_{0}$ is not zero. By Proposition 6.1.1.3 there exists a multiplicator $g$ such that $v_{0}+w_{0} \in \operatorname{Fr}\left[g \Phi_{\Lambda}\right]$. Let $\Lambda_{0}$ be the set of zeros of $g$. Since $\nu_{0} \in \operatorname{Fr} \Lambda_{0}$, $\bar{\Delta}\left(\Lambda_{0}\right)>0$, because $\nu_{0} \neq 0$ and by the definitions in Section 3.3.1.

We can shift a little zeros of $g$ and suppose without lack of generality that they are simple and $\Lambda_{0} \cap \Lambda=\varnothing$.

The condition for a multiplicator gives the inequality:

$$
h_{g \Phi_{\Lambda}}(\phi) \leq h_{G}(\phi), \forall \phi .
$$

It implies

$$
m\left(\lambda, G, v_{\Pi}\right)=r h_{G}(\phi)-v_{\Pi} \geq 0
$$

for all $v_{\Pi} \in \operatorname{Fr}\left[g \Phi_{\Lambda}\right]$. It means that $m\left(\lambda, G, v_{\Pi}\right)$ has zero as a minorant $\forall v_{\Pi} \in$ $\operatorname{Fr}\left[g \Phi_{\Lambda}\right]$, i.e., $D_{G} \supset \operatorname{Fr}\left[g \Phi_{\Lambda}\right]$. So the domain $G$ is maximal although the system $\exp \Lambda$ is replaced with the system $\exp \left(\Lambda \cup \Lambda_{0}\right)$. This contradicts the extremal overcompleteness. Hence, $\nu_{0} \equiv 0$ and $w_{0}=\mathcal{G}_{G} v_{0} \in$ HARM.

Necessity is proved. Let us prove sufficiency.
Let the condition $\mathcal{G}_{G} v \in \operatorname{HARM} \forall v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ hold. Suppose that there exists $\Lambda_{0}$ such that $\bar{\Delta}_{\Lambda_{0}}>0$ and $G$ is maximal for the system $\exp \left(\Lambda \cup \Lambda_{0}\right)$.

Theorem 6.3.8.1, 2 implies

$$
\begin{equation*}
D_{G} \supset \operatorname{Fr}\left[\Phi_{\Lambda_{1}}\right], \tag{6.3.9.11}
\end{equation*}
$$

where $\Lambda_{1}=\Lambda \cup \Lambda_{0}$.
For every $v_{0} \in \operatorname{Fr}\left[\Phi_{\Lambda_{0}}\right]$ one can find $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$ such that

$$
v_{1}:=v_{0}+v \in \operatorname{Fr}\left[\Phi_{\Lambda_{1}}\right] .
$$

The condition $\bar{\Delta}_{\Lambda_{0}}>0$ implies that one can choose $v_{0}$ for which the Riesz measure $\nu_{0} \not \equiv 0$. For $w_{1}=\mathcal{G}_{G} v_{1}$ one has the inequality $w_{1} \leq r h_{G}-v_{1}$ by (6.3.9.11), so $w_{1}+v_{0} \leq r h_{G}-v$ holds. Hence $w_{v}:=\mathcal{G}_{G} v$ satisfies the inequality

$$
\begin{equation*}
\left(w_{1}+v_{0}\right)(\lambda) \leq w_{v}(\lambda), \forall \lambda \in \mathbb{C} . \tag{6.3.9.12}
\end{equation*}
$$

Let us show that (6.3.9.12) is impossible.Indeed, since $w_{v} \in \operatorname{HARM} w:=w_{1}+$ $v_{0}-w_{v} \leq 0$ and $w \in U[\rho]$. Thus $w \equiv 0$. However the Riesz measure $\nu_{w} \geq \nu_{0} \not \equiv 0$, hence $w \not \equiv 0$. This contradiction proves sufficiency.
6.3.10 Now we prove Theorems 6.3.3.1, 6.3.4.1 and 6.3.5.1. We need some auxiliary assertions.

Lemma 6.3.10.1 Let $v:=r h_{1}(\phi)$ and $G_{1}$ be the conjugated diagram of $h_{1}$. Then the following holds:

1. $\left\{G_{1}\right.$ is freely enclosed in $\left.G\right\} \Longleftrightarrow\left\{\mathcal{G}_{G} v\right.$ is non-minimal $\}$;
2. $\left\{G_{1}\right.$ is enclosed to $G$ but not free enclosed $\} \Longleftrightarrow\left\{\mathcal{G}_{G} v\right.$ is minimal $\}$;
3. $\left\{G_{1}\right.$ is rigidly enclosed in $\left.G\right\} \Longleftrightarrow\left\{\mathcal{G}_{G} v \in \mathrm{HARM}\right\}$;
4. $\left\{G_{1}\right.$ is not enclosed in $\left.G\right\} \Longleftrightarrow\left\{\mathcal{G}_{G} v\right.$ does not exist $\}$.

To prove this lemma we need the following two lemmas.
Lemma 6.3.10.2 Let $v:=r h(\phi)$. Then $\mathcal{G}_{G} v=r h_{1}(\phi)$ where $h_{1}$ is the maximal trigonometrically convex minorant of the function

$$
m(\phi, G, h):=h_{G}(\phi)-h(\phi) .
$$

Proof. Let $v_{1}=\mathcal{G}_{G} v$. Since $v_{[t]}=v$ for all $t>0$,

$$
\left(v_{1}\right)_{[t]}=\mathcal{G}_{G} v_{[t]}=\mathcal{G}_{G} v
$$

by Theorem 6.1.1.2, 2 .
Thus the function

$$
\hat{v}_{1}:=\left(\sup _{t}\left(v_{1}\right)_{[t]}\right)^{*}(\lambda) \geq v_{1}(\lambda)
$$

and is also a subharmonic minorant belonging to $U[1]$. Thus $v_{1}=\hat{v}_{1}$. However, the function $\hat{v}_{1}$ is invariant with respect to the transformation $(\bullet)_{[t]}$. Hence it has the form $r h_{1}(\phi)$. The maximality of $h_{1}(\phi)$ follows from the maximality $v_{1}$.

Lemma 6.3.10.3 In order that $v:=r h_{1}$ be a minimal function, it is necessary and sufficient that $G_{1}$, the conjugate diagram of $h_{1}$, be a segment (in particular, a point).

Proof. Let $v=r h_{1}$ be minimal and let $G_{1}$ be the conjugate diagram of $h_{1}$. If $G_{1}$ is not the segment, then it contains some disc of radius $\delta>0$. Hence there exists a trigonometric function $A \cos \left(\phi-\phi_{0}\right)$ such that

$$
\delta+A \cos \left(\phi-\phi_{0}\right) \leq h_{1}(\phi) .
$$

Multiplying this inequality by $r$, we obtain that $v-\delta r$ has a harmonic (and hence subharmonic) minorant. This contradicts minimality.

Inversely, suppose $v$ is not minimal. Then there exists $\delta>0$ and t.c.f. $h_{2}(\phi)$ such that

$$
\begin{equation*}
h_{2}(\phi) \leq h_{1}(\phi)-\delta . \tag{6.3.10.1}
\end{equation*}
$$

For every t.c.f. $h_{2}$ there exists a trigonometric function $A \cos \left(\phi-\phi_{0}\right)$ such that

$$
\begin{equation*}
h_{2}(\phi)+A \cos \left(\phi-\phi_{0}\right) \geq 0 . \tag{6.3.10.2}
\end{equation*}
$$

This corresponds to a shift of the diagram which contains zero. From (6.3.10.1) and (6.3.10.2) we obtain

$$
\delta-A \cos \left(\phi-\phi_{0}\right) \leq h_{1}(\phi),
$$

which means that $G_{1}$ contains some disc of radius $\delta>0$. So it is not a segment.

Proof of Lemma 6.3.10.1. $G$ is freely enclosed iff the following assertion holds: there exists $\delta>0$ and a trigonometrical function $A \cos \left(\phi-\phi_{0}\right)$ such that the inequality

$$
\begin{equation*}
h_{1}(\phi)+\delta-A \cos \left(\phi-\phi_{0}\right) \leq h_{G}(\phi) \tag{6.3.10.3}
\end{equation*}
$$

holds.

Exercise 6.3.10.1 Prove this.
Let $\mathcal{G}_{G} v$ be non-minimal. By Lemma 6.3.10.2 it has the form $w_{2}=r h_{2}$, where $h_{2}$ is the maximal trigonometrically convex minorant of $m\left(\phi, G, h_{1}\right)$. There exists $\delta>0$ such that the function $w_{2}-\delta r$ has the maximal subharmonic minorant $v_{3}=r h_{3}(\phi)$. Let $A \cos \left(\phi-\phi_{0}\right)$ be a trigonometric function for which

$$
h_{3}(\phi)+A \cos \left(\phi-\phi_{0}\right) \geq 0 .
$$

In addition,

$$
h_{3}(\phi) \leq h_{2}(\phi)-\delta, h_{2}(\phi) \leq h_{G}-h_{1}(\phi)
$$

From this we obtain (6.3.10.3) and hence that $G_{1}$ is free enclosed.
Inversely, let $G_{1}$ be freely enclosed in $G$. From (6.3.10.3) it follows that

$$
\begin{equation*}
\delta-A \cos \left(\phi-\phi_{0}\right) \leq h_{G}(\phi)-h_{1}(\phi) \tag{6.3.10.4}
\end{equation*}
$$

Multiplying (6.3.10.4) by $r$, we obtain that $m(\lambda, G, v)$ has a minorant $v_{0}=r(\delta-$ $\left.A \cos \left(\phi-\phi_{0}\right)\right)$ which obviously is non-minimal. Hence, $\mathcal{G}_{G} v$ is non-minimal.
$G_{1}$ is enclosed in $G$ with sliding, hence there does not exist $\delta>0$ such that (6.3.10.3) is fulfilled, but there exists a segment with support function

$$
E(\phi)=B|\sin \phi|+A \cos \left(\phi-\phi_{0}\right),
$$

such that the inequality

$$
\begin{equation*}
h_{1}(\phi)+E(\phi) \leq h_{g}(\phi) \tag{6.3.10.5}
\end{equation*}
$$

holds.
Exercise 6.3.10.2 Prove this.
Let $\mathcal{G}_{G} v$ be minimal. By Lemma 6.3.10.2 it has the form $w_{2}=r h_{2}$ and by Lemma 6.3.10.3, $h_{2}=E(\phi)$. Thus $E(\phi) \leq\left(h_{G}-h_{1}\right)(\phi)$, which is equivalent to (6.3.10.5).

Prove 2, suppose $G$ is not freely enclosed and hence only (6.3.10.5) is possible. If $\mathcal{G}_{G} v$ were non-minimal, (6.3.10.3) would follow, as it was proved above. This contradicts the supposition.

The rigid enclosure is equivalent only to the inequality of the form

$$
h(\phi)-A \cos \left(\phi-\phi_{0}\right) \leq h_{G}(\phi) \forall \phi,
$$

and impossibility of enclosure is equivalent to the impossibility of even such an inequality. Thus all other assertions of the lemma can be proved analogously.
Exercise 6.3.10.3 Do this in detail.
Proof of Theorem 6.3.3.1. Regularity of $\Lambda$ means that $\operatorname{Fr}\left[\Phi_{\Lambda}\right]=\left\{v_{0}\right\}$ where $v_{0}=r h_{\Lambda}$. Thus $J_{G}(\Lambda)=\left\{\mathcal{G}_{G} v_{0}\right\}$ and all the assertions of Theorem 6.3.3.1 follows from Theorem 6.3.8.1 and Lemma 6.3.10.1.

Exercise 6.3.10.4 Check this in detail.

For proving Theorem 6.3.4.1 we need an additional
Lemma 6.3.10.4 Let $\Lambda$ be an indicator set, $v_{1}=r h_{1}, v_{2}=r h_{2}$. Then

$$
\left\{J_{G}(\Lambda) \notin \mathrm{MIN}\right\} \Longleftrightarrow\left\{\mathcal{G}_{G} v_{1} \quad \text { and } \mathcal{G}_{G} v_{2} \text { are non-minimal }\right\} .
$$

Proof. Suppose $w_{1}:=\mathcal{G}_{G} v_{1}$ and $w_{2}:=\mathcal{G}_{G} v_{2}$ are not minimal, i.e., $w_{1}-\delta r$ and $w_{1}-\delta r$ have subharmonic minorants $g_{1}$ and $g_{2}$.

Then $c g_{1}+(1-c) g_{2}$ is a minorant of the function $c w_{1}+(1-c) w_{2}-\delta r$, i.e., $J_{G}(\Lambda) \notin$ MIN. The inverse implication is trivial.

Exercise 6.3.10.5 Prove this.
Proof of Theorem 6.3.4.1. Suppose $\exp \Lambda$ is not complete. By Theorem 6.3.8.1 $J_{G} \notin$ MIN. By Lemma 6.3.10.4, $\mathcal{G}_{G} v_{1}$ and $\mathcal{G}_{G} v_{2}$ are not minimal. Hence $G_{1}$ and $G_{2}$ are freely enclosed in $G$ by Lemma 6.3.10.1. Since every one of these assertions is reversible, the inverse implication also holds. Analogously the other cases are proved.

Exercise 6.3.10.6 Prove all this in detail.
6.3.11 To prove Theorem 6.3.5.1 we need some auxiliary assertions.

Lemma 6.3.11.1 Let $\phi_{0}$ be a maximum point of t.c.f. $h(\phi)$ and $h\left(\phi_{0}\right) \geq 0$. Then

$$
\begin{equation*}
h(\phi) \geq h\left(\phi_{0}\right) \cos \left(\phi-\phi_{0}\right),\left|\phi-\phi_{0}\right| \leq \pi / 2 \tag{6.3.11.1}
\end{equation*}
$$

Proof. We write $y(\phi):=h\left(\phi_{0}\right) \cos \left(\phi-\phi_{0}\right)$. We have $y\left(\phi_{0}\right)=h\left(\phi_{0}\right)$ and $y(\phi)$ is a trigonometric function. If $y\left(\phi_{1}\right)=h\left(\phi_{1}\right)$ for some $\phi_{1}$ such that $\left|\phi_{1}-\phi_{0}\right|<\pi / 2$ this contradicts Theorem 3.2.5.2. If $y(\phi)$ does not intersect $h(\phi)$, this contradicts Theorem 6.2.3.4 applied to the function $h(\phi)-y(\phi)$, which is an $L_{\rho}$-subfunction with $\rho=1$.

Lemma 6.3.11.2 Let $H(\phi)$ be a trigonometric function on the interval $I=(\alpha, \beta)$ of length $\leq \pi$, such that $H(\phi)=0$ at one of the ends of $I$. Then every one of the conditions

1. $H\left(\phi_{0}\right)=0, \phi_{0} \in(\alpha ; \beta) ; \quad$ 2. $H(\phi)$ is zero on $\partial I$ with tangency;
implies $H(\phi) \equiv 0, \phi \in I$.

## Exercise 6.3.11.1 Prove this.

Lemma 6.3.11.3 Let $g \geq 0$ be a continuous periodic function, and let $\Theta_{\Lambda}, I_{\Lambda}, d_{\Lambda}$ be defined as in Theorem 6.3.5.1. In order that its maximal t.c.minorant be a trigonometrical function, it is necessary and sufficient satisfying at least one of the conditions:

1. $d_{\Lambda}<\pi$;
2. $d_{\Lambda}=\pi$ and $g(\phi)$ is zero with tangency on $\partial I$.

Proof. Necessity. Suppose $d_{\Lambda}>\pi$. Without loss of generality we can suppose that $I_{\Lambda}=(\alpha ;-\alpha)$, where $\alpha>\pi / 2$.

Set $\cos ^{+} \phi:=\max (\cos \phi, 0)$,

$$
\begin{equation*}
a=\inf \left(\frac{g(\phi)}{\cos ^{+} \phi}: \phi \in(-\alpha ; \alpha)\right) \tag{6.3.11.2}
\end{equation*}
$$

We have $a>0$. Set

$$
h(\phi):= \begin{cases}a_{1} \cos \phi, & |\phi| \leq \pi / 2  \tag{6.3.11.3}\\ 0 & |\phi|>\pi / 2\end{cases}
$$

where $a_{1} \leq a$.
The function $h(\phi)$ is a t.c.minorant of $g(\phi)$ and it is not a trigonometric function, which contradicts the supposition. Thus $d_{\Lambda} \leq \pi$.

Suppose $d_{\Lambda}=\pi$ and the condition to be zero with tangency on $\partial I$ does not hold. Then for $a$ defined by (6.3.11.2) the condition $a>0$ holds and $h(\phi)$ defined by (6.3.11.3) is a non-trigonometric minorant of $g$ that contradicts the supposition.

Sufficiency. Let the first condition hold and let $I=(\alpha ; \beta)$ be an arbitrary interval belonging to $\Theta_{\Lambda}$; let $h(\phi)$ be the maximal t.c.minorant of $g(\phi)$.

Set

$$
H(\phi):=h\left(\phi_{0}\right) \cos \left(\phi-\phi_{0}\right)
$$

where $\phi_{0}$ is the maximum point of $h(\phi)$ on $I$. From inequality (6.3.11.1) and the conditions $g(\alpha)=g(\beta)=0$ follows $H(\alpha)=H(\beta)=0$. Then, by Lemma 6.3.11.2, we obtain $H(\phi) \equiv 0$. Thus $h\left(\phi_{0}\right)=0$ and $h(\phi) \equiv 0$ for $\phi \in(\alpha ; \beta)$, i.e., $h(\phi)$ is trigonometric.

Let the second condition be fulfilled. Lemma 6.3.11.1 implies that $H(\phi)$ is zero with tangency on $\partial I$. By Lemma 6.3 .11 .2 we obtain that $h(\phi) \equiv 0$.

Proof of Theorem 6.3.5.1. Necessity. Let us note that if $v \in \operatorname{Fr}\left[\Phi_{\Lambda}\right]$, then for $c \in[0 ; 1]$ we have the equality

$$
\begin{equation*}
m(\lambda, G, v)=r\left(c\left(h_{2}-h_{1}\right)^{+}+(1-c)\left(h_{2}-h_{1}\right)^{-}\right)(\phi):=\operatorname{rm}(\phi, c) \tag{6.3.11.4}
\end{equation*}
$$

Let $\exp \Lambda$ be extremely overcomplete in $A(G)$. By Theorem 6.3.8.1 $J_{G} \subset$ HARM, i.e., for every $c \in[0 ; 1]$ the maximal t.c.minorant of the function $m(\phi, c)$ is trigonometric. Since

$$
\forall c \in[0 ; 1], \Theta_{\Lambda}=\{\phi: m(\phi, c)>0\}
$$

the necessity follows from Lemma 6.3.11.3.
Sufficiency. It follows directly from Lemma 6.3.11.3.
Exercise 6.3.11.2 Explain this.
6.3.12 To prove Theorem 6.3.6.1 we need

Lemma 6.3.12.1 If $w \in U[1]$ is non-minimal, then

$$
\mathbb{C}(w):=\left\{w_{[t]}: 1 \leq t \leq e^{P}\right\} \notin \operatorname{MIN} .
$$

It follows from Theorem 6.1.1.2, 2
Exercise 6.3.12.1 Explain this in detail.
Proof of Theorem 6.3.6.1. By Theorem 6.1.1.2, 2. $J_{G}(\Lambda)=\mathbb{C}\left(\mathcal{G}_{G} v\right)$. Thus Lemma 6.3.12.1 implies that $J_{G}(\Lambda) \notin$ MIN if and only if $\mathcal{G}_{G} v$ is not minimal. Thus Theorem 6.3.8.1 implies Theorem 6.3.6.1, 1 and 2.

Suppose $J_{G}(\Lambda) \subset$ HARM. Hence, $\mathcal{G}_{G} v=r H_{0}(\phi)$, where $H_{0}$ is trigonometric. Inversely, Lemma 6.3.12.1 implies $J_{G}(\Lambda)=\left\{r H_{0}(\phi)\right\}$.

Proof of Theorem 6.3.7.1. Let $\rho(\Lambda, G)>1$. Suppose $w_{q}:=g_{G} q$ exists. By definition of $\rho(\Lambda, G)$ we have $w_{q}(z) \leq 0$ for $z \in \partial D(\Lambda, G)$. By Theorem 6.2.3.3 (Maximum principle) $w_{q} \leq 0$ for $z \in D(G, \Lambda)$. Also $w_{q} \leq 0$ for $z \in \mathbb{T}_{P}^{2} \backslash D(G, \Lambda)$ by definition of $D(\Lambda, G)$. By Theorem 6.2.3.4 $w_{q} \equiv 0$ and hence is minimal. So $\exp \Lambda$ is complete by Theorem 6.3.6.1. If $\rho(\Lambda, G)=1$, then the system $\exp \Lambda$ is complete for every $G_{n} \ni G$, because of strict monotonicity of $\rho(\bullet)$ (see Section 6.2 .2 ) so $G$ is the maximal domain.

For the proof of Theorem 6.3.7.3 we need an auxiliary assertion. We suppose that $D$ is an image on $\mathbb{T}_{P}^{2}$ by the map (6.2.1.3) of the domain $G$ with a smooth boundary.

Theorem 6.3.12.2 Let $D \subset \mathbb{T}_{P}^{2}$ and $\rho(D) \leq 1$. Then $\rho\left(\mathbb{T}_{P}^{2} \backslash \bar{D}\right)>1$, if $D_{\Lambda} \neq$ $\{\Re z>0\}$.

For the proof we need the following assertion which was proved originally by A. Eremenko and M. Sodin:

Theorem 6.3.12.3 (Eremenko, Sodin) Let $\Gamma$ be a Jordan curve, connecting 0 and $\infty, T \Gamma=\Gamma$ for some $T>1$. Let $D_{+}, D_{-}$be domains, into which $\Gamma$ divides the plane, and let $\rho_{1}, \rho_{2}$ be the orders of the minimal harmonic functions in $D_{+}$and $D_{-}$respectively. Then

$$
\frac{1}{\rho_{1}}+\frac{1}{\rho_{2}} \leq 2
$$

and equality is attained only if $\Gamma$ consists of two rays.
We will give prove this theorem in Section 6.3.14.
Proof of Theorem 6.3.12.2. Let $\rho_{1}=\rho(D)$, and suppose $q_{1}(z)$ is a solution of boundary problem (6.2.2.1), $\rho_{2}=\rho\left(\mathbb{T}_{P}^{2} \backslash \bar{D}\right), q_{2}(z)$ is a solution of the corresponding boundary problem. Then the image of the boundary under the map $\lambda=e^{z}$ (we denote it as $\Gamma$ ) satisfies the conditions of Theorem 6.3.12.3 and the functions
$v_{1}(\lambda):=q_{1}(\log \lambda)|\lambda|^{\rho_{1}}$ and $v_{2}(\lambda):=q_{2}(\log \lambda)|\lambda|^{\rho_{2}}$ are positive harmonic functions in $D_{+}, D_{-}$with orders $\rho_{1}$ and $\rho_{2}$ respectively. By Theorem 6.3.12.3 we obtain

$$
1 / \rho(D)+1 / \rho\left(\mathbb{T}_{P}^{2} \backslash \bar{D}\right) \leq 2
$$

and equality holds only if $\Gamma$ is a pair of rays, i.e., $D_{\Lambda}=\{\Re z>0\}$.
Proof of Theorem 6.3.7.3. Necessity. Suppose $\rho\left(\Lambda, G_{\Lambda}\right)<1$. Let us prove that $\exp \Lambda$ is not complete. To this end we construct an $L_{1}$-minorant of $m\left(z, G_{\Lambda}, \Lambda\right)$ and prove that it is not minimal.

Let $D_{0} \Subset D\left(G_{\Lambda}, \Lambda\right)$ be a domain with smooth boundary for which $\rho\left(D_{0}\right)=1$. This is possible because of strict monotonicity $\rho(D)$ (Section 6.2.2). Let $q_{0}$ be a solution of the problem (6.2.2.1) satisfying the condition

$$
0<\max \left\{q_{0}(z): z \in D_{0}\right\} \leq \min \left\{m\left(z, G_{\Lambda}, \Lambda\right): z \in D_{0}\right\}-2 \epsilon
$$

for sufficiently small $\epsilon$. By Theorem 6.3.12.2, $\rho\left(\mathbb{T}_{P}^{2} \backslash \bar{D}_{0}\right)>1$. Thus the potential

$$
\Pi(z)=-\int_{D} G_{\rho}(z, \zeta, D) \nu(d \zeta)
$$

exists and $\nu$ can be chosen in such way that $\operatorname{supp} \nu \Subset \mathbb{T}_{P}^{2} \backslash \bar{D}_{0}$. By Proposition 6.2.3.6,

$$
\frac{\partial q_{0}}{\partial n}>0, z \in D_{0}
$$

Thus $\nu$ can be chosen in such a way that

$$
-\frac{\partial \Pi}{\partial n}<\min \frac{\partial q_{0}}{\partial n}, z \in \partial D_{0}
$$

Then the function

$$
q(z)= \begin{cases}q_{0}(z), & z \in D_{0} \\ \Pi(z), & z \in \mathbb{T}_{P}^{2} \backslash D_{0}\end{cases}
$$

is an $L_{1}$-subfunction on $\mathbb{T}_{P}^{2}$.
Exercise 6.3.12.2 Explain this in detail, exploiting Theorem 2.7.2.1.
The function $q(z)$ satisfies the condition

$$
q(z) \leq m\left(z, G_{\Lambda}, \Lambda\right)-2 \epsilon, \forall z \in \mathbb{T}_{P}^{2}
$$

because of negative potential. Hence,

$$
q_{1}(z):=q(z)+\eta
$$

for some $\eta>0$ also is a minorant of $m\left(z, G_{\Lambda}, \Lambda\right)$ and it is not minimal. Necessity is proved. Sufficiency follows from Theorem 6.3.7.1.
6.3.13 Now we pass to the proof of Theorem 6.3.7.4 and construction of Example 6.3.7.2.

Proof of Theorem 6.3.7.4. The set $\mathbb{T}_{P}^{2} \backslash D_{0}$ is closed. Let $\phi(z)$ be an infinitely differentiable function equal to zero on $\mathbb{T}_{P}^{2} \backslash D_{0}$ and positive on $D_{0}$. Set

$$
q(z):=h_{0}(y)-\epsilon \phi(z)
$$

where $h_{0}(y)$ is a t.c.f., corresponding to $G_{0}$, and let $\epsilon$ be small enough to satisfy $L_{1} q(z)>0, z \in \mathbb{T}_{P}^{2}$. It is possible, because $L_{1} h_{0}(y)>0$ by the condition of the theorem.

Then we have

$$
m\left(z, G_{0}, q\right)=\epsilon \phi(z), \quad \text { hence } \quad\left\{z: m\left(z, G_{0}, q\right)>0\right\}=D_{0}
$$

Take

$$
v(\lambda):=|\lambda| q(\log \lambda)
$$

and construct an entire function $\Phi_{\Lambda}$ for which

$$
\operatorname{Fr}\left[\Phi_{\Lambda}\right]=\left\{v_{[t]}: 1 \leq t \leq e^{P}\right\} .
$$

It is easy to check that the zero distribution of this function has all the properties demanded by Theorem 6.3.7.4.

Exercise 6.3.13.1 Check this.
Proof of (6.3.7.8). Consider the problem

$$
\begin{equation*}
L_{1} q(z)=0,\left.q\right|_{x=(2 \pi / P) y}=0 \tag{6.3.13.1}
\end{equation*}
$$

Let us pass in the equation to new coordinates

$$
\left\{\begin{array}{l}
\xi=x \cos \alpha+y \sin \alpha \\
\eta=-x \sin \alpha+y \cos \alpha, \tan \alpha=2 \pi / P
\end{array}\right.
$$

Then the equation takes the form:

$$
\left[\frac{\partial^{2}}{\partial \xi^{2}}+\frac{\partial^{2}}{\partial \eta^{2}}+2 \rho\left(\cos \alpha \frac{\partial}{\partial \xi}-\sin \alpha \frac{\partial}{\partial \eta}\right) r^{2}\right] R(\xi, \eta)=0
$$

The condition of being zero on $D_{0}$ is

$$
R(\xi, 2 \pi l \cos \alpha)=0 . l \in \mathbb{Z}
$$

The condition of periodicity gives

$$
R(\xi+(P / \cos \alpha) k, \eta)=R_{1}(\xi, \eta), k \in \mathbb{Z}
$$

We search for a solution that does not depend on $\xi$. We have

$$
R^{\prime \prime}(\eta)-2 \rho \sin \alpha R^{\prime}(\eta)+\rho^{2} R(\eta)=0, R(0)=R(2 \pi \cos \alpha)=0
$$

Further,

$$
R(\eta)=C_{1} e^{(\rho \sin \alpha) \eta} \cos ((\rho \cos \alpha) \eta)+C_{1} e^{(\rho \sin \alpha) \eta} \sin ((\rho \cos \alpha) \eta)
$$

Exploiting the boundary condition, we have

$$
\rho_{\min }=\left(2 \cos ^{2} \alpha\right)^{-1}=\frac{1}{2}\left[1+\left(\frac{2 \pi}{P}\right)^{2}\right]
$$

The corresponding eigenfunction is

$$
R=\exp \left(\rho_{\min } \sin \alpha\right) \eta \sin \left(\left(\rho_{\min } \cos \alpha\right) \eta\right)
$$

It is zero on $\mathbb{T}_{P}^{2} \backslash D_{0}$ and positive in $D_{0}$, so it is determined up to a constant multiple.
6.3.14 We are going to prove Theorem 6.3.12.3. Actually we prove

Theorem 6.3.14.1 Let $\Gamma_{1}, \Gamma_{2}, \ldots, \Gamma_{n}$ be Jordan curves, such that

1. $\Gamma_{i}, i=1,2, \ldots, n$ connect 0 and $\infty$;
2. there exists a number $T,|T|>1$ (not necessarily real) for which $T \Gamma_{i}=\Gamma_{i}, i=$ $1,2, \ldots, n$.
Let $D_{i}, i=1,2, \ldots, n$ be domains into which the plane is divided, and let $\rho_{i}$ be the order of the minimal harmonic function in $D_{i}$. Then

$$
\begin{equation*}
\sum_{i} 1 / \rho_{i} \leq 2 \tag{6.3.14.1}
\end{equation*}
$$

and equality holds if and only if $\Gamma_{i}$ are a logarithmic spirals (or rays, when $\left.T \in \mathbb{R}_{+}\right)$.

Proof. Denote by $H_{i}$ the minimal harmonic function in $D_{i}$. Then $H_{i}=\Im \phi_{i}$ where $\phi_{i}: D_{i} \mapsto \Pi^{+}$is a conformal map of $D_{i}$ to the upper half-plane, $\phi(0)=0$. The maps $g_{i}:=\phi_{i}\left(T \phi_{i}^{-1}\right): \Pi^{+} \mapsto \Pi^{+}$are continued by isomorphism to $\mathbb{C}$, and $g_{i}(0)=0$. Thus $g_{i}(z)=\sigma_{i} z$, where $\sigma_{i}>1$. Hence, $\phi_{i}(T z)=\sigma_{i} \phi_{i}(z)$ or

$$
T h_{i}(z)=h_{i}\left(\sigma_{i} z\right), h_{i}:=\phi_{i}^{-1}: \Pi^{+} \mapsto D_{i} .
$$

Now we exploit the following inequality from [Lev]

$$
\begin{equation*}
\sum_{i=1}^{n} \frac{1}{\log \sigma_{i}} \leq \frac{2 \log T}{|\log T|^{2}} \leq \frac{2}{\log T} \tag{6.3.14.2}
\end{equation*}
$$

The equalities in (6.3.14.2) are attained only when $\Gamma_{i}$ are logarithmic spirals or rays.

Since $\rho_{i}=\log \sigma_{i} / \log |T|$ (6.3.14.2) implies (6.3.14.1).

## Notation

$2.1 \mathbb{R}^{m}, M(f, x, \varepsilon), f^{*}(x), C^{+}(E), C^{-}(E), \chi_{G}, \chi_{F}, \Gamma_{A}, F^{A}, K, M(f, K)$, $K_{n}, K_{\max }$.
$2.2 \sigma(G), \mu, G_{0}(\mu), \operatorname{supp} \mu, \mathcal{M}(G), \mu_{F}(E), \nu, \mathcal{M}^{d}, \nu^{+}, \nu^{-},|\nu|, \operatorname{supp} \phi$, $\xrightarrow{*}, \stackrel{\circ}{E}, \bar{E}, \sigma\left(\mathbb{R}^{m_{1}} \times \mathbb{R}^{m_{2}}\right), \Phi_{1} \otimes \Phi_{2}, \mu_{1} \otimes \mu_{2}$.
$2.3 \varphi_{n} \xrightarrow{\mathcal{D}} \varphi, \alpha(t), \alpha_{\varepsilon}(x), \psi_{\varepsilon}(x),\langle f, \varphi\rangle,\left\langle\delta_{x}, \varphi\right\rangle,\left\langle\delta_{x}^{(n)}, \varphi\right\rangle,\langle\mu, \varphi\rangle,\langle\alpha f, \varphi\rangle$, $\left\langle f_{1}+f_{2}, \varphi\right\rangle,\left\langle\frac{\partial}{\partial x_{k}} f, \varphi\right\rangle, \quad f_{\epsilon}(x),\left.\quad f\right|_{G_{1}}, \widehat{\cos \rho}(\phi)$.
$2.4 \Delta_{\boldsymbol{x}^{0}}, \mathcal{E}_{m}(x), \theta_{m}, G(x, y, \Omega), G\left(x, y, K_{a, R}\right)$.
$2.5 \Pi(x, \mu, D), G_{N}(x, y), \Pi_{N}(x, \mu, D), \Pi(x, \mu), \Pi(z, \mu), \boldsymbol{c a p}_{G}(K, D)$, $\boldsymbol{\operatorname { c a p }}_{m}(K), \boldsymbol{\operatorname { c a p }}_{m}(D), \overline{\operatorname{cap}}_{m}(E), \underline{\operatorname{cap}}_{m}(E), \boldsymbol{\operatorname { c a p }}_{l}(K)$.
$2.6 \mathcal{M}(x, r, u), \mathcal{N}(x, r, u), E^{\epsilon}, D^{-\epsilon}, u_{\epsilon}(x), K_{R}, M(r, u), \mu(r, u), \mathcal{M}(r, u)$, $N(r, u), M(z)$.
$2.7 \tilde{u}(x), \mu_{x}(t), E\left(\alpha, \alpha^{\prime}, \epsilon, \mu\right), E_{n, \delta_{0}}$.
$2.8 a(r), \rho[a], \sigma[a], \rho(r), \sigma[a, \rho(r)], V(r), L(r), \delta S H\left(\mathbb{R}^{m}\right), T(r, u), \rho_{T}[u]$, $\sigma_{T}[u], \sigma_{T}[u, \rho(r)], \rho_{M}[u], \sigma_{M}[u], \sigma_{M}[u, \rho(r)], \quad \rho[\mu], \bar{\Delta}[\mu], \bar{\Delta}[\mu, \rho(r)]$, $N(r, \mu), \rho_{N}[\mu], \delta \mathcal{M}\left(\mathbb{R}^{m}\right)$.
$2.9 H(z, \cos \gamma, m), G\left(x, y, \mathbb{R}^{m}\right), D_{k}(x, y), H(z, \cos \gamma, m, p), G_{p}(x, y, m)$, $G_{p}(z, \zeta, 2), \Pi(x, \mu, p), \delta S H(\rho), \Pi_{<}^{R}(x, \nu, \rho-1), \Pi_{>}^{R}(x, \nu, \rho), \delta_{R}(x, \nu, \rho)$, $\delta_{R}(z, \nu, \rho), \delta_{R}(x, u, \rho), M(r, \delta), \bar{\Delta}_{\delta}[u, \rho], \Omega[u, \rho(r)], T(r, \lambda,>), T(r, \lambda,<)$.
$3.1 V_{t}, P_{t}, S H\left(\mathbb{R}^{m}, \rho, \rho(r)\right), S H(\rho(r)), u_{t}(x), \operatorname{Fr}\left[u, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right], U[\rho, \sigma]$, $U[\rho], v_{[t]}, \quad \mathcal{M}\left(\mathbb{R}^{m}, \rho(r)\right), \quad \mu \in \mathcal{M}(\rho(r)), \quad \operatorname{Fr}\left[\mu, \rho(r), V_{\bullet}, \mathbb{R}^{m}\right], \quad \operatorname{Fr}[\mu]$, $\mathcal{M}[\rho, \Delta], \quad \mathcal{M}[\rho], \quad \nu_{[t]}$.
$3.2 h(x, u), \underline{h}(x, u), l_{\boldsymbol{x}^{0}}, x^{0}(x), T_{\rho}, G_{I}(\phi, \psi), \Pi_{I}(\phi, d s), T C_{\rho}, C o_{\Omega}$.
$3.3 \bar{\Delta}(G, \mu), \bar{\Delta}(E, \mu), \underline{\Delta}(K, \mu), \underline{\Delta}(E, \mu), C o_{\Omega}(I), \bar{\Delta}^{\mathrm{cl}}(E), \underline{\Delta}^{\mathrm{cl}}(E), \Omega^{G}(\epsilon)$, $\Omega^{K}(\epsilon)$.
$4.1 T^{t},\left(T^{\bullet}, M\right), d(\bullet \bullet), \Omega\left(T^{\bullet}\right), \mathbb{C}(m), \Omega(m), A(m), T_{t} v$.
$4.2 U_{0}, \beta(x), b_{0}, k(s), R_{\epsilon} v(x), \quad \operatorname{Str}(\delta), v(x \mid t), v(\bullet, t)$
$4.3 \quad w(\bullet \mid t), w(\bullet \bullet \bullet)$.
$4.4 \boldsymbol{u},(\boldsymbol{u})_{t}, \operatorname{Fr}[\boldsymbol{u}], \boldsymbol{U}[\rho]$.
5.1 $M(r, f), T(r, f), \rho_{T}[f], \rho_{M}[f], \sigma_{T}[f, \rho(r)], \sigma_{M}[f, \rho(r)], n\left(K_{r}\right), n(r)$, $\rho[n], \bar{\Delta}[n], N(r, n), \rho_{N}[n], \bar{\Delta}_{N}[n], p[n], \operatorname{Fr}[f], \operatorname{Fr}[n], \operatorname{Mer}(\rho, \rho(r))$, $T(r, f), \rho_{T}[f], \sigma_{T}[f, \rho(r)]$.
$5.2 \alpha-\overline{\mathrm{mes}} C, C_{0}^{\alpha}, C_{0}^{0}$.
$5.3\|g\|_{p}$.
$5.4 \underline{h}_{1}(\phi, f), \underline{h}_{2}(\phi, f), \underline{h}(\phi, f)$.
5.5 $N(\delta, \mathcal{X}),(\mathcal{X}) \int f d \delta,(\mathcal{X}) \int_{E} f d \delta, \delta\left(\Theta^{F}\right), D_{r, \Theta}, \delta_{z}\left(D_{r, \Theta}\right), A^{c l}\left(\delta, \chi_{\Theta}\right)$.
$5.7 \mathcal{F}(u), H_{\phi}(u), T(u), M_{\alpha}(u), M(u), I_{\alpha \beta}(u), I(u, g), \overline{\mathcal{F}}[f], \underline{\mathcal{F}}[f], \chi_{H}$, $\chi_{I}, \chi_{F o}$.
5.8 $\quad K_{S_{1}}, \quad S_{1}, G(t, \gamma, \rho), \hat{\mathbb{G}}\left(s, S_{1}-S\right),(\mathcal{F} \nu)(s)$.
6.1 $H(z), m(z, v, H), \mathcal{G}_{H}, D_{H}, U_{\text {ind }}, \hat{U}_{\text {ind }}$.
$6.2 \mathbb{T}_{P}^{2}, \mathcal{D}^{\prime}\left(\mathbb{T}_{P}^{2}\right), q(z), L_{\rho}, E_{\rho}(\bullet-\zeta), E_{\rho}^{\prime}(\bullet-\zeta), q_{D}, G_{\rho}(z, \zeta, D), \mathcal{H}_{\rho}(q)$.
6.3 $\Lambda, \Phi_{\Lambda}(\lambda), \exp \Lambda, A(G), h_{\Lambda}(\phi), G_{\Lambda}, \alpha G_{1}+\beta G_{2}, \Theta_{\Lambda}, I_{\Lambda}, d_{\Lambda}, h_{G}(\phi)$, $m(\lambda, G, v), H(\lambda), q_{\Lambda}(z), D(G, \Lambda), \rho(\Lambda, G), g_{G} q, \mathcal{G}_{G}, D_{G}$, MIN, $J_{G}(\Lambda)$, HARM, $m(\phi, G, h), E(\phi)$.

## List of Terms

2.1 upper semicontinuous regularization upper semicontinuous function
lower semicontinuous function
2.2 measure
mass distribution
support of $\mu$
$\mu$ is concentrated on $E \in \sigma(G)$
restriction of $\mu$ onto $F \in \sigma(G)$
charge
positive and negative, respectively, variations of $\nu$
full variation of $\nu$
variation
Borel function
restriction of $\mu$ on the set $E$
product of measures
2.3 linear space
topological space
linear continuous functional on $\mathcal{D}$
Schwartz distribution
Dirac delta-function
the nth derivative of the Dirac delta-function regular distribution
positive distribution
product of a distribution $f$ by an infinitely differentiable function $\alpha(x)$
sum of distributions $f_{1}$ and $f_{2}$
partial derivative of distribution
sequence of distributions $f_{n}$ converges to a distribution $f$
regularization of the distribution $f$
restriction of distribution $f \in \mathcal{D}^{\prime}(G)$ to $G_{1} \subset G$
fundamental solution of $L$ at the point $y$ spherical operator
2.4 harmonic distribution

Lipschitz boundary, Lipschitz domain
harmonic measure
spherical function of a degree $\rho$
Green potential of $\mu$ relative to $D$
Newton potential
logarithmic potential
2.5 balayage, sweeping

Green capacity of the compact set $K$ relative to the domain $D$
Wiener capacity
external and inner capacity of any set $E$
capacible set
logarithmic capacity
irregular point
equilibrium mass distribution
$h$-Hausdorff measure
Carleson measure
2.6 mean value of $u(x)$ on the sphere $S_{x, r}:=\{y:|y-x|=r\}$
subharmonic function
the least harmonic majorant of $u$ in $K$
Riesz measure of the subharmonic function $u$
2.7 precompact family of functions
a sequence $f_{n}$ of locally summable functions converges in $L_{\text {loc }}$ quasi-everywhere convergence
a sequence of functions $u_{n}$ converges to a function $u$ relative to $\alpha$-Carleson measure
a point $x \in \mathbb{R}^{m}\left(\alpha, \alpha^{\prime}, \epsilon\right)$-normal with respect to the measure $\mu$
2.8 order of $a(r)$
type number of $a(r)$
$a(r)$ of minimal type
$a(r)$ of normal type
$a(r)$ of maximal type
convergence exponent for the sequence $\left\{r_{j}\right\}$
a proximate order with respect to order $\rho$
equivalent proximate orders
type number with respect to a proximate order
proper proximate order

Nevanlinna characteristic
order of $u(x)$ with respect to $T(r)$
characteristics $\rho_{M}[u], \sigma_{M}[u], \sigma_{M}[u, \rho(r)]$
convergence exponent of $\mu$
upper density of $\mu$
genus of $\mu$
$N$-order of $\mu$
$N$-type of $\mu$
2.9 Gegenbauer polynomials

Chebyshev polynomials
primary kernel
canonical potential
zero distribution
canonical Weierstrass product
3.1 limit set of the function $u(x)$
limit set of the mass distribution $\mu$
3.2 indicator of growth of $u$ lower indicator
$\rho$-subspherical function
$\rho$-trigonometrically convex ( $\rho$-t.c.)
fundamental relation of indicator
3.3 upper (lower) density of $\mu$
subadditivity of $\bar{\Delta}(E, \bullet)$
superadditivity of $\underline{\Delta}(E, \bullet)$
semi-additivity
generalized semi-additivity
monotonic function of $E \in \mathbb{R}^{m}$
$t$-extension of $E$
to be dense in
angular densities
4.1 dynamical system
$(\epsilon, s)$-chain from $m$ to $m^{\prime}$
chain recurrent dynamical system
non-wandering point
attractor
completely regular growth
polygonally connected set
periodic dynamical system
4.2 partition of unit
4.3 pseudo-trajectory
asymptotically dynamical pseudo-trajectory with dynamical asymptotics $T_{\bullet}$ (a.d.p.t.)
piecewise continuous pseudo-trajectory $w(\bullet \mid \bullet)$
$\omega$-dense pseudo-trajectory
4.4 subharmonic curve
5.1 entire function of order $\rho$ and normal type with respect to proximate order $\rho(r)$ entire function with prescribed limit set meromorphic function of order $\rho$ and normal type with respect to a proximate order $\rho(r)$
5.2 relative Carleson $\alpha$-measure
5.3 lower indicator of entire function
5.4 maximal interval of $\rho$-trigonometricity
strictly $\rho$-t.c.f.
concordant $h$ and $g$
5.5 upper density of zeros of entire function $(\mathcal{X})$-integral with respect to a nonnegative measure $\delta$
5.6 completely regular growth function CRG-function
regular zero distribution
regular zero distribution with integer $\rho$
completely regular growth functions
along curves of regular rotation
curve of regular rotation
5.7 growth characteristic
continuity, positive homogeneity asymptotic characteristics of growth
total family of growth characteristics non-rarefied set
rarefied set
thinly closed set
independent family of characteristic
6.1 ideally complementing $H$-multiplicator entire function is of minimal type
with respect to a proximate order $\rho(r), \rho(r) \rightarrow \rho$
limit set of indicators
the maximum principle for $U[\rho]$ is valid in the domain $G$
6.2 automorphic
connected on spirals
spectrum
strictly monotonic
minimal $v \in U[\rho]$
6.3 function of exponential type
completeness
maximality
extremal overcompleteness
maximal domain of completeness
extremely overcomplete system $\exp \Lambda$
trigonometrically convex function (t.c.f)
conjugate indicator diagram
regular set
$G_{\Lambda}$ is enclosed in $G$
enclosed with sliding
enclosed hardly
enclosed freely
indicator limit set
indicator set
zero with tangency
$\Lambda$ is periodic
$w \in U[1]$ is minimal
$U \subset U[1]$ is minimal

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OWS 38: Bobenko, A.I. / Schröder, P. / Sullivan, J.M. / Ziegler, G.M. (Eds.), Discrete Differential Geometry (2008). ISBN 978-3-7643-8620-7

Discrete differential geometry is an active mathematical terrain where differential geometry and discrete geometry meet and interact. It provides discrete equivalents of the geometric notions and methods of differential geometry, such as notions of curvature and integrability for polyhedral surfaces. Current progress in this field is to a large extent stimulated by its relevance for computer graphics and mathematical physics. This collection of essays, which documents the main lectures of the 2004 Oberwolfach Seminar on the topic, as well as a number of additional contributions by key participants, gives a lively, multi-facetted introduction to this emerging field.
OWS 37: Galdi, G.P. I Rannacher, R. I
Robertson, A.M. I Turek, S., Hemodynamical Flows (2008). ISBN 978-3-7643-7805-9
OWS 36: Cuntz, J. / Meyer, R. / Rosenberg, J.M., Topological and Bivariant $K$-theory (2007). ISBN 978-3-7643-8398-5

Topological K-theory is one of the most important invariants for noncommutative algebras. Bott periodicity, homotopy invariance, and various long exact sequences distinguish it from algebraic K-theory. We describe a bivariant K-theory for bornological algebras, which provides a vast generalization of topological K-theory. In addition, we discuss other approaches to bivariant K-theories for operator algebras. As applications, we study K-theory of crossed products, the Baum-Connes assembly map, twisted K-theory with some of its applications, and some variants of the Atiyah-Singer Index Theorem.

OWS 35: Itenberg, I. / Mikhalkin, G. /
Shustin, E., Tropical Algebraic Geometry (2007). ISBN 978-3-7643-8309-1

Tropical geometry is algebraic geometry over the semifield of tropical numbers, i.e., the real
numbers and negative infinity enhanced with the (max,+)-arithmetics. Geometrically, tropical varieties are much simpler than their classical counterparts. Yet they carry information about complex and real varieties.
These notes present an introduction to tropical geometry and contain some applications of this rapidly developing and attractive subject. It consists of three chapters which complete each other and give a possibility for non-specialists to make the first steps in the subject which is not yet well represented in the literature. The intended audience is graduate, post-graduate, and Ph.D. students as well as established researchers in mathematics.

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[^0]:    ${ }^{1} \mathcal{E}_{m}$ is slightly different from the fundamental solution (see, (2.3.6.2)), but this is traditional in Potential Theory

[^1]:    ${ }^{2}$ We will give the construction of such a function for the case of finite order (item 2.9.2), but it is possible actually always, see, for example, [HK, Thm. 4.1]

[^2]:    ${ }^{1}$ See Exercise 3.3.1.1

[^3]:    ${ }^{2}$ Exercise 3.3.1.2

[^4]:    ${ }^{3}$ See Exercise 3.3.1.3
    ${ }^{4}$ See Exercise 3.3.1.4

[^5]:    ${ }^{1}$ For the definition of $\mathrm{Co}_{(\alpha, \beta)}\left(I_{t}\right)$, see Exercise 3.3.1.5.

