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FOURIER TRANSFORMS OF HOLOMORPHIC FUNCTIONS AND APPLICATION TO NEWTON INTERPOLATION SERIES, II

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by
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ABSTRACT

This paper treats a generalization of the Martineau-Ehrenpreis theorem and applies it to the derivation of the Newton interpolation series for the largest possible class of functions. By means of Fourier transformation the Martineau-Ehrenpreis theorem establishes the isomorphism between analytic functionals with compact carrier and some space of entire functions. In this paper the analytic functionals are carried by unbounded convex sets with respect to some class of weightfunctions and its Fourier transforms are no longer entire functions, but they are holomorphic in cones in $\mathbb{C}^{n}$.

KEY WORDS \& PHRASES: Fourier transformation; analytic functionals carried by unbounded convex sets; holomorphic functions of several complex variables; cohomology with bounds; the MartineauEhrenpreis theorem on Fouriertransforms of analytic functionals; Newton interpolation series in several variables for non-entire functions.

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## 1. INTRODUCTION

This paper is the last of two papers dealing with Fourier transforms of holomorphic functions and the Newton interpolation series.

In [10] KIOUSTELIDIS derived the Newton series with the aid of Fourier transformation. The advantage of this method against the classical one (Cauchy's integral formula, NÖRLUND [13], GELFORD [5]) is that it treats the case of several variables as well. However, his treatment is valid for entire functions only. The aim of this paper is to show that this restriction is not due to the method, but that the method (namely the formalism of Fourier transformation) can be extended so as to include all possible nonentire functions for which the Newton series is valid.

In the first paper [14] the Newton series has been derived for functions, holomorphic in tubular radial domains, of polynomial growth in $|\operatorname{Re} z|$ and of exponential growth in $|\operatorname{Im} z|$. Such functions are the Fourier transforms of tempered distributions with support in unbounded convex sets according to a well known theorem (see [16]) generalizing the theorem of PALEY-WIENER-SCHWARTZ. In this case, however, one only uses the real part of the domain of convergence of the Fourier transformed Newton series. In [10] KIOUSTELIDIS has considered complex compact subsets of this domain using a Paley-Wiener type theorem, namely the theorem of EHRENPREIS [2] and MARTINEAU [12] dealing with Fourier transforms of analytic functionals with compact carrier. These Fourier transforms are entire functions of exponential type in $|z|$ and for such functions the Newton series is derived.

Generalizing the Ehrenpreis-Martineau theorem the main theorem of this paper states that holomorphic functions of exponential type in cones are the Fourier transforms of analytic functionals carried by unbounded convex sets with respect to some class of weightfunctions. One can formulate two versions of this theorem (based on formula (5.5) (i) and (5.5) (ii) respectively) and surprisingly it turns out that the apparently weaker version (i) equals the stronger version (ii). A particular case of version (i) has already been proved by KAWAI in [9]. However, this case cannot be handled very well in the derivation of the Newton series. Therefore, it still has sense to present the theorem as it is done here.

The proof of the main theorem is very different from the proof in [16] of the similar theorem in part $I$ [14]. In fact the last theorem in $2 n$ variables is used in proving the former in the $n$-dimensional case. The pattern of the proof is actually the same as that of Ehrenpreis' fundamental principle [3], only here one deals with non-entire functions. While in the Ehrenpreis-Martineau theorem the injectivity of the map $F$ (Fourier transformation) presents no problem, it seems to be the most difficult part of the generalization given here. For this part and for the transition from version (i) to version (ii) cohomology with bounds is used.

Together with the theorem on Fourier transforms some other theorems are given dealing with estimates for products of a polynomial matrix with a holomorphic non-entire vectorfunction similar to the case of entire functions given by HÖRMANDER in [7]. These theorems as well as the main theorem itself may be useful in other applications, for example if one is interested in solutions of systems of partial differential equations that can be written as boundary values of functions holomorphic in tubular radial domains.

The main theorem yields all the tools for deriving the Newton series for non-entire functions in several variables. Now the domain of convergence in $\mathbb{C}^{\mathrm{n}}$ is used completely, so that the most general class of functions is obtained for which the Newton series holds. This generalizes the case of one variable in NÖRLUND [13].

In section 2 the Ehrenpreis-Martineau theorem is discussed, and section 3 describes how the Newton series can be derived from this theorem as it is done by KIOUSTELIDIS in [10]. Section 4 deals with some properties of unbounded convex sets. Section 5 gives the space of holomorphic functions in cones in $\mathbb{C}^{n}$ of exponential type and the space of their Fourier transforms, which turns out to be the dual of some other space of holomorphic functions. These spaces are topologized in such a way that they are reflexive and that Fourier transformation is an isomorphism. A part of the version (i) of this isomorphism is also proved. In section 6 the main theorem of this paper, i.e. version (ii) of this isomorphism, is stated and the problems used to prove the main theorem are formulated. In section 7 these problems are solved, formulated so as to make them useful in other applications too. Here cohomology with bounds is derived and used. Section 8 gives some
corollaries and particular cases. Especially those concerned with functions holomorphic in tubular radial domains prepare section 9 , where the Newton series is derived for these functions. The appendix deals with the problem how to extend local relations between holomorphic functions to global relations. It uses cohomology as derived from the existence theorems for the $\bar{\partial}$-operator given by HÖRMANDER in [7]. Furthermore special coverings of open sets in $\mathbb{C}^{\text {n }}$ are constructed, adapted to the case of non-entire functions. Finally a short description of Ehrenpreis' fundamental principle is given.

## 2. ENTIRE FUNCTIONS OF EXPONENTIAL TYPE

This section deals with the relation between an entire function of exponential type and the carrier of its Fourier transform. It contains nothing new, but it is merely a rearrangement of some theorems of [7], stated in the appendix, in a way to make it suitable for generalization in section 5 .

Let $\Omega \subset \mathbb{C}^{\mathrm{n}}$ be an open set and let $\mathrm{A}(\Omega)$ be the space of in $\Omega$ holomorphic functions with the topology of uniform convergence on compact subsets K of $\Omega$. Elements $\mu$ of the strong dual $A^{\prime}(\Omega)$ of $A(\Omega)$ are called analytic functionals in $\Omega$. $A(\Omega)$ with the norm
(2.1) $\|f\|_{K}=\sup _{\zeta \in K}|f(\zeta)|,. \quad K \subset \subset \Omega$
is a linear subspace of $C(K)$, the space of continuous functions on the compact set $K$. Therefore, in view of the Hahn-Banach theorem and the theorem of Riesz, each analytic functional in $\Omega$ can be represented as a measure in a compact set K of $\Omega$. We say that an analytic functional $\mu$ in $\Omega$ is concentrated on the compact set K of $\Omega$, when for all $\mathrm{f} \in \mathrm{A}(\Omega)$

$$
|<\mu, f>| \leq M\|f\|_{K}
$$

with some positive constant $M$. In that case $\mu$ can be represented as a measure in $K$. Thus every analytic functional $\mu$ in $\mathbb{C}^{n}$ can be considered as an analytic functional in $\Omega$, where $\Omega$ is an arbitrary open neighborhood of the
compact set $\mathrm{K} \mu$ is concentrated on. We denote the space of analytic functionals in $\mathbb{C}^{\mathfrak{n}}$ concentrated on the compact sets of $\Omega$ as

$$
\mathrm{A}_{\Omega}^{\prime}\left(\mathbb{C}^{\mathrm{n}}\right)
$$

Conversely, analytic functionals in $\Omega$ are analytic functionals in $\mathbb{C}^{\mathrm{n}}$ too by means of their action on the restrictions to $\Omega$ of entire functions. This correspondence is $1-1$, when $\Omega$ is a Runge domain (see def. A5), for then $A\left(\mathbb{C}^{n}\right)$ is dense in $A(\Omega)$. Hence there is a map of $A^{\prime}(\Omega)$ onto $A_{\Omega}^{\prime}\left(\mathbb{C}^{n}\right)$, which is 1-1 when $\Omega$ is a Runge domain. For example, when $n=1$ one can think of $\Omega=\mathbb{C} \backslash\{0\}$; then the map

$$
\mathrm{A}^{\prime}(\Omega) \rightarrow \mathrm{A}_{\Omega}^{\prime}(\mathbb{C})
$$

is surjective, but not injective. Here $A_{\Omega}^{\prime}(\mathbb{C})=A^{\prime}(\mathbb{C})$, since by the maximum principle for every compact neighborhood $\hat{K}$ of 0 with boundary $K$ in $\mathbb{C} \backslash\{0\}$

$$
\begin{equation*}
\|f\|_{\mathrm{K}}=\|\mathrm{f}\|_{\hat{K}}, \tag{2.2}
\end{equation*}
$$

when f is entire.
We now give a more rigorous exposition of the foregoing. Let $\Omega$ be an open set in $\mathbb{C}^{\mathrm{n}}$ and K a compact subset of $\Omega$. Denote by

$$
\mathrm{A}(\overline{\mathrm{~K}})
$$

the space of functions holomorphic in a neighborhood of $K$ with the norm (2.1) and by

$$
\mathrm{A}_{\mathrm{K}}(\Omega)
$$

the space of functions holomorphic in $\Omega$ with the same norm. It is clear that $A_{K}(\Omega)$ is a linear subspace of $A(\bar{K})$ and that both spaces are not Banach spaces. Since we are only interested in their duals, it doesn't matter if we consider these spaces or their completions, the Banach spaces $\bar{A}(\bar{K})$ and
$\bar{A}_{K}(\Omega)$, respectively, consisting of functions continuous on $K$ and holomorphic in the interior of K if this is not empty. We denote by

$$
\mathrm{K} \hookrightarrow \Omega
$$

a sequence $\left\{K_{m}\right\}_{m=1}^{\infty}$ of compact subsets of $\Omega$ with int $K_{m} \subset K_{m} \subset$ int $K_{m+1} \subset$ $\subset K_{m+1} \subset \Omega$ and with $\bigcup_{m=1}^{\infty} K_{m}=\Omega$. Then we have the following characterization of the space $A(\Omega)$
$\mathrm{A}(\Omega)=\underset{\mathrm{Kcc}}{\operatorname{proj}} \lim \underset{\mathrm{K}}{\mathrm{K}}(\Omega)=\underset{\mathrm{Kcc} \Omega}{\operatorname{proj}} \lim \overline{\mathrm{A}}_{\mathrm{K}}(\Omega)=\underset{\mathrm{Kcc}}{\mathrm{proj}} \lim \mathrm{A}(\overline{\mathrm{K}})=\underset{\mathrm{Kcc}}{\operatorname{proj}} \lim \overline{\mathrm{A}}(\overline{\mathrm{K}})$.
Both $\bar{A}_{K}(\Omega)$ and $\bar{A}(\bar{K})$ are closed linear subspaces of $A_{\infty}(1 ; K)$, see [14] B. 4 or [18], so that according to [14] C. 6 and C. 7 the maps $\bar{A}_{K}(\Omega) \rightarrow \bar{A}_{S}(\Omega)$ and $\overline{\mathrm{A}}(\overline{\mathrm{K}}) \rightarrow \overline{\mathrm{A}}(\overline{\mathrm{S}})$ are compact, $\mathrm{S} \subset \subset \mathrm{K}$. Thus $\mathrm{A}(\Omega)$ is an $\overline{\mathrm{S}}$-space (see [14] F.8), which is nuclear according to [14] G.7. Since $\bar{A}_{K}(\Omega)$ is dense in $\bar{A}_{S}(\Omega)$, the dual can be represented as

$$
\begin{equation*}
\mathrm{A}^{\prime}(\Omega)=\underset{\mathrm{K} \subseteq \subseteq \Omega}{\operatorname{ind} \lim _{\mathrm{K}} A_{\mathrm{K}}^{\prime}(\Omega)} \tag{2.3}
\end{equation*}
$$

according to [14] F. 12. However, in general an element of $A_{K}^{\prime}(\Omega)$ does not uniquely determine an analytic functional in any neighborhood $\Omega^{\prime} \subset \Omega$ of K . This is true for distributions: distributions in 0 with support in $K$ are also distributions in $0^{\prime}, K \subset \subset 0^{\prime} \subset 0$. Only when $A(\Omega)$ is dense in $A_{K}\left(\Omega^{\prime}\right)$, representation (2.3) of $A^{\prime}(\Omega)$ is the inductive limit of all analytic functionals in any open neighborhood $\Omega^{\prime}$ of $K$ concentrated on $K$. So we must find a sequence $K \hookrightarrow \Omega$ for which $A\left(\bar{K}_{m+1}\right)$ is dense in $A\left(\bar{K}_{m}\right)$, for then $A(\Omega)$ is dense in $A(\bar{K})$ (see [14] lit.[2], $\S 26,2.5)$, thus also in $A_{K}\left(\Omega^{\prime}\right) \subset A(\bar{K})$. In that case (2.3) can be written as

$$
\begin{equation*}
A^{\prime}(\Omega)=\underset{\mathrm{K} c \subsetneq}{ } \operatorname{ind}^{\lim } \mathrm{A}^{\prime}(\overline{\mathrm{K}}) \tag{2.4}
\end{equation*}
$$

the inductive limit of analytic functionals concentrated on $K$. It is not possible to find such a sequence $\left\{K_{m}\right\}_{m=1}^{\infty}$ for all domains $\Omega$. Only for pseudoconvex domains $\Omega$ we will find one or, actually for domains of holomorphy,
but according to theorem A. 3 the domains of holomorphy are just the open pseudoconvex domains.

We define for any compact subset $K$ of $\Omega$ the set

$$
\hat{\mathrm{K}}_{\Omega}=\left\{\zeta\left|\zeta \in \Omega,|\mathrm{f}(\zeta)| \leq\|\mathrm{f}\|_{\mathrm{K}} \text { for all } \mathrm{f} \in \mathrm{~A}(\Omega)\right\} \xlongequal{\text { not }} \hat{\mathrm{K}} \subset \Omega,\right.
$$

see (A1). Hence for $f \in A(\Omega)$ we have $\|f\|_{K}=\|f\|_{\hat{K}}$, thus $A_{K}(\Omega)=A_{\hat{K}}(\Omega)$. $\Omega$ is a domain of holomorphy, if

$$
\begin{equation*}
\hat{\mathrm{R}}_{\Omega} \subset \subset \Omega, \text { whenever } \mathrm{K} \subset \subset \Omega \text {, } \tag{2.5}
\end{equation*}
$$

according to theorem A.1. In the sequel we will assume that $\Omega$ is pseudoconvex expressing that (2.5) is satisfied. Then the restriction map from $\mathrm{A}_{\widehat{\mathrm{K}}}(\Omega)$ into $\mathrm{A}(\overline{\hat{\mathrm{K}}})$ exists. According to theorem $\mathrm{A} .4 \mathrm{~A}_{\widehat{\mathrm{K}}}(\Omega)$ is a dense linear subspace of $A(\overline{\hat{K}})$. Hence
(2.6) $\quad \overline{\mathrm{A}}_{\mathrm{K}}(\Omega)=\overline{\mathrm{A}}\left(\overline{\hat{K}}_{\Omega}\right) \quad(\Omega$ pseudoconvex).

Thus for any sequence $\mathrm{K} \hookrightarrow \Omega, \hat{\mathrm{K}}_{\Omega}$ is the desired sequence satisfying (2.4). We have obtained that the closure of the space

$$
\mathrm{A}_{\Omega^{\prime}}(\Omega) \stackrel{\text { def }}{=} \underset{\mathrm{Kcc}}{\mathrm{Kroj}} \lim _{\Omega^{\prime}} A_{K}(\Omega)
$$

equals

$$
\overline{\mathrm{A}}_{\Omega^{\prime}}(\Omega)=\underset{\mathrm{Kccc}}{\operatorname{proj}} \lim \overline{\mathrm{~A}}_{\mathrm{K}}(\Omega)=\underset{\mathrm{Kcc}}{\operatorname{proj}} \lim _{\Omega^{\prime}} \overline{\mathrm{A}}\left(\overline{\hat{\mathrm{~K}}}_{\Omega}\right) .
$$

$A_{\Omega^{\prime}}(\Omega)$ is a prē-F $\bar{S}-$ space, that means its closure is an $\overline{\mathrm{S}}$-space. If $\Omega^{\prime}$ is such that $\mathrm{K} \hookrightarrow \Omega^{\prime}$ implies $\hat{\mathrm{K}}_{\Omega} \hookrightarrow \Omega$ we get $\mathrm{A}_{\Omega^{\prime}}(\Omega)=\mathrm{A}(\Omega)$; for example $\mathrm{A}_{\Omega}(\Omega)=$ $=\mathrm{A}(\Omega)$; another example is (2.2). However, we will consider cases where $A_{\Omega},(\Omega) \neq A(\Omega)$.

The strong dual of $A_{\Omega},(\Omega)$ is the LS-space

$$
\begin{equation*}
\mathrm{A}_{\Omega^{\prime}}^{\prime}(\Omega)=\underset{\mathrm{K} \mathrm{\subset} \mathrm{\hookrightarrow}}{\operatorname{ind}} \lim _{\Omega^{\prime}} \mathrm{A}^{\prime}\left(\overline{\hat{\mathrm{K}}}_{\Omega}\right) \tag{2.7}
\end{equation*}
$$

which yields (2.4) when $\Omega^{\prime}=\Omega$.
Let $\Omega_{1} \subset \Omega_{2}$ be both pseudoconvex open sets in $\mathbb{C}^{n}$; then the restriction maps

$$
\mathrm{A}\left(\Omega_{2}\right) \rightarrow \mathrm{A}_{\Omega_{1}}\left(\Omega_{2}\right) \hookrightarrow \mathrm{A}\left(\Omega_{1}\right)
$$

are continuous. The first map is a surjection from a Fréchet space onto a pré-Fréchet space and since $A_{K}\left(\Omega_{2}\right)$ is a linear subspace of $A_{K}\left(\Omega_{1}\right)$, the second map is in fact the embedding of the linear subspace $A_{\Omega_{1}}\left(\Omega_{2}\right)$ into $\mathrm{A}\left(\Omega_{1}\right)$. The transposed maps are the continuous maps

$$
\mathrm{A}^{\prime}\left(\Omega_{1}\right) \rightarrow \mathrm{A}_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right) \hookrightarrow \mathrm{A}^{\prime}\left(\Omega_{2}\right),
$$

where the first map is surjective according to the Hahn-Banach theorem and the second map is injective. We always have $\hat{\mathrm{K}}_{\Omega_{1}} \subset \hat{\mathrm{~K}}_{\Omega_{2}}$, but if

$$
\begin{equation*}
\hat{\mathrm{K}}_{\Omega_{1}}=\hat{\mathrm{K}}_{\Omega_{2}}, \tag{2.8}
\end{equation*}
$$

then in view of (2.3), (2.6) and (2.7) we have (see theorem A.7)

$$
\begin{equation*}
\mathrm{A}^{\prime}\left(\Omega_{1}\right)=\mathrm{A}_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right) . \tag{2.9}
\end{equation*}
$$

When each component of $\Omega_{2}$ contains a component of $\Omega_{1}$, for example when both are connected, $\mathrm{A}^{\prime}\left(\Omega_{1}\right)$ is dense in $\mathrm{A}^{\prime}\left(\Omega_{2}\right)$, for then $\mathrm{A}\left(\Omega_{2}\right)=\mathrm{A}^{\prime \prime}\left(\Omega_{2}\right)$ is mapped injectively into $\mathrm{A}\left(\Omega_{1}\right)=\mathrm{A}^{\prime \prime}\left(\Omega_{1}\right)$. We do not have this in the case of distributions: $E^{\prime}\left(0_{1}\right)$ is not dense in $E^{\prime}\left(0_{2}\right), 0_{1} \subset 0_{2} ; E^{\prime}\left(\overline{0}_{1}\right)$ is even a closed linear subspace of $E^{\prime}\left(0_{2}\right), \overline{0}_{1} \subset 0_{2}$ and $0_{1}$ convex (see [14], G.5).

The linear hull $L$ of the following set of entire functions in $\zeta$

$$
\left\{\mathrm{e}^{\mathrm{iz} \cdot \zeta^{2}}\right\}_{z \in \mathbb{C}^{n}}
$$

is dense in $\mathrm{A}\left(\Omega_{2}\right)$, when $\Omega_{2}$ is a Runge domain, so this set is dense in $\mathrm{A}_{\Omega_{1}}\left(\Omega_{2}\right)$ too. Indeed, differentiating $\mathrm{e}^{\mathrm{iz} \cdot \zeta}$ with respect to $z$ and setting $z=0$, we get $i \zeta$, so that we can approximate the polynomials by elements
of L. Therefore, the map

$$
F: \mu \in A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right) \leadsto f(z)=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta_{>}}\right.
$$

is an injective map from $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ into some set of entire functions $f$. Let $H_{K}$ be the function from $\mathbb{C}^{n}$ into $\mathbb{R}$

$$
H_{K}(z)=\sup _{\zeta \in K} \operatorname{Im}(-z \cdot \zeta), \quad K \subset \subset C^{n}
$$

We have $H_{K}=H_{c h(K)}$, where $c h(K)$ is the convex hull of $K$, see section 4. When $\mu$ is concentrated on $K, f=F(\mu)$ satisfies

$$
\begin{equation*}
|f(z)| \leq M \exp H_{K}(z) . \tag{2.10}
\end{equation*}
$$

Hence we define the Banach space (see [14] B.4)

$$
\operatorname{Exp}(K)=A_{\infty}\left(\exp -H_{K}(z) ; \mathbb{C}^{n}\right)
$$

and the LS-space

$$
\underset{\operatorname{Exp}}{\tilde{K}}(\Omega)=\underset{K \hookrightarrow \Omega}{\operatorname{ind}} \lim \operatorname{Exp}(K)
$$

We have $\operatorname{Exp}(\Omega)=\tilde{\operatorname{Exp}}(\operatorname{ch}(\Omega))$ and according to [14] G. $7 \tilde{\operatorname{Exp}}(\Omega)$ is nuclear.
Hence $F$ is an injective map from $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ into $\tilde{\operatorname{Exp}}\left(\Omega_{1}\right)$ when $\Omega_{2}$ is a Runge domain. Also $F$ is a bounded map, which follows from (2.10) and the fact that $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ being an LS-space is regular, see [14] F. 15 and F. 16 . Since $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ is bornological, $F$ is continuous. We will see that if $\Omega_{1}$ is convex, $F$ is surjective and its inverse is continuous too. Convex sets are Runge domains (see [16] 16.11), hence with $\Omega_{1}$ convex (2.8) is satisfied according to theorem A.6, so that then (2.9) holds.

THEOREM 2.1. Let $\Omega$ be a convex domain in $\mathbb{C}^{\text {n }}$. The map $F$ from $A^{\prime}(\Omega)$ into $\operatorname{Exp}(\Omega)$ given by

$$
F(\mu)(z)=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta}\right\rangle, \quad \mu \in A^{\prime}(\Omega)
$$

is an isomorphism.

Before proving this theorem we write $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ in a different way. We have introduced the notion of an analytic functional in $\Omega$ concentrated on the compact set $K \subset \subset \Omega$ and $A_{K}(\Omega)$ was a linear subspace of $C(K)$. However, $A(\Omega)$ also is a linear subspace of $E(\Omega)$, the space of all $C^{\infty}$-functions in $\Omega \subset \mathbb{C}^{\mathrm{n}}=\mathbb{R}^{2 \mathrm{n}}$ with the topology of uniform convergence of all derivatives on compact subsets. Indeed, all the derivatives of holomorphic functions converge on compact sets of $\Omega$ when the functions converge. So we can give A( $\Omega$ ) the topology induced by $E(\Omega)$ and each $\mu \in A^{\prime}(\Omega)$ can be extended to an element of some $E^{\prime}(K)$. Then $\mu$ is a distribution with compact support $K$ and for $f \in A(\Omega)$ we get

$$
|<\mu, f>| \leq M \sup _{|\alpha| \leq k}\left\|D^{\alpha} f\right\|_{K} ;
$$

so Cauchy's formula yields for all $\varepsilon>0$

$$
\begin{equation*}
|<\mu, f>| \leq M_{\varepsilon}\|f\|_{K_{\varepsilon}}, \tag{2.11}
\end{equation*}
$$

where $K_{\varepsilon}$ is a closed $\varepsilon$-neighborhood of $K$ in $\mathbb{C}^{n}$ with $K_{\varepsilon} \subset \subset \Omega$ and $M_{\varepsilon}$ is a positive constant depending on $\varepsilon$. When (2.11) holds we say that the analytic functional $\mu$ in $\Omega$ is carmied by $K$. Thus an analytic functional in $\Omega$ carried by $K$ is concentrated on any neighborhood of $K$ in $\Omega$. Sometimes it is said that $\mu$ is carried by such a neighborhood, see [15]: An analytic functional can be carried by several compact sets, but it is not true that it is carried by the intersection of all carriers, unlike the notion of support of a distribution, see [7] 4.5.

We will now describe the topology of $A(\Omega)$ using the concept of carrier, although this makes the description more complicated. Analytic functionals concentrated on compact sets are easier to describe, but analytic functionals carried by compact sets are easier to handle and are more natural as we will see.

Let $K \subset \subset \Omega$ be a compact subset of $\Omega$. We define the pré-LS-space (this means that its closure is an LS-space)

$$
\begin{equation*}
\mathrm{A}_{\overline{\mathrm{K}}}(\Omega)=\underset{\varepsilon \downarrow 0}{\operatorname{ind}} \lim \mathrm{~A}_{\varepsilon}(\Omega) \tag{2.12}
\end{equation*}
$$

The closure $\bar{A}_{K}(\Omega)$ of $A_{K}(\Omega)$ does not consist of holomorphic functions in $K$, but the closure $\overline{\mathrm{A}}_{\overline{\mathrm{K}}}(\Omega)$ of $\mathrm{A}_{\overline{\mathrm{K}}}(\Omega)$ consists of functions each holomorphic in a neighborhood of $K$. $\overline{\mathrm{A}}_{\overline{\mathrm{K}}}(\Omega)$ consists of all holomorphic functions in a neighborhood of $K$ when $\Omega$ is pseudoconvex and $K=\hat{K}_{\Omega}$ according to (2.6), for example when $K$ is convex. The dual of $A_{\bar{K}}(\Omega)$ is the $\overline{\mathrm{S}}$-space

$$
\begin{equation*}
{\underset{\bar{K}}{\prime}}_{\prime}^{(\Omega)}=\underset{\varepsilon \downarrow 0}{\operatorname{proj}} 1 \mathrm{im} \mathrm{~A}_{\varepsilon}^{\prime}(\Omega), \tag{2.13}
\end{equation*}
$$

the space of analytic functionals in $\Omega$ carried by K . Now $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ is the inductive limit of the spaces $\frac{\mathrm{A}_{\mathrm{K}}^{\prime}}{\mathrm{K}}(\Omega)$, namely

$$
\underset{\mathrm{Kcc}}{\operatorname{ind} \Omega_{1}} \lim \mathrm{~A}_{\mathrm{K}}^{\prime}\left(\Omega_{2}\right)=\mathrm{A}_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)=\underset{\mathrm{Kcc}}{\operatorname{ind} \Omega_{1}} \lim A_{\mathrm{K}}^{\prime}\left(\Omega_{2}\right) .
$$

Indeed, each $\mathrm{A}_{\overline{\mathrm{K}}}^{\prime}\left(\Omega_{2}\right)$ can be mapped continuously into $A_{\mathrm{K}_{\varepsilon}}^{\prime}\left(\Omega_{2}\right)$ and into $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ successively and conversely for all $\varepsilon>0$ each $A_{K}^{\prime}\left(\Omega_{2}\right)$ can be mapped continuously into $A_{K_{\varepsilon}}^{\prime}\left(\Omega_{2}\right)$, thus into $\mathrm{A}_{\overline{\mathrm{K}}}^{\prime}\left(\Omega_{2}\right)$, see [14] F.6. In this representation $A_{\Omega_{1}}^{\prime}\left(\Omega_{2}\right)$ is an LS-space too: a neighborhood of zero in $A_{\overline{\mathrm{K}}}^{\prime}\left(\Omega_{2}\right)$, that is a neighborhood of zero in some $A_{K_{\varepsilon}}^{\prime}\left(\Omega_{2}\right)$, is mapped into a relatively compact set of $A_{S_{\eta}}^{\prime}\left(\Omega_{2}\right)$ for any $\eta>0, K \subset \subset S$ and $\varepsilon$ small enough, thus into a relatively compact set of $A_{\bar{S}}^{\prime}\left(\Omega_{2}\right)$. This is in contrast with distributions, where the inductive limit $E^{\prime}(0)=$ ind $\lim E^{\prime}(K), K C 0$, is strict, when 0 and $K$ are convex, see [14] G.5.

Along the same lines one can see that $\tilde{\operatorname{Exp}}(\Omega)$ can be represented as the LS-space

$$
\tilde{\operatorname{Exp}}(\Omega)=\underset{\mathrm{Kcc} \Omega}{\operatorname{ind}} \lim \operatorname{Exp}\left(\mathrm{~K}_{0}\right)
$$

with

$$
\operatorname{Exp}\left(K_{0}\right) \stackrel{\text { def }}{=} \underset{\varepsilon \downarrow 0}{\operatorname{proj}} \lim A_{\infty}\left(\exp \left(-H_{K}(z)-\varepsilon\|z\| ; \mathbb{C}^{\text {n }}\right)=\underset{\varepsilon \downarrow 0}{\operatorname{proj}} \lim \operatorname{Exp}\left(K_{\varepsilon}\right) .\right.
$$

We will now prove theorem 2.1.
PROOF OF THEOREM 2.1. It is clear that $F$ maps $A_{\bar{K}}^{\prime}\left(\Omega_{2}\right)$ continuous $1 y$ and injectively into $\operatorname{Exp}\left(K_{0}\right)$. It is sufficient to prove that $F$ is a surjective map between the Fréchet-spaces $\mathrm{A}_{\overline{\mathrm{K}}}^{\prime}\left(\Omega_{2}\right)$ and $\operatorname{Exp}\left(\mathrm{K}_{0}\right)$, for then $F^{-1}$ is continuous according to the open mapping theorem. When $K$ is convex, this is exactly theorem 4.5.3 of [7] with $\Omega_{2}=\mathbb{C}^{n}$. Thus $F$ is an isomorphism between $A_{\Omega}^{\prime}\left(\mathbb{C}^{n}\right)=$ $=A^{\prime}(\Omega)$ and $\tilde{\operatorname{Exp}}(\Omega)$, when $\Omega$ is convex.

Theorem 2.1 is due to EHRENPREIS [2] and MARTINEAU [12]; for polydiscs this theorem can also be found in [15] and [18], where analytic functionals are used concentrated on compact sets. The notion of carrier of an analytic functional enables us to prove the continuity of $F^{-1}$ by the open mapping theorem. In [7], [15] and [18] $\mathrm{e}^{z \cdot \zeta}$ is used instead of $\mathrm{e}^{i z \cdot \zeta}$, but we use $e^{i z \cdot \zeta}$ in view of the the generalization in section 5 .

In the sequel we will start with a space of holomorphic functions of exponential growth. Let $a(y, x)$ be a continuous function of $z=x+i y$ on the unit sphere of $\mathbb{C}^{n}=\mathbb{R}^{2 n}$, such that the following function, which is homogeneous of degree one, is convex

$$
a(z) \stackrel{\text { def }}{=} a\left(\frac{y}{\|z\|}, \frac{x}{\left\|_{z}\right\|}\right)\|z\| .
$$

In that case we call $a(y, x)$ itself convex, see section 4 . This function determines a convex compact set $K \subset \mathbb{C}^{n}$ by

$$
K=\left\{\zeta \mid \zeta=\xi+i n,-y \cdot \xi-x \cdot n \leq \tilde{a}(z), z=x+i y \in C^{n}\right\}
$$

see section 4 . With this compact set $K$ we denote the space $\operatorname{Exp}(K)$ also as

$$
\operatorname{Exp}(\mathrm{a}) \stackrel{\text { def }}{=} \operatorname{Exp}(\mathrm{K})
$$

The function $a(y, x)+\varepsilon$ on the unit sphere determines the function $a(z)+\varepsilon\|z\|$ on $\mathbb{C}^{\mathrm{n}}$ and we have

$$
\operatorname{Exp}(a+\varepsilon)=\operatorname{Exp}\left(K_{\varepsilon}\right)
$$

Similarly we denote

$$
\operatorname{Exp}(\mathrm{a}+0) \stackrel{\text { def }}{=} \operatorname{Exp}\left(\mathrm{K}_{0}\right)
$$

We conclude this section with a corollary about the difference between analytic functionals and distributions expressed in properties of the spaces of their Fourier transforms.

Let $a$ and $b$ be two convex functions with $\check{a}(y, x) \leq \mathscr{B}(y, x)$ on $\mathbb{C}^{n}$ and let they determine the compact sets $K$ and $S$, respectively; then $K \subset S$. The LSspace $\overline{\mathrm{A}}_{\overline{\mathrm{S}}}(\Omega)$ is reflexive and it is mapped injectively into $\overline{\mathrm{A}}_{\overline{\mathrm{K}}}(\Omega)$ :

$$
\mathrm{A}_{\overline{\mathrm{S}}}^{\prime \prime}(\Omega) \subset \mathrm{A}_{\overline{\mathrm{K}}}^{\prime \prime}(\Omega) .
$$

Hence (see [15], corollary 5 to th. 18.1) ${\underset{\sim}{\mathcal{K}}}_{\prime}^{(\Omega)}$ is dense in $A_{\bar{S}}^{\prime}(\Omega)$ and since $F$ is an isomorphism $\operatorname{Exp}(a+0)$ is dense in $\operatorname{Exp}(b+0)$.

In [14] section 2 we have seen that the space of Fourier transforms of distributions with support in some compact set $K_{1}$ in $\mathbb{R}^{n}$ is a closed linear subspace of the space of Fourier transforms of distributions with support in a compact set $S_{1}$ with $K_{1} \subset S_{1}$. Let us take the example when $K_{1}=K$ and $S_{1}=S$ : let $K$ and $S$ be balls in the real part of $\mathbb{C}^{n}$ with radius $a$ and $b$ respectively $(a<b)$. Then $a(y, x)$ becomes $a l l y l l$ and $\tilde{a}(z)=a l l y$, so that we get

$$
\begin{aligned}
& \underset{\varepsilon \downarrow 0}{\operatorname{proj}} \lim A_{\infty}\left(e^{-(a+\varepsilon)\|y\|-\varepsilon\|x\|} ; \mathbb{C}^{n}\right)=\operatorname{Exp}(a+0) \xrightarrow{\text { dense }} \operatorname{Exp}(b+0) \\
& \underset{\mathrm{m} \rightarrow \infty}{\text { ind }} \lim _{\infty} A_{\infty}\left(\frac{e^{-a\|y\|}}{(1+\|z\|)^{m}} ; \mathbb{C}^{n}\right)=H\left(a ; C^{n}\right) \xrightarrow[\begin{array}{l}
\text { linear } \\
\text { subspace }
\end{array}]{ } H\left(b ; C^{n}\right) .
\end{aligned}
$$

The difference between these maps is not a consequence of the different structure of the spaces Exp and H. For we can make them look similar without changing them. Namely, it follows from the second map that

$$
\underset{\varepsilon \downarrow 0}{\operatorname{proj}} \lim \underset{m \rightarrow \infty}{\operatorname{ind}} \lim _{\infty}\left(\frac{e^{-(a+\varepsilon)\|y\|}}{(1+\|z\|)^{m}} ; \mathbb{C}^{n}\right)=H\left(a ; \mathbb{C}^{n}\right)
$$

and since the following injections are continuous

$$
\begin{aligned}
& A_{\infty}\left(e^{-(a+\varepsilon)\|y\|-\varepsilon\|x\|} ; \mathbb{C}^{n}\right) \longrightarrow \underset{m \rightarrow \infty}{\operatorname{ind} \lim _{m}}\left(\frac{e^{-(a+\varepsilon)\|y\|-\varepsilon\|x\|}}{(1+\|z\|)^{m}} ; \mathbb{C}^{n}\right) \longrightarrow \\
& \longrightarrow A_{\infty}\left(e^{-(a+n)\|y\|-\eta\|x\|} ; \mathbb{C}^{n}\right)
\end{aligned}
$$

with $\varepsilon<\eta$, we get, according to [14] F.6,
(2.14) $\quad \operatorname{proj} \underset{\varepsilon \downarrow 0}{\lim } \underset{m \rightarrow \infty}{\operatorname{ind}} \lim _{\mathrm{m}}^{\infty}\left(\frac{\mathrm{e}^{-(a+\varepsilon)\|y\|-\varepsilon\|x\|}}{(1+\|z\|)^{m}} ; \mathbb{C}^{n}\right)=\operatorname{Exp}(a+0)$.

It follows from [14] C. 3 that $H\left(a ; \mathbb{C}^{n}\right)$ can be mapped continuously into $\operatorname{Exp}(a+0)$. Hence the transposed map between the inverse Fourier transforms of these spaces, which are reflexive, is continuous:

$$
\mathrm{A}_{\overline{\mathrm{K}}}\left(\mathbb{C}^{\mathrm{n}}\right) \rightarrow E(\mathrm{~K})
$$

$A_{\bar{K}}\left(\mathbb{C}^{n}\right)$ is the space of all real analytic functions on $K \subset \mathbb{R}^{n}$ in view of (2.6) and the fact that $K=\hat{K}_{\mathbb{C}^{n}}$, because $K$ is convex. Actually, every compact set in $\mathbb{R}^{n}$ is polynomially convex in $\mathbb{C}^{n}$, see [11]. Since an analytic function, that vanishes in an open set in $\mathbb{R}^{n}$, vanishes, the above map is injective. Therefore, the map from $H\left(a ; C^{n}\right)$ into $\operatorname{Exp}(a+0)$ is dense and this implies that the distributions with support in a compact set $K$ in $\mathbb{R}^{n}$ are dense in the space of analytic functionals carried by $K$. Even since $\operatorname{Exp}(a+0)$ is dense in $\operatorname{Exp}(b+0)$, the distributions with support in $K$ are dense in the space of analytic functionals carried by the connected compact set $S$ in $\mathbb{C}^{n}$, where $K$ is any compact subset of $S$ for example $K$ may consist of only one point of $S$.
3. NEWTON SERIES FOR ENTIRE FUNCTIONS OF EXPONENTIAL TYPE

In this section we derive the Newton interpolation series (see [14]) for entire functions of exponential type. The same is treated by KIOUSTELIDIS in [10]. However, the form given here yields a stronger result on the convergence and serves as a good introduction to the generalization in section 9 .

Each vector $h$ in $\mathbb{C}^{n}$ determines a convex open set $\Omega_{h}$ in $\mathbb{C}^{n}$ by
(3.1) $\Omega_{h}=\left\{\zeta\left|\zeta \in \mathbb{C}^{n},\left|e^{-h \cdot \zeta}-1\right|<1\right\}\right.$.

LEMMA 3.1. For all $z \in \mathbb{C}^{\mathrm{n}}$ and $\mathrm{s} \in \mathbb{C}$ the sequence

$$
e^{i z \cdot \zeta} \sum_{k=0}^{N}\binom{\mathrm{~s}}{\mathrm{k}}\left(e^{-\mathrm{h} \cdot \zeta}-1\right)^{k}
$$

converges for $\mathrm{N} \rightarrow \infty$ to $\mathrm{e}^{\mathrm{i}(\mathrm{z}+\mathrm{ish}) \cdot \zeta}$ in the space $\mathrm{A}\left(\Omega_{\mathrm{h}}\right)$ regarded as functions of $\zeta$.

PROOF. The series converges uniformly on compact subsets of $\Omega_{h}$, which is the convergence of the space $\mathrm{A}\left(\Omega_{h}\right)$.

For $h \in \mathbb{C}^{\mathrm{n}}$ let f be a function in $\tilde{\operatorname{Exp}}\left(\Omega_{h}\right)$ and let $s \in \mathbb{C}$. Then using theorem 2.1 and lemma 3.1 we derive the Newton series

$$
\begin{align*}
f(z+i s h) & =\left\langle\mu_{\zeta}, e^{i(z+i s h) \cdot \zeta}\right\rangle=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta} \lim _{N \rightarrow \infty} \sum_{k=0}^{N}\binom{s}{k}\left(e^{-h \cdot \zeta}-1\right)^{k}\right\rangle= \\
& =\lim _{N \rightarrow \infty} \sum_{k=0}^{N}\binom{s}{k}\left\langle\mu_{\zeta}, e^{i z \cdot \zeta}\left(e^{-h \cdot \zeta}-1\right)^{k}\right\rangle= \\
& =\sum_{k=0}^{\infty}\binom{s}{k}\left\langle\mu_{\zeta}, \sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} e^{\left.i(z+i m h) \cdot \zeta_{>}\right\rangle=}\right.  \tag{3.2}\\
& =\sum_{k=0}^{\infty}\binom{s}{k} \sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} f(z+i m h)=\sum_{k=0}^{\infty}\binom{s}{k} \Delta_{i h}^{k} f(z),
\end{align*}
$$

where $\Delta_{i h} f(z) \xlongequal{\text { def }} f(z+i h)-f(z)$, so that

$$
\Delta_{i h}^{k} f(z)=\sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} f(z+i m h)
$$



$$
\mu_{\zeta} \sum_{k=0}^{N}\binom{s}{k}\left(e^{-h \cdot \zeta}-1\right)^{k}
$$

converges weakly in each $A_{K_{\varepsilon / 2}}^{\prime}\left(\Omega_{h}\right)$, hence it converges strongly in each $A_{K_{\varepsilon}}^{\prime}\left(\Omega_{h}\right)$, thus in $A_{\bar{K}}^{\prime}\left(\Omega_{h}\right)$. Therefore (3.2) with $f(z)=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta_{\nu}}\right.$ converges in
the topology of $\operatorname{Exp}\left(\mathrm{K}_{0}\right)$. So we have found that if $f$ satisfies for some $K \subset \subset \Omega_{h}$

$$
\begin{equation*}
\forall \delta>0:|f(z)| \leq M_{\delta} \exp \left(H_{K}(z)+\delta\left\|_{z}\right\|\right), \quad z \in \mathbb{C}^{n}, \tag{3.3}
\end{equation*}
$$

the series (3.2) converges according to:

$$
\begin{align*}
& \forall \varepsilon>0, \quad \forall \delta>0, \quad \exists N_{0}(\varepsilon, \delta) \geq N_{1}(s), \quad \forall N \geq N_{0}, \quad \forall z \in \mathbb{C}^{n} \\
& \left|f(z+i s h)-\sum_{k=0}^{N}\binom{s}{k} \Delta_{i h}^{k} f(z)\right|<\varepsilon A(s) \exp \left(H_{K}(z)+\delta\|z\|\right) \tag{3.4}
\end{align*}
$$

where $N_{1}(s)$ is determined by (5.1) in [14] and $A(s)$ by (5.4) in [14]. Thus there is certainly uniform convergence on compact subsets of $\mathbb{C}^{n}$, which is the convergence given in [10].

There exists a $\rho<1$ with for $\zeta \in K_{\varepsilon}\left|e^{-h \cdot \zeta}-1\right| \leq \rho$, so that

$$
\left|<\mu_{\zeta}, e^{i z \cdot \zeta} \sum_{k=0}^{N}\binom{s}{k}\left(e^{-h \cdot \zeta}-1\right)^{k}>\left|\leq C \exp \left(H_{K}(z)+\varepsilon\|z\|\right) \sum_{k=0}^{N}\right|\binom{s}{k}\right| \rho^{k} .
$$

Hence the series (3.2) converges absolutely.
We restate the results in
 Newton series (3.2) is valid; the series converges absolutely in the topology of $\tilde{\operatorname{Exp}}\left(\Omega_{h}\right)$ or more precisely (3.2) converges according to (3.4) when f satisfies (3.3); the series (3.2) converges uniformly in s on compact subsets of $\mathbb{C}$.

For a more detailed description of the function $H_{K}(z)$ when $K \subset \subset \Omega_{h}$, see KIOUSTELIDIS [10] Satz 9, when $K$ is given by (40) and (41).
4. CONVEX SETS

In this section we describe how a closed convex set 0 in $\mathbb{R}^{n}$ determines an open convex cone $C$ in $\mathbb{R}^{n *}$ and a homogeneous convex function $a n$ on $C$ and
that, conversely, $C$ and a determine a closed convex set 0 in $\mathbb{R}^{n}$.
A closed convex set 0 in $\mathbb{R}^{n}$ is the intersection of closed halfspaces. Let $H$ be the largest collection of halfspaces in $\mathbb{R}^{n}$ such that 0 is the intersection of halfspaces $H \in H$. Let $y$ be the unit vector perpendicular to the hyperplane $\partial H$ bordering a halfspace $H \in H$, in other words $y \in \mathbb{R}^{n *}$ is the linear functional which vanishes on the translation of $\partial H$ to the origin. We identify the action of a linear functional $y$ on $\xi \in \mathbb{R}^{n}$ with the innerproduct: $\langle y, \xi\rangle=y \cdot \xi$. Then the halfspace $H$ can be written as

$$
H_{y}(a)=\{\xi \mid-y \cdot \xi \leq a\}
$$

with $\mathrm{y} \in \mathbb{R}^{\mathrm{n} *}$ and a a real number. Thus we have

$$
0 \subset H_{y}(a) \subset H_{y}(b) \quad \text { when } b \geq a
$$

that is $H_{y}(a) \in H$ implies $H_{y}(b) \in H$.
The normals $y$ to $\partial H$ vary in a set $p r C$ on the unit sphere $S$ of $\mathbb{R}^{n *}$, when $H$ varies in $H$. For each $y \in \operatorname{prC}$ let $a(y)$ be the smallest of all the numbers a with $0 \subset H_{y}(a)$. Thus for each $y_{0} \in \operatorname{prC}$ and each sequence $a_{k} \uparrow a\left(y_{0}\right)$, there is a sequence $\xi_{k} \in 0$ with
(4.1) $\quad-y_{0} \cdot \xi_{k}=a_{k} \leq a\left(y_{0}\right)$.

Let $C$ be the cone in $\mathbb{R}^{n *}$ determined by pr $C$

$$
C=\{y \mid y \neq 0, \tilde{y} \in \operatorname{pr} C\}
$$

with the notation $\tilde{y}=y /\|y\|$. Hence any closed convex set 0 in $\mathbb{R}^{n}$ determines a cone $C$ in $\mathbb{R}^{n *}$ and a function $a(y)$ on pr $C$ such that
(4.2) $0=\{\xi \mid-y \cdot \xi \leq a ̃(y) \xlongequal{\text { def }} a(\tilde{y})\|y\|, y \in C\}$.

It is clear that for $y \in C$ the function

$$
\begin{equation*}
I_{0}(y) \stackrel{\text { def }}{=} \sup _{\xi \in 0}-y \cdot \xi \tag{4.3}
\end{equation*}
$$

satisfies $I_{0}(y) \leq x(y)$ and that $0=\left\{\xi \mid-y \cdot \xi \leq I_{0}(y), y \in C\right\}$. Since $a(y)$ is the smallest possible function determining the set 0 , we have

$$
\begin{equation*}
\tilde{a}(y)=I_{0}(y), \quad y \in C \tag{4.4}
\end{equation*}
$$

The cone $C$ is convex, for

$$
-\left(t y_{1}+(1-t) y_{2}\right) \cdot \xi=-t y_{1} \cdot \xi-(1-t) y_{2} \cdot \xi \leq t I_{0}\left(y_{1}\right)+(1-t) I_{0}\left(y_{2}\right)
$$

with $\xi \in 0,0 \leq t \leq 1$ and $y_{1}, y_{2} \in C$, hence

$$
0 \subset H_{t y_{1}}+(1-t) y_{2}\left(t I_{0}\left(y_{1}\right)+(1-t) I_{0}\left(y_{2}\right)\right)
$$

From this it also follows that the function $I_{0}(y)$ is convex, that is

$$
I_{0}\left(t y_{1}+(1-t) y_{2}\right) \leq t I_{0}\left(y_{1}\right)+(1-t) I_{0}\left(y_{2}\right), \quad y_{1}, y_{2} \in C
$$

Taking into account (4.4) we find that $\check{a}(y)$ is convex, hence continuous and a(y) is bounded from below on pr C. We say that the continuous function $a(y)$ on $\operatorname{prC}$ is convex, when the function $\mathfrak{a}(y)$, which is homogeneous of degree one, is convex on $C$.

It is possible that the cone $C$ is contained in a linear subspace of $\mathbb{R}^{\mathrm{n} *}$ of lower dimension. Therefore, we consider $C$ in the lowest dimensional space containing it. Then we speak of the interior int $C$ of $C$ and we show that the open cone int $C$ determines the same convex set $O$ as $C$. We denote the closed convex set 0 determined by a cone $C$ and a convex function $a(y)$ on pr C according to (4.2) by $O(\mathrm{a} ; \mathrm{C})$.

It is clear that $O(a ; C) \subset O(a ; i n t C)$. Now let $\xi_{0}$ be a point outside $O(a ; C)$, then there is a vector $y_{0} \in \operatorname{prC}$ such that

$$
-\xi_{0} \cdot y_{0}>\left(y_{0}\right)
$$

Hence there is an $\varepsilon>0$ with

$$
-\xi_{0} \cdot y_{0}>a\left(y_{0}\right)+\varepsilon
$$

Since $a(y)$ is continuous on pr C, there is a $y \in$ pr int C with

$$
\left|a(y)-a\left(y_{0}\right)\right|<\varepsilon / 2 \quad \text { and } \quad\left\|y-y_{0}\right\|<\varepsilon /\left(2\left\|\xi_{0}\right\|\right) .
$$

Hence

$$
-\xi_{0} \cdot y=-\xi_{0} \cdot y_{0}-\xi_{0} \cdot\left(y-y_{0}\right)>a\left(y_{0}\right)+\varepsilon-\varepsilon / 2>a(y),
$$

thus $\xi_{0} \notin O(a ;$ int $C)$ by (4.2).
So we have found that each closed convex set in $\mathbb{R}^{\mathrm{n}}$ determines an open convex cone $C$ in $\mathbb{R}^{\mathrm{n} *}$ (open relatively to a linear subspace of $\mathbb{R}^{\mathrm{n} *}$ ) and a convex function $a(y)$ on pr C. Now we will prove that, conversely, each open convex cone $C$ in $\mathbb{R}^{n *}$ and each convex function a on pr C determine a closed convex set 0 in $\mathbb{R}^{n}$ by (4.2) that satisfies (4.4).

Indeed, 0 is convex and closed being the intersection of closed halfspaces and we only have to prove (4.1). Since ã is convex and C is open, we can find for each $y_{0} \in \operatorname{prC}$ a linear function on $C$, say $\alpha \cdot y$ for some vector $\alpha$, with $\alpha \cdot y \leq a(y), y \in C$ and $\alpha \cdot y_{0}=a\left(y_{0}\right)$. Then the point $\xi_{0}=-\alpha \in \mathbb{R}^{n}$ satisfies $-\xi_{0} \cdot y=\alpha \cdot y \leq \tilde{a}(y)$ for all $y \in C$, thus $\xi_{0} \in 0$. Furthermore, $-\xi_{0} \cdot y_{0}=\alpha \cdot y_{0}=a\left(y_{0}\right)$, hence (4.1) holds. We have also obtained that in (4.1) we may take $a_{k}=a\left(y_{0}\right)$ and $\xi_{k}=\xi_{0} \in 0$ for all $k$, when $y_{0} \in \operatorname{pr}$ int C. COROLLARY 4.1. Each closed convex set 0 in $\mathbb{R}^{n}$ determines an open (with respect to some linear subspace of $\mathbb{R}^{n *}$ ) convex cone $C$ in $\mathbb{R}^{n *}$ and a continuous convex function a on pr C by (4.3) such that (4.2) holds. Conversely, each open convex cone $C$ in $\mathbb{R}^{\mathrm{n} *}$ and each convex function a on pr $C$ determine a closed convex set $0(a ; C)$ in $\mathbb{R}^{n}$ by (4.2), such that (4.4) is satisfied.

We give some examples. Let C be an open cone in the first quadrant of $\mathbb{R}^{2 *}$. We consider the function an on some straight line $\ell \cap C \subset \mathbb{R}^{2 *}$.


Then the dual cone $C^{\star}$ is $C^{\star}=\{\xi \mid-y \cdot \xi \leq 0, y \in C\} \subset \mathbb{R}^{2}$. We have the following cases with different behaviour of the convex function a near the boundary of $\ell \cap C$ :

I


ã(y) is vertical at the boundary of $C$ and the boundary of 0 is asymptotically parallel to the boundary of $C^{*}$.

II


ã (y) is not vertical at the boundary of $C$ and outside some compact set $K$ the boundary of 0 is parallel to the boundary of $C^{*}$.

III


$\mathfrak{a}(\mathrm{y})$ tends to infinity when y approaches the boundary of C and the distance between the boundaries of 0 and $C^{*}$ increases to infinity.

In cases $I$ and III we say that the function $a$ is vertical at the boundary of $C$. In case II, when a is not vertical at the boundary, we may consider
the closed convex set $0^{\prime}=0 \cap K^{\prime}$, where $K^{\prime}$ is a compact set the interior of which contains $K$. Then $0^{\prime}$ is compact, so the cone $C^{\prime}$ determined by it is $\mathbb{R}^{\mathrm{n} *}$ and the convex function $a^{\prime}$ on $\mathrm{pr} \mathbb{R}^{\mathrm{n} *}=\mathrm{S}$ determined by $0^{\prime}$ coincides with a on pr C, that is

$$
\begin{equation*}
a^{\prime}(y)=a(y), \quad y \in \operatorname{pr} C \tag{4.5}
\end{equation*}
$$

Thus when $\mathfrak{a}$ is not vertical at the boundary of $C$, it can be extended to a convex homogeneous function $a^{\prime}$ on the whole of $\mathbb{R}^{n *}$.

Finally we describe how we can exhaust 0 by closed convex sets $0_{m}$ not touching the boundary of 0 , namely $0=U_{m=1}^{\infty} O_{m}$ with $O_{m}$ closed convex sets satisfying $0_{m} \subset$ int $0_{m+1} \subset 0_{m+1} \subset$ int $0 \subset \mathbb{R}^{n}, m=1,2, \ldots$. Let $\left\{a_{m}\right\}_{m=1}^{\infty}$ be an increasing sequence of convex functions on pr C with $a_{m}(y)<a_{m+1}(y)<$ $<a(y)$ for $y \in \operatorname{prC}$ and $\lim _{m \rightarrow \infty} a_{m}(y)=a(y), y \in \operatorname{prC}$ and moreover, either there are positive numbers $\varepsilon_{m}$ with $a(y)-a_{m}(y) \geq \varepsilon_{m}, y \in \operatorname{pr} C$, or all the functions $a_{m}$ are vertical at the boundary of $C$. Then the sets $O_{m}=0\left(a_{m} ; C\right)$ satisfy the conditions. When the functions $a_{m}$ are vertical at the boundary of $C$, the boundary of each $0_{m}$ approaches the boundary of 0 . Otherwise the boundaries of $O_{m}$ and 0 have a distance greater than $\varepsilon_{m}$.

When $C$ does not contain a straight line (then $C^{*}$ is not contained in a proper linear subspace of $\mathbb{R}^{n}$ and conversely) let $C_{k}$ be a sequence of open convex subcones of $C$ with pr $C_{k} \subset \operatorname{pr} \bar{C}_{k} \subset \operatorname{pr} C_{k+1} \subset \operatorname{pr} \bar{C}_{k+1} \subset C \subset \mathbb{R}^{n *}$ and $U_{k=1}^{\infty} C_{k}=C$. We call $C_{k}$ a relatively compact open subcone of $C$ and we write $C_{k} \subset \subset$. When $C^{*}$ is contained in a linear subspace of $\mathbb{R}^{n}$ we take open cones $C_{k}^{*}$ in this subspace with $\operatorname{pr} C_{k}^{*} \supset \operatorname{pr} C_{k+1}^{*} \supset \operatorname{pr} C_{k+1}^{*} \supset C^{*}$ and the cones $C_{k}$ are defined by the interior of the duals of $C_{k}^{*}: C_{k}=\operatorname{int}\left(C_{k}^{*}\right)$. Also in this case we call $C_{k}$ a relatively compact subcone of $C$.

The functions $a+1 / k$ on $\operatorname{pr} C$ are also convex and the sets $0\left(a+1 / k ; C_{k}\right)$ are of type II. Then

$$
O(a ; C)=\bigcap_{k=1}^{\infty} 0\left(a+1 / k ; C_{k}\right)
$$

but none of the sets $0\left(a+1 / k ; C_{k}\right)$ is contained in $0(a+1 ; C)$.
In the next sections we will regard $\mathbb{C}^{n}$ as a $2 n$-dimensional real space
$\mathbb{R}^{2 \mathrm{n}}$ by the identification

$$
\zeta=\xi+\mathrm{i} n \in \mathbb{C}^{\mathrm{n}} \leftrightarrow(\xi, n) \in \mathbb{R}^{\mathrm{n}} \times \mathbb{R}^{\mathrm{n}}=\mathbb{R}^{2 \mathrm{n}} .
$$

We identify the action of an element $z$ in the complex dual space $\mathbb{C}^{\text {n* }}$ with the ordinary product of complex numbers

$$
\langle z, \zeta\rangle=z \cdot \zeta=z_{1} \zeta 1+\ldots+z_{n} \zeta_{n}
$$

and we identify $\mathbb{C}^{\mathrm{n} *}$ with $\mathbb{R}^{2 \mathrm{n}}$ by

$$
z=x+i y \in \mathbb{C}^{n *} \leftrightarrow(y, x) \in \mathbb{R}^{n *} \times \mathbb{R}^{n *}=\mathbb{R}^{2 n *}
$$

Then regarded as $2 n$-dimensional real vectors the action of $z \in \mathbb{R}^{2 n *}$ on $\zeta \in \mathbb{R}^{2 \mathrm{n}}$ is

$$
\operatorname{Im} z \cdot \zeta=(y, x) \cdot(\xi, n)=y \cdot \xi+x \cdot n .
$$

When $C$ is a cone in $\mathbb{R}^{\mathrm{n} *}$ the set $\mathrm{T}^{\mathrm{C}}=\mathbb{R}^{\mathrm{n}}+\mathrm{iC}$ is a cone in $\mathbb{C}^{\mathrm{n} *}$ containing a straight line. The dual cone ( $\left.\mathrm{T}^{\mathrm{C}}\right)^{*}$ is the cone $\mathrm{C}^{*}$ contained in the imaginary subspace of $\mathbb{C}^{\mathrm{n}}$. Relatively compact subcones, constructed in the above way, are $\mathbb{R}^{n *}+i C_{k}$, where $C_{k} \subset c$.

## 5. FUNCTIONS OF EXPONENTIAL TYPE HOLOMORPHIC IN CONES

In this section we discuss the space of functions of exponential type, holomorphic in cones in $\mathbb{C}^{\mathrm{n}}$ and the space of their Fourier transforms (sometimes called Fourier-Borel transforms or Fourier-Laplace transforms).

Let $C$ be an open convex cone in $\mathbb{C}^{\mathrm{n}}$, which is identified with $\mathbb{R}^{2 \mathrm{n}}$ by $z=x+i y \in \mathbb{C}^{n} \leftrightarrow(y, x) \in \mathbb{R}^{n} \times \mathbb{R}^{n}=\mathbb{R}^{2 n}$. Let $a(z)$, regarded as a function of ( $\mathrm{y}, \mathrm{x}$ ), be continuous on pr C , such that the function

$$
\tilde{a}(z)=\|z\| a\left(\frac{z}{\|z\|}\right)=\|(y, x)\| a\left(\frac{y}{\|(y, x)\|}, \frac{x}{\|(y, x)\|}\right)
$$

is convex in $C$. A function $f$ holomorphic in $C \subset \mathbb{C}^{n}$ is of exponential type $a$, when for all $\varepsilon>0$ and $\delta>0$ and for all open relatively compact subcones $C^{\prime}$ cc $C$ constants $M\left(\varepsilon, \delta, C^{\prime}\right)$ exist such that

$$
|f(z)| \leq M\left(\varepsilon, \delta, C^{\prime}\right) e^{\tilde{a}(z)+\varepsilon\|z\|}, \quad z \in C^{\prime},\|z\|>\delta
$$

We denote the space of all such functions $f$ by $\operatorname{Exp}(a+0 ; C)$ or sometimes by Exp and we give this space a topology of an $\mathrm{F} \overline{\mathrm{S}}$-space by means of

$$
\operatorname{Exp}(a+0 ; C)=\underset{k \rightarrow \infty}{\operatorname{proj}} \lim _{\infty}\left(e^{-a ̃(z)+1 / k\|z\|} ; C(k)\right)
$$

where $C(k)=C_{k} \cap\{z \mid\|z\|>1 / k\}$ and $\left\{C_{k}\right\}_{k=1}^{\infty}$ is an increasing sequence of open relatively compact subcones of $C$ exhausting $C$ (see section 4). According to [14] G. 7 this space can also be written as a projective limit of Hilbert spaces $(\operatorname{Exp}(a+0 ; C)$ is nuclear).

We will construct a reflexive space $A^{\prime}$, which is the dual of some space A of holomorphic functions, such that Exp is the Fourier transform of $A^{\prime}$. Assume that there is a continuous map $F^{t}$ from Exp' into the completion $\bar{A}$ of $A$

$$
\begin{equation*}
F^{\mathrm{t}}: \operatorname{Exp}{ }^{\prime} \rightarrow \overline{\mathrm{A}} \tag{5.1}
\end{equation*}
$$

then the transposed map $F$ is a continuous map between the duals. So, since Exp is reflexive we get

$$
\begin{equation*}
F: A^{\prime} \rightarrow \operatorname{Exp} \tag{5.2}
\end{equation*}
$$

and since $A^{\prime}$ is reflexive, $F^{t}$ is the transposed map of $F$.
In order to get information about $\bar{A}$ we investigate Exp'. According to
[14] C. 3 and $F .6$ we can write Exp also as the FS-space

$$
\operatorname{Exp}(a+0 ; C)=\underset{k \rightarrow \infty}{\operatorname{proj}} \lim A_{\infty, 0}(\exp (-\tilde{a}(z)-1 / k\|z\|) ; \overline{C(k)})
$$

where $A_{\infty, 0}(M ; \bar{\Omega})$ consists of functions holomorphic in $\Omega$ and continuous on $\bar{\Omega}$
with $\mathrm{M}|\mathrm{f}|<\infty$ on $\bar{\Omega}$ and with $\mathrm{M}|\mathrm{f}|=0$ at infinity. Hence by [14] B. 5 Exp' is the inductive limit of spaces, whose elements $\sigma$ can be represented as bounded measures $\sigma(z)$ in $C(k)$, namely for $f \in \operatorname{Exp}$

$$
\begin{equation*}
\langle\sigma, f\rangle=\int_{\overline{C(k)}} f(z) \exp (-\tilde{a}(z)-1 / k\|z\|) d \sigma(z) \tag{5.3}
\end{equation*}
$$

and

$$
\int_{\overline{C(k)}}|d \sigma(z)|<\infty .
$$

Next we define the map $F^{t}$. Therefore we regard $\mathbb{C}^{n}$ with elements $z=(y, x)$ as the dual $\mathbb{R}^{2 n *}$ of some other space $\mathbb{R}^{2 n}$, whose elements are denoted by $(\xi, n)$ and which is identified with $\mathbb{C}^{n}$ by $\zeta=\xi+i n$. Then $\operatorname{Im} z \cdot \zeta=\langle(y, x),(\xi, \eta)\rangle$. The cone $C \subset \mathbb{R}^{2 n *}$ and the convex function $a(z)$ on pr $C$ determine a closed convex set $\Omega$ in $\mathbb{C}^{\mathrm{n}} \cong \mathbb{R}^{2 \mathrm{n}}$ by (4.2)

$$
\begin{equation*}
\Omega \stackrel{\text { not }}{=} \Omega(a ; C) \xlongequal{=}\{\zeta \mid \operatorname{Im} z \cdot \zeta \leq a(z), z \in \operatorname{pr} C\} . \tag{5.4}
\end{equation*}
$$

Furthermore let us introduce the closed convex sets $\Omega_{k}$ either by
(5.5) (i) $\quad \Omega_{k}=\Omega\left(a+1 / k ; C_{k+2}\right)$
or by
(5.5)(ii) $\Omega_{k}=\Omega(a+1 / k ; C)$.

In both cases $\Omega=\bigcap_{k=1}^{\infty} \Omega_{k}$. It is easy to see that $e^{i z \cdot \zeta}$ belongs to Exp if $\zeta$ belongs to $\Omega$. Therefore, we can define the map $F^{t}$ (5.1) by

$$
\begin{equation*}
\sigma \in \operatorname{Exp}^{\prime}: F^{t}(\sigma)(\zeta)=\left\langle\sigma_{z}, \mathrm{e}^{\mathrm{iz} \cdot \zeta}\right\rangle \quad \text { for } \zeta \in \Omega \tag{5.6}
\end{equation*}
$$

The representation (5.3) yields for some $k$

$$
\begin{equation*}
\phi(\zeta)=F^{t}(\sigma)(\zeta)=\int_{\overline{C(k)}} e^{-a(z)-1 / k\|z\|} e^{i z \cdot \zeta} d \sigma(z) . \tag{5.7}
\end{equation*}
$$

In both cases (5.5)(i) and (5.5)(ii) $\phi$ is holomorphic in int $\Omega_{k}$ and satisfies there for some $K>0$
(5.8) $\quad|\phi(\zeta)| \leq K \exp \left(-\delta_{k} 1 / k\|\zeta\|\right), \quad \zeta \in \Omega_{k}$,
where $\delta_{k}=\sin \alpha_{k}$ the minimum distance in radials between pr $C_{k+1}$ and pr $C_{k}$, see [14] proof of lemma 6.3. Indeed, for

$$
\zeta \in U=C_{k+1}^{*} \cap\left\{\zeta \mid\|\zeta\| \geq-a / \delta_{k}, \quad a=\min _{z \in \operatorname{pr} C_{k}}(a(z)+1 / k)\right\}
$$

we have

$$
\begin{equation*}
-\tilde{a}(z)-1 / k\|z\|+\operatorname{Im}-z \cdot \zeta \leq-\left(a+\delta_{k}\|\zeta\|\right)\|z\| \leq-a / k-\delta_{k} 1 / k\|\zeta\| \tag{5.9}
\end{equation*}
$$

when $z \in \overline{C(k)}$ and the set $\Omega_{k} \cap \overline{U^{c}}$ is compact in both cases (5.5)(i) and (5.5)(ii), so that (5.8) follows.

Therefore, we introduce the weightfunctions

$$
M_{k}(\zeta)=\exp \delta_{k} \frac{1}{k}\|\zeta\|
$$

Then it follows from (5.9) that the map $F^{t}$ given by (5.6) is a bounded, hence continuous, map from Exp' into the LS-space

$$
\begin{equation*}
\underset{k \rightarrow \infty}{\operatorname{ind}} \lim A_{\infty}\left(M_{k} ; \text { int } \Omega_{k}\right) \tag{5.10}
\end{equation*}
$$

in both cases (5.5)(i) and (5.5) (ii), see [14] F.11, F. 16 and C.7.
Our aim is to choose such a space A that the map $F$ (5.2) is an isomorphism. In [14] we have seen that, when the support of a distribution is contained in all the sets $\Omega_{k}$, it is contained in $\Omega$. But when the analytic functionals in $A^{\prime}$ are concentrated on all the sets $\Omega_{k}$, we cannot immediately conclude that they are concentrated on $\Omega$. Therefore, we do not yet know which of the alternatives (5.5)(i) or (5.5) (ii) we should take.

Now we will define linear subspaces $A_{i}$ and $A_{i i}$ of (5.10), depending on (5.5)(i) or (5.5)(ii) respectively, such that the map $F$ (5.2) is injective. In fact we give the linear hull L of the set $\left\{e^{i z \cdot \zeta}\right\} \quad{ }_{z \in C}$ the topology of the
space (5.10).
Let $L_{k}$ be the linear hull of the set

$$
\left\{\mathrm{e}^{\mathrm{iz} \mathrm{\cdot} \mathrm{\zeta}}\right\}_{z \in \mathrm{C}(\mathrm{k})}
$$

of functions in $\zeta$. We provide $L_{k}$ with the norm

$$
\|\cdot\|_{k}=\sup _{\zeta \in \Omega_{k}} M_{k}(\zeta)|\cdot(\zeta)|
$$

and denote it by

$$
A_{M_{k}} ; \Omega_{k}\left(L_{k}\right) .
$$

Then $A_{i}$ and $A_{i i}$ are defined as

$$
A \stackrel{\text { not }}{=} A_{M ; \bar{\Omega}}(L) \xlongequal{\text { def }} \underset{k \rightarrow \infty}{\lim } A_{M_{k}} ; \Omega_{k}\left(L_{k}\right) \text {, }
$$

where $\Omega_{k}$ is given by (5.5)(i) or (5.5)(ii), respectively. The closure $\bar{A}$ of both spaces in an LS-space, namely

$$
\begin{equation*}
\bar{A}=\operatorname{ind} \underset{k \rightarrow \infty}{\lim } \bar{A}_{M_{k}} ; \Omega_{k}\left(L_{k}\right), \tag{5.11}
\end{equation*}
$$

since LS-spaces are complete (see [14] F.14). $\bar{A}$ consists of functions each holomorphic in a neighborhood int $\Omega_{k}$ of $\Omega$ (compare 2.12). The duals $A_{i}^{\prime}$ of $A_{i}$ and $A_{i i}^{\prime}$ of $A_{i i}$ are $F \bar{S}-$ spaces

$$
\begin{equation*}
A^{\prime} \stackrel{\text { not }}{=} A_{M ; \bar{\Omega}}^{\prime}(L)=\underset{k \rightarrow \infty}{\operatorname{proj}} \lim A_{M_{k}}^{\prime} ; \Omega_{k}\left(L_{k}\right) \tag{5.12}
\end{equation*}
$$

(compare 2.13). We only have to check that A is not too small, in other words that (5.1) still holds. By letting $\sigma(z)$ be $\delta$-functions we see that $L$ is contained in the range of $F^{t}$ and when we write (5.7) as a defining "Lebesgue sum", it follows from (5.9) that this sum converges in the topology of $\bar{A}_{M_{k+1}} ; \Omega_{k+1}\left(L_{k+1}\right)$ to $\phi(\zeta)$. Hence $F^{t}(5.1)$ has dense image, so that $F(5.2)$ is an injective map from $A_{i}^{\prime}$ and from $A_{i}^{\prime}$, given by (5.12), into Exp.

The definition of $F(5.2)$ as the transposed map of (5.1) yields for $\mu \in A^{\prime}$

$$
\begin{aligned}
\forall \sigma \in \operatorname{Exp} & :<\sigma_{z}, F\left(\mu_{\zeta}\right)(z)>\stackrel{\text { def }}{=}\left\langle\mu_{\zeta}, F^{t}\left(\sigma_{z}\right)(\zeta)>=\left\langle\mu_{\zeta},\left\langle\sigma_{z}, e^{i z \cdot \zeta} \gg=\right.\right.\right. \\
& =<\sigma_{z},<\mu_{\zeta}, e^{i z \cdot \zeta} \gg
\end{aligned}
$$

by Fubini's theorem, so that $F(5.2)$ may as well be defined as

$$
\begin{equation*}
F(\mu)(z)=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta}\right\rangle, \quad z \in C \tag{5.13}
\end{equation*}
$$

like (5.6).
Now $F$ is a continuous injective map from $A_{i}^{\prime}$ and from $A_{i i}^{\prime}$ into Exp. Since $A_{i}$ can be continuously embedded into $A_{i i}, A_{i}^{\prime}$ is a priori larger than $A_{i i}^{\prime}$. So it is easier to prove that $F$ is also surjective from $A_{i}^{\prime}$ onto Exp. In that case the inverse map would be continuous according to the open mapping theorem, because $A_{i}^{\prime}$ and Exp are $F \bar{S}$-spaces and $F$ would be an isomorphism between $A_{i}^{\prime}$ and Exp. If we can also prove that $F$ is a surjective map from $A_{i i}^{\prime}$ onto Exp, then $A_{i i}^{\prime}$ too would be isomorphic to Exp, so that $A_{i j}^{\prime}=A_{i}^{\prime}$ and $\bar{A}_{i}=\bar{A}_{i i}$. First we will prove the apparently weaker version (i), theorem 5.1, of the main result of this paper, namely that $F$ is an isomorphism between $A_{i}^{\prime}$ and Exp. Then in a next section we investigate the spaces $\bar{A}_{i}$ and $\bar{A}_{i i}$ and finally in section 7 we will show that $F\left(A_{i i}^{\prime}\right)=\operatorname{Exp}$, which is the stronger version (ii) of the main theorem of this paper, theorem 6.1.

THEOREM 5.1. Let a be a convex function on pr $C$ for some open convex cone $C$ in $\mathbb{C}^{\mathrm{n}}$ and let $\Omega$ and $\Omega_{\mathrm{k}}$ be the closed convex sets in $\mathbb{C}^{\mathrm{n}}$ determined by (5.4) and (5.5)(i) respectively. Then the map $F$ from $A_{M}^{\prime} ; \bar{\Omega}(L)$ (5.12) into $\operatorname{Exp}(\mathrm{a}+0 ; \mathrm{C})$ given by (5.13) is an isomorphism.

PROOF. We only have to prove the surjectivity of the map $F$. So given an $f \in \operatorname{Exp}$, we have to find for each $k=1,2, \ldots$ a linear functional $\mu^{k}$ on $L_{k}$ with

$$
\mathrm{f}(z)=\left\langle\mu \frac{\mathrm{k}}{\zeta}, \mathrm{e}^{\mathrm{i} z \cdot \zeta}\right\rangle, \quad z \in \overline{\mathrm{C}(\mathrm{k})}
$$

and with

$$
\left|<\mu_{\zeta}^{k}, \phi(\zeta)>\left|\leq K_{k} \sup _{\zeta \in \Omega_{k}} \exp \left(\delta_{k} \frac{1}{k}\|\zeta\|\right)\right| \phi(\zeta)\right|
$$

for $\phi \in A_{M_{k}} ; \Omega_{k}\left(L_{k}\right)$.
Like in the proof of the theorem for entire functions (see HÖRMANDER [7], EHRENPREIS [3]) we try to extend $f$ as a holomorphic function $F$ in $2 n$ complex variables $v$ satisfying a certain bound and apply a Paley-WienerSchwartz type theorem. Precisely, choose an integer $k$ and assume that we have found a function $F_{k}$ of the complex variables $v \xlongequal{n o t}\left(v^{1}, v^{2}\right) \in \mathbb{C}^{n} \times \mathbb{C}^{n}=\mathbb{C}^{2 n}$ holomorphic in $\mathbb{R}^{2 n}+i C(q)$ with $q>\max \left(k+2,(k+1) / \delta_{k+1}\right)$ that satisfies for some m
(5.14) $\quad\left|F_{k}\left(\nu^{1}, \nu^{2}\right)\right| \leq M_{k}(1+\|\nu\|)^{m} \exp (\tilde{a}(\operatorname{Im} v)+1 / q\|\operatorname{Im} v\|)$
for $\operatorname{Im} v \in C(q)$ and

$$
\begin{equation*}
\mathrm{F}_{\mathrm{k}}(\mathrm{z}, \mathrm{iz})=\mathrm{f}(\mathrm{z}) \quad \text { for } \mathrm{z} \in \mathrm{C}(\mathrm{q}) \tag{5.15}
\end{equation*}
$$

Then we can apply theorem 9.1 of [14] (remark 9.1 and formula (9.5) or in fact the $5^{\text {th }}$ line from below on page 61 , since $\left.1 / q<\delta_{k+1} /(k+1)\right)$, which says that $\mathrm{F}_{\mathrm{k}}$ can be written as

$$
F_{k}(\nu)=\left\langle\mu \sum_{\xi, \eta}^{k}, e^{i \nu^{1} \cdot \xi+i \nu^{2} \cdot \eta}\right\rangle, \quad \operatorname{Im} \nu \in C(k+1)
$$

with $\mu^{k} \in\left(S_{m+n+1}^{k+1}\left(a+1 / q ; C_{q}\right)\right)^{\prime}$. This means that for $\phi \in S^{k+1 *}\left(a+1 / q ; C_{q}\right)$

$$
\left|<\mu_{\xi, \eta}^{\mathrm{k}}, \phi(\xi, \eta)>\right| \leq K_{k}^{\prime} \quad \sup _{\substack{|\mathrm{p}| \leq m+n+1 \\(\xi, n) \in O\left(a+1 / q ; C_{q}\right)}}
$$

$$
(1+\|(\xi, n)\|)^{m+n+1} \exp \left(\delta_{k+1} \frac{1}{k+1}\|(\xi, n)\|\right)\left|D^{p_{\phi}}(\xi, n)\right|
$$

Identifying $\mathbb{R}^{2 n}$ with $\mathbb{C}^{n},(\xi, \eta) \leftrightarrow \zeta=\xi+i n$, we get, because $C(k) \subset C(k+1)$,

$$
f(z)=\left\langle\mu_{\zeta}^{k}, e^{i z \cdot \zeta}\right\rangle, \quad z \in \overline{C(k)}
$$

with

$$
\begin{aligned}
& \left|<\mu_{\zeta}^{\mathrm{k}}, \phi(\zeta)>\right| \leq \mathrm{K}_{\mathrm{k}}^{\prime} \sup _{\substack{|\mathrm{p}| \leq \mathrm{m}+\mathrm{n}+1 \\
\zeta\left(\mathrm{a}+1 / \mathrm{q} ; \mathrm{C}_{\mathrm{q}}\right)}}^{\left.(1+\|\zeta\|)^{\mathrm{m}+\mathrm{n}+1}{\exp \left(\delta_{\mathrm{k}+1}\right.} \frac{1}{\mathrm{k}+1}\|\zeta\|\right) \quad\left|\mathrm{D}^{\mathrm{p}} \phi(\zeta)\right| \leq} \\
& \leq \mathrm{K}_{\mathrm{k}} \sup _{\zeta \in \Omega_{\mathrm{k}}} \exp \left(\delta_{\mathrm{k}} \frac{1}{\mathrm{k}}\|\zeta\|\right)|\phi(\zeta)|
\end{aligned}
$$

for any $\phi \in A_{M_{k} ; \Omega_{k}}\left(L_{k}\right) \subset S^{k+1^{*}}\left(a+1 / q ; C_{q}\right)$, since int $\Omega\left(a+1 / q ; C_{q}\right) \subset \Omega_{k}$ because $q>k+2$.

Now we have to find an $\mathrm{F}_{\mathrm{k}}$ satisfying (5.14) and (5.15). For an arbitrary $\ell$ we will construct a function $F_{k, \ell}$ that satisfies

$$
\begin{equation*}
\left|F_{k, \ell}\left(\nu^{1}, \nu^{2}\right)\right| \leq M_{k, \ell}(1+\|v\|)^{m} \exp (\tilde{a}(\operatorname{Im} v)+1 / \ell\|\operatorname{Im} v\|), \operatorname{Im} v \in C(q) . \tag{5.16}
\end{equation*}
$$

The construction follows the same pattern of the proof of theorem 4.4.3 of [7], only here we have to be careful near the boundary of $\mathbb{R}^{n}+i C$.

Since an open domain $\mathbb{R}^{n}+i B$ in $\mathbb{C}^{n}$ is pseudoconvex (domain of holomorphy) only if $B$ is convex (theorem A.2), we will use domains of the form $\mathbb{R}^{n}+\mathrm{i} \operatorname{ch}(\mathrm{C}(\mathrm{q})$ ), where $\mathrm{ch}(\mathrm{C}(\mathrm{q}))$ is the convex hull of $\mathrm{C}(\mathrm{q})$. This does not change anything, because for all $q$ there is a $p$ with $C(q) \subset \operatorname{ch}(C(q)) \subset C(p)$.

Let

$$
\begin{align*}
C(q+1)_{\delta, j}= & \left\{(y, x)\left|x_{1}=x_{1}^{0}, \ldots, x_{j}=x_{j}^{0}, \quad\right| x_{k}-x_{k}^{0} \mid<\delta\right.  \tag{5.17}\\
& \text { for } \left.k=j+1, \ldots, n \text { and }\left(y, x^{0}\right) \in \operatorname{ch}(C(q+1))\right\} \subset \mathbb{R}^{2 n}
\end{align*}
$$

Then $\operatorname{ch}(C(q+1))=C(q+1)_{\delta, n} \subset \ldots \subset C(q+1)_{\delta, j} \subset \ldots \subset C(q+1)_{\delta, 0}$. We can choose $\delta>0$ so small that there exists an integer $p>q+1$, such that

$$
\begin{equation*}
C\left(q^{+1}\right)_{\delta, 0} \subset C(p) \tag{5.18}
\end{equation*}
$$

Let $\psi_{k}$ be a $C^{2}$-function in $\mathbb{C}$ between 0 and 1 which is equal to 1 in the disc with radius $1 / 2 \delta$ and vanishes outside the disc with radius $\delta$. We write the coordinates in $\mathbb{C}$ as $w=u+i v$. Then there is a constant $K_{k}$ with

$$
\left|\frac{\partial \psi}{\partial \bar{w}}(w)\right| \leq K_{k}, \quad w \in \mathbb{C} .
$$

Let us define the $(0,1)$-form (see appendix section II) $\psi^{\prime}(w)=\partial \psi / \partial \bar{w}(w) d \bar{w}$ and let $w_{j}=i v{ }_{j}^{1}-\nu_{j}^{2}$, then $d \bar{w}_{j}=-i d \bar{\nu}{ }_{j}^{-1}-d \bar{\nu}_{j}^{2}$. When $f$ is regarded as an element of $A_{\infty}(\exp (-a ̃(z)-1 / \ell\|z\|) ; C(p))$ we define the function $F_{k, \ell}$ as follows:

$$
\begin{aligned}
& F_{k, \ell}\left(\nu^{1}, \nu^{2}\right)=F_{k, \ell}(\nu)= \\
& =\prod_{j=1}^{n} \psi_{k}\left(w_{j}\right) f\left(\nu^{1}\right)-\sum_{j=1}^{n}\left[\prod_{m=j+1}^{n} \psi_{k}\left(w_{m}\right)\right]_{j} U_{j}^{\ell}\left(\nu^{1} ; v_{1}^{2}, \ldots, v_{j}^{2}\right)
\end{aligned}
$$

for certain functions $U_{j}^{\ell}$ in $n+j$ complex variables. When $\operatorname{Im} v \in C(q+1)$, $\Pi_{j=1}^{n} \psi_{k}\left(w_{j}\right)$ vanishes for $v^{1}=\left(\operatorname{Im} v^{1}, \operatorname{Re} \nu^{1}\right) \notin C(p)$ according to (5.17) and (5.18), thus $F_{k, \ell}$ is defined for $v \in \mathbb{R}^{2 n}+i \operatorname{ch}(C(q+1))$. When $v^{2}=i v^{1}$, that is $w_{j}=0$ for $j=1, \ldots, n$, we get

$$
F_{k, \ell}\left(\nu^{1}, i \nu^{1}\right)=f\left(\nu^{1}\right) \quad \text { for } \nu^{1} \in \operatorname{ch}(C(q+1))
$$

so that (5.14) is certainly satisfied.
Now we choose the functions $U_{j}^{\ell}$ with a suitable bound such that $F_{k, \ell}$ is holomorphic, that is such that $\bar{\partial} \mathrm{F}_{\mathrm{k}, \ell}=0$. We can write $\mathrm{F}_{\mathrm{k}, \ell}$ in a different way, namely let $H_{0}^{\ell}\left(\nu^{1}\right)=f\left(\nu^{1}\right)$ and let

$$
\begin{aligned}
& H_{j}^{\ell}\left(v^{1} ; v_{1}^{2}, \ldots, v_{j}^{2}\right)= \\
& =\psi_{k}\left(w_{j}\right) H_{j-1}^{\ell}\left(v^{1} ; v_{1}^{2}, \ldots, v_{j-1}^{2}\right)-w_{j} U_{j}^{\ell}\left(v^{1} ; v_{1}^{2}, \ldots, v_{j}^{2}\right)
\end{aligned}
$$

for $j=1, \ldots, n$ successively, then $H_{n}^{\ell}=F_{k, \ell}$. If $H_{j-1}^{\ell}$ is holomorphic for $\left(\operatorname{Im} v_{1}^{1}, \ldots, \operatorname{Im} v_{n}^{1}, \operatorname{Im} v_{1}^{2}, \ldots, \operatorname{Im} v_{j-1}^{2}, \operatorname{Re} v_{1}^{1}, \ldots, \operatorname{Re} v_{n}^{1} ; \operatorname{Re} v_{1}^{1}, \ldots, \operatorname{Re} v_{j-1}^{1}\right.$, $\left.\operatorname{Re} v_{1}^{2}, \ldots, \operatorname{Re} v_{j-1}^{2}\right) \in C(q+1)_{\delta, j-1}^{n} \times_{\ell} \mathbb{R}^{2(j-1)} \xlongequal{{ }^{\text {def }}} B_{j-1} \subset \mathbb{C}^{n+j-1}$, which is true for $j=1$ by (5.18), then $H_{j}^{\ell}$ is holomorphic in $B_{j}$ when $U_{j}^{\ell}$ satisfies

$$
\begin{equation*}
\bar{\partial} U_{j}^{\ell}=H_{j-1}^{\ell}\left(v^{1} ; v_{1}^{2}, \ldots, v_{j-1}^{2}\right) \psi^{\prime}\left(w_{j}\right) / w_{j} n g_{j}^{\ell} . \tag{5.19}
\end{equation*}
$$

It follows when $j=n$, that $F_{k, \ell}$ is holomorphic in $B_{n}=\mathbb{R}^{2 n}+i \operatorname{ch}(C(q-1)) \subset$ $\subset \mathbb{C}^{2 n}$. Since by assumption $H_{j-1}^{\ell k, \ell}$ is holomorphic in $B_{j-1}^{n}, 1 / w_{j}$ is holomorphic outside any neighborhood of zero, $\psi^{\prime}\left(\mathrm{w}_{\mathrm{j}}\right)=0$ in a neighborhood of zero and since

$$
\begin{aligned}
& \bar{\partial} \psi^{\prime}\left(w_{j}\right)=\bar{\partial} \frac{\partial \psi}{\partial \bar{w}_{j}}\left(i v_{j}^{1}-v_{j}^{2}\right)\left[-i d \bar{v}_{j}^{1}-d \bar{v}_{j}^{2}\right]= \\
& =\left(i \frac{\partial^{2} \psi}{\partial \bar{v}_{j} \partial \bar{w}_{j}}-\frac{\partial^{2} \psi}{\partial \bar{v}_{j}^{1} \partial \overline{w_{j}}}\right) d \bar{v}_{j}^{1} \wedge d \bar{v}_{j}^{2}=\left(-i \frac{\partial^{2} \psi}{\partial \bar{w}_{j}^{2}}+i \frac{\partial^{2} \psi}{\partial \bar{w}_{j}^{2}}\right) d \bar{v}_{j}^{1} \wedge d \bar{v}_{j}^{2}=0,
\end{aligned}
$$

we get $\bar{\partial} g_{j}^{\ell}=0$. Furthermore the domain $B_{j}$ is convex, thus pseudoconvex. Therefore we can apply theorem A. 10 in order to solve (5.19). As a weightfunction we may take $\left(1+\|z\|^{2}\right)^{-3(j-1)} \exp \left(-2 \tilde{a}(z)-2 / \ell\left\|_{z}\right\|\right)$, since $\tilde{a}(z)+1 / \ell\|z\|$ is a convex function and $\log \left(1+\|z\|^{2}\right)$ is plurisubharmonic. Write $z^{j}=$ $=\left(\operatorname{Im} \nu_{1}^{1}, \ldots, \operatorname{Im} \nu_{n}^{1}\right.$, $\operatorname{Im} \nu_{n}^{2}, \ldots, \operatorname{Im} \nu_{j}^{2}$, $\left.\operatorname{Re} \nu_{j+1}^{1}, \ldots, \operatorname{Re} \nu_{n}^{1}\right) \in \mathbb{R}^{2 n}$ and $\nu[j]=$ $=\left(\nu^{1} ; \nu_{1}^{2}, \ldots, \nu_{j}^{2}\right) \in \mathbb{C}^{n+j}$ and let $\left.\lambda(\nu L j]\right)$ be the Lebesgue measure in $\mathbb{C}^{n+j}$. Then by theorem A. 10 there exists a solution $U_{j}^{\ell}$ of (5.19) with

$$
\begin{aligned}
& \int_{B_{j}}\left|U_{j}^{\ell}(\nu[j])\right|^{2} \frac{\exp \left(-2 \check{a}\left(z^{j}\right)-2 / \ell\left\|z^{j}\right\|\right)}{\left(1+\|\nu[j]\|^{2}\right)^{3 j-1}} d \lambda(\nu[j]) \leq \\
& \leq \int_{B_{j}}\left|g_{j}^{\ell}(\nu[j])\right|^{2} \frac{\exp \left(-2 \widetilde{a}\left(z^{j}\right)-2 / \ell\left\|_{z}{ }^{j}\right\|\right)}{\left(1+\|\nu[j]\|^{2}\right)^{3(j-1)}} d \lambda(\nu[j])
\end{aligned}
$$

Since $\left.\tilde{a}\right|_{C_{p}}$ can be extended to a convex homogeneous function $\tilde{a}^{\prime}$ on $\mathbb{R}^{2 n}$, see (4.5), we get for $x, y \in C_{p}$

$$
\begin{aligned}
\tilde{a}(x)-\tilde{a}(y) & =2 \tilde{a}\left(\frac{y}{2}+\frac{x-y}{2}\right)-\tilde{a}(y) \leq \tilde{a}(y)+\tilde{a}^{\prime}(x-y)-\tilde{a}(y)= \\
& =\|x-y\| a^{\prime}(\widetilde{x}-\tilde{y}) \leq A\|x-y\|
\end{aligned}
$$

and also

$$
\tilde{a}(y)-\tilde{a}(x) \leq\|x-y\| a^{\prime}(\widetilde{y-x}) \leq A\|x-y\|
$$

for some constant $A$. We set $M=\exp (2 \delta A+2 \delta / \ell)$ and $C_{k}=2\left(\delta^{2}+16 K_{k}^{2}\right) \pi M$, then we can estimate $H_{j}^{\ell}$ in terms of $H_{j-1}^{\ell}$ using $(a+b)^{2} \leq 2\left(a^{2}+b^{2}\right)$, hence $\left|w_{j}\right|^{2} /\left(1+\|\nu[j]\|^{2}\right)^{j} \leq 2$,

$$
\begin{aligned}
& \int_{B_{j}}\left|H_{j}^{\ell}(\nu[j])\right|^{2} \frac{\exp \left(-2 \tilde{a}\left(z^{j}\right)-2 / \ell\left\|z^{j}\right\|\right)}{\left(1+\|\nu[j]\|^{2}\right)^{3 j}} d \lambda(\nu[j]) \leq \\
& \leq 2\left\{\pi \delta^{2} M \int_{B_{j-1}}\left|H_{j-1}^{\ell}(\nu[j-1])\right|^{2} \frac{\exp \left(-2 \tilde{a}\left(z^{j-1}\right)-2 / \ell\left\|z^{j-1}\right\|\right)}{\left(1+\|\nu[j-1]\|^{2}\right)^{3(j-1)}} d \lambda(\nu[j-1])+\right. \\
& \left.+\int_{B}\left|g_{j}^{\ell}(\nu[j])\right|^{2} \frac{\exp \left(-2 \tilde{a}\left(z^{j}\right)-2 / \ell\left\|z^{j}\right\|\right)}{\left(1+\|\nu[j]\|^{2}\right)^{3(j-1)}} d \lambda(\nu[j])\right\} \leq \\
& \leq C_{k} \int_{j-1}\left|H_{j-1}^{\ell}(\nu[j-1])\right|^{2} \frac{\exp \left(-2 \check{a}\left(z^{j-1}\right)-2 / \ell\left\|z^{j-1}\right\|\right)}{\left(1+\|\nu[j-1]\|^{2}\right)^{3(j-1)}} d \lambda(\nu[j-1]) .
\end{aligned}
$$

Since for $j=n \quad B_{n}=\mathbb{R}^{2 n}+i \operatorname{ch}(C(q+1)), H_{n}^{\ell}=F_{k, \ell}, \nu[n]=\left(\nu^{1} ; \nu^{2}\right)=\nu$, $z^{n}=\operatorname{Im} \nu$ and for $j=0 \quad B_{0} \subset C(p), H_{0}^{l}=f, \nu[0]=\nu^{l}=z \in \mathbb{C}^{n}, z^{0}=(y, x)=z$, it follows that

$$
\begin{aligned}
& \quad \int_{\mathbb{R}^{2 n}+i C(q+1)}\left|F_{k, \ell}(v)\right|^{2} \frac{\exp (-2 a(\operatorname{Im} v)-2 / \ell\|\operatorname{Im} v\|)}{\left(1+\|v\|^{2}\right) 3 n} d \lambda(\nu) \leq \\
& \leq C_{k}^{n} \int_{C(p)}|f(z)|^{2} \exp (-2 a(z)-2 / \ell\|z\|) d \lambda(z) .
\end{aligned}
$$

According to condition $\mathrm{HS}_{2}$ (see [14] G.7) we can estimate the sup-norm by the $L^{2}$-norm and we find that (5.16) is satisfied with $m=3 n$, since also Exp can be written as projective limit of Hilbert spaces.

REMARK 5.1. If we could choose for all $k$ the functions $F_{k, \ell}$ satisfying (5.15) and (5.16) such that $F_{k+1, \ell}(\nu)=F_{k, \ell}(\nu)$ for $\nu \in \mathbb{R}^{2 n+i C(k), \text { thus if there is one }}$ function $F_{\ell}$ holomorphic in $\mathbb{R}^{2 n}+i C$ satisfying for all $k$

$$
\begin{equation*}
\left|F_{\ell}(\nu)\right| \leq M_{k}(1+\|\nu\|)^{m} \exp (\tilde{a}(\operatorname{Im} v)+1 /(\ell+1)\|\operatorname{Im} v\|), \quad \operatorname{Im} v \in C(k) \tag{5.20}
\end{equation*}
$$ for some $m$ and

$$
\begin{equation*}
F_{\ell}(z, i z)=f(z) \quad \text { for } z \in C \tag{5.21}
\end{equation*}
$$

then $\mathrm{F}_{\ell}$ would belong to $\mathrm{H}^{*}(\mathrm{a}+1 /(\ell+1) ; \mathrm{C})$ and we would have according to theorem 9.1 [14]

$$
F_{\ell}(\nu)=\left\langle\mu_{\xi, \eta}^{\ell}, e^{i \nu^{1} \cdot \xi+i \nu^{2} \cdot \eta}\right\rangle
$$

with $\mu^{\ell} \in S^{*}(a+1 /(\ell+1) ; C)^{\prime}$. In that case

$$
f(z)=\left\langle\mu_{\zeta}^{\ell}, e^{i z \cdot \zeta}\right\rangle \quad \text { for } z \in C
$$

and since $\delta_{\ell} / \ell \| \zeta^{\|}$is uniformly continuous on $\Omega_{\ell}$, we get for $\phi \in \mathrm{A}_{\ell} ; \Omega_{\ell}\left(\mathrm{L}_{\ell}\right)$


$$
\leq \mathrm{K}_{\ell} \sup _{\zeta \in \Omega_{\ell}} \exp \left(\delta_{\ell} \cdot / \ell\|\zeta\|\right)|\phi(\zeta)|
$$

where $\Omega_{\ell}$ is now given by (5.5)(ii). Hence $F$ would be a surjective map from $A_{i i}^{\prime}$ onto Exp.
6. FORMULATION OF THE PROBLEMS AND STATEMENT OF THE MAIN RESULT

In this section we reconsider the procedure followed in the last section and we formulate the problems to be solved in order that $\bar{A}_{i}$ and $\bar{A}_{i i}$ given in (5.11) are indeed also given by (5.10), that is the space of $a Z Z$ holomorphic functions satisfying the growth conditions in the sets int $\Omega_{k}$, where $\Omega_{k}$ is given by (5.5) (i) or (5.5) (ii), respectively. Theorem 5.1 says
that $F$ is surjective from $A_{i}^{\prime}$ onto $\operatorname{Exp}$, from which it follows that $F^{t}$ (5.1) defined by (5.6) is injective. Since $\bar{A}_{i}$ is reflexive (see (5.11)) the linear hull of the set $\left\{e^{i z \cdot \zeta}\right\}_{\zeta \in \Omega}$ is dense in $\operatorname{Exp}(a+0 ; C)$. On the other hand when the above problems are solved, we see that the linear hull of the set $\left\{e^{i z \cdot \zeta}\right\}_{z \in C}$ is dense in the space (5.10) and the main result (theorem 6.1) follows.

Anticipating on the results we will get we mean in this section by $A_{i}^{\prime}$ and $A_{i i}^{\prime}$ the dual $A^{\prime}$ of the spaces $\bar{A}_{i}$ and $\bar{A}_{i i}$, respectively, given by (5.10)

$$
\begin{equation*}
\bar{A}=\underset{k \rightarrow \infty}{\operatorname{ind} \lim A_{\infty}}\left(M_{k} ; \Omega_{k}\right), \tag{6.1}
\end{equation*}
$$

where $\Omega_{k}$ is given in (5.5)(i) or (5.5)(ii), respectively. Then $\bar{A}$ is an LS-space. There are also other possibilities of writing $\bar{A}_{i}$ as an inductive limit of spaces $A_{i}^{k}$, or $\bar{A}_{i i}$ as inductive limit of spaces $A_{\ell}$. We will choose appropriate spaces $A_{i}^{k}$ and $A_{\ell}$. In the above $A_{i}^{k}=A_{\infty}\left(M_{k} ; \Omega_{k}\right)$.

In the last section we have embedded (a linear subspace, namely (5.11), of) $A_{i}^{k}$ into the space $S^{\star k+1}\left(\Omega_{k+1}\right)$. Roughly we can say that $A_{i}^{k}$ consists of those elements $\phi$ in $S^{\star k+1}\left(\Omega_{k+1}\right)$ with $\bar{\partial} \phi=0$ and that any element $\mu$ of $S^{*}(\Omega)^{\prime}$ that satisfies $\mu=\bar{\partial}^{t} \sigma_{k}$ for some $\sigma_{k} \in\left(S^{k *}\left(\Omega_{k}\right)^{\prime}\right)^{n}$, is zero when restricted to $A_{i}^{k}$. Hence the elements of $A_{i}^{\prime}$ can be identified with the equivalence classes of the elements in $S^{*}(\Omega)$ ', when two elements in $S^{*}(\Omega)$ ' are equivalent if their difference $\mu$ can be written as $\mu=\bar{\partial}^{t} \sigma_{k}$ in each $S^{\star k}\left(\Omega_{k}\right)^{\prime}$ for some $\sigma_{k} \in\left(S^{\star k}\left(\Omega_{k}\right)^{\prime}\right)^{n}$. Now we investigate this more precisely.

First we write $\bar{A}_{i}$ as inductive limit of spaces having the topology of $S^{k *}\left(\Omega_{k}\right)$, that is $A_{i}^{k}$ now is the closed linear subspace of

$$
s^{k} \xlongequal{\text { not }} s^{k *}\left(a+1 / k ; C_{k+2}\right) \xlongequal{\text { def }} \underset{m \rightarrow \infty}{\operatorname{proj}} 1 \operatorname{im} W_{\infty}^{m}, 0\left((1+\|\zeta\|)^{m} \exp \left(\delta_{k} \frac{1}{k}\|\zeta\|\right) ; \Omega_{k}\right)
$$

consisting of the functions holomorphic in int $\Omega_{k}$ and $C^{\infty}$ on $\Omega_{k}$ with the topology heredited from $\mathrm{S}^{\mathrm{k}}$. Therefore, according to [15] prop. 35.5 (a) the following sets can be identified

$$
\left(A_{i}^{k}\right)^{\prime}=\left(S^{k}\right)^{\prime} /\left(A_{i}^{k}\right)^{0}
$$

and according to [15] prop. 35.6 this is also true for the topologies, when we provide these spaces with the weak ${ }^{*}$ topology and the quotient topology with respect to the weak topology, respectively:

$$
\left(A_{i}^{k}\right)_{\sigma}^{\prime} \cong\left(S^{k}\right)_{\sigma}^{\prime} /\left(A_{i}^{k}\right)^{0}
$$

On the other hand $A_{i}^{k}$ is the kernel of the map $\bar{\partial}=\left(\partial / \partial \bar{\zeta}_{1}, \ldots, \partial / \partial \bar{\zeta}{ }_{n}\right)$

$$
\bar{\partial}: S^{k} \longrightarrow\left(S^{k}\right)^{n}
$$

so that according to [15] prop. $35.4\left(\mathrm{~A}_{\mathrm{i}}^{\mathrm{k}}\right)^{0}$ is the weak* closure in $\left(\mathrm{S}^{\mathrm{k}}\right)_{\sigma}^{\prime}$ of Im $\bar{\partial}^{t}$. Since $S^{k}$ is reflective (it is an $\bar{F} \bar{S}$-space), the weak ${ }^{*}$ closure of $\operatorname{Im} \bar{\partial}^{\mathrm{t}}$ is equal to the closure in the strong topology in ( $\left.\mathrm{S}^{\mathrm{k}}\right)^{\prime}$, [15] prop. 35.2. We denote the closure in $\left(S^{k}\right)^{\prime}$ of the range of the map

$$
\bar{\partial}^{-t}:\left(\left(S^{k}\right)^{\prime}\right)^{n} \rightarrow\left(S^{k}\right)^{\prime}
$$

by $\overline{\operatorname{Im} \bar{\partial}_{k}^{t}}$. So we get

$$
\left(A_{i}^{k}\right)_{\sigma}^{\prime} \cong\left(S^{k}\right)_{\sigma}^{\prime} / \frac{}{\operatorname{Im} \bar{\partial}_{k}^{t}}
$$

Finally we will obtain isomorphisms also for the strong topologies.
Therefore, we consider this spaces only with the topology of weakly* converging sequences denoted by $\left(A_{i}^{k}\right)_{\sigma, s}^{\prime}$ and $\left(S^{k}\right)_{\sigma, s}^{\prime}$. Since $S^{k}$ is a Montel space we get

$$
\left(S^{k}\right)_{\sigma, s}^{\prime}=\left(S^{k}\right)_{b, s}^{\prime}
$$

where $\left(S^{k}\right)_{b}^{\prime}$ is the dual of $S^{k}$ provided with the strong topology. According to [14] theorem 9.1 ( 9.6 ) and page $615^{\text {th }}$ line from below) the Fourier transformation $F$ maps $\left(S^{k+1}\right)_{b}^{\prime}$ continuously into

$$
H^{k}=\underset{m \rightarrow \infty}{\operatorname{ind}} \lim _{H^{m *}}\left(a+1 /(k+1) ; C_{k}, k\right)
$$

with

$$
H^{m *}\left(a+1 / k ; C_{k}, k\right) \xlongequal{\text { def }} A_{\infty}\left((1+\|v\|)^{-m} \exp (-a ̃(\operatorname{Im} v)-1 /\|\operatorname{Im} v\|) ; \mathbb{R}^{\left.2 n_{+i} C(k)\right)}\right.
$$

and $F^{-1}$ maps $H^{k}$ continuously into $S^{p *}\left(a+1 /(k+1) ; C_{k}\right)$, when $1 / k<\delta_{p} / \mathrm{p}$, hence into ( $S^{\mathrm{P}}$ ) ${ }_{b}$.

Let $\mathrm{W}=\left(\mathrm{w}_{1}, \ldots, \mathrm{w}_{\mathrm{n}}\right)$, where $\mathrm{w}_{\mathrm{j}}=\mathrm{i} v_{\mathrm{j}}^{1}-v_{\mathrm{j}}^{2}, \mathrm{j}=1, \ldots, \mathrm{n}$ and let $\mathrm{W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{k}}$ be the subspace of $H^{k}$ consisting of functions $f(\nu)$ that can be written as

$$
f(v)=\sum_{j=1}^{n} w_{j} g_{j}(v)
$$

with $g_{j} \in H^{k}, j=1, \ldots, n$. Then

$$
F \overline{\operatorname{Im} \bar{\partial}_{\mathrm{k}+1}^{\mathrm{t}}} \subset \overline{\mathrm{~W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{k}}} \text { and } F^{-1} \cdot \mathrm{~W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{k}-1} \subset \overline{\operatorname{Im} \bar{\partial}_{\mathrm{p}}^{\mathrm{t}}}
$$

when $1 /(k-1)<\delta_{p} / p$. Furthermore

$$
\overline{\mathrm{W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{k}}} \subset \mathrm{~W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{k}-1}
$$

for, let $f_{\alpha} \in W \cdot \vec{H}^{k}$ be a Cauchy net converging to $f \in H^{k}$. Then $f_{\alpha}=W \cdot \vec{g}_{\alpha}$ with $\vec{g}_{\alpha} \in\left(H^{k}\right)^{n}$, so that $f_{\alpha}$ and hence $f$ vanishes on

$$
v_{k}=\left\{\mathbb{R}^{2 n}+i C(k)\right\} \cap\left\{\dot{v} \mid i v_{j}^{1}-v_{j}^{2}=0, j=1, \ldots, n\right\} .
$$

The inclusion follows if we have shown
PROBLEM 6.1. A function $\mathrm{f} \in \mathrm{H}^{\mathrm{k}}$ vanishing on $\mathrm{V}_{\mathrm{k}}$ can be written as

$$
\mathrm{f}(v)=\mathrm{w} \cdot \overrightarrow{\mathrm{~g}}^{\mathrm{k}-1}(v), \quad v \in \mathbb{R}^{2 \mathrm{n}}+\mathrm{iC}(\mathrm{k}-1)
$$

with $\overrightarrow{\mathrm{g}}^{\mathrm{k}-1} \in\left(\mathrm{H}^{\mathrm{k}-1}\right)^{\mathrm{n}}$; in particular there is a positive N such that $\mathrm{g}^{\mathrm{k}-1} \epsilon$ $\epsilon H^{m+N *}\left(a+1 / k ; C_{k-1}, k-1\right)$, when $f \in H^{m *}\left(a+1 / k ; C_{k}, k\right)$.

In that case we have
where $\left(H^{k}\right)_{S}$ means that $H^{k}$ is provided with the topology of convergent sequences.

Furthermore weakly ${ }^{*}$ converging sequences in $\left(A_{i}^{k+1}\right)^{\prime}$ converge strongly in $\left(A_{i}^{k}\right)^{\prime}$ because $A_{i}^{k}$ is an LS-space, so that

$$
\underset{k \rightarrow \infty}{\operatorname{proj}} \lim \left(A_{i}^{k}\right)_{\sigma, s}^{\prime}=\underset{k \rightarrow \infty}{\operatorname{proj}} \lim \left(A_{i}^{k}\right)_{b, s}^{\prime}=\left(A_{i}\right)_{b, s}^{\prime}=A_{i}^{\prime},
$$

where the last equality follows from the fact that the topology in the metric space (namely the $\overline{\mathrm{S}}$-space) $A_{i}^{\prime}$ is determined by convergent sequences. Thus $F$ is an isomorphism between


In case (ii) when $\Omega_{k}$ is defined by (5.5)(ii) we define S-spaces with $L^{2}$-norms rather than with sup-norms. For, a continuous map from one Hilbert space into another is weakly compact, so that projective and inductive limits of sequences of Hilbert spaces (called FS*-spaces and DFS*-spaces, respectively) are reflexive and they are dual to each other (see [19]). Also we apply th. 15 of [19], where the isomorphism holds for the strong topologies and not only for the weak ${ }^{*}$ topologies as in [15] prop. 35.6. Hence we do not have to restrict ourselves to the topology of convergent sequences (this also applies to case (i)).

Let

$$
S_{\ell}^{\mathrm{m}}=\underset{\mathrm{k} \rightarrow \infty}{\operatorname{ind}} \lim W_{2}^{\mathrm{m}}\left((1+\|\zeta\|)^{\mathrm{m}} \exp \delta_{\mathrm{k}} 1 / \mathrm{k}\|\zeta\| ; \Omega_{\ell}\right)
$$

and let $A_{\ell}^{m}$ be the closed linear subspace $S_{\ell}^{m} \cap A\left(\Omega_{\ell}\right)$. In virtue of [19] th. 7 the topology of $\operatorname{ind} \lim _{k \rightarrow \infty}^{m}(k)$ with

$$
A_{\ell}^{m}(k) \stackrel{\text { def }}{=} A\left(\Omega_{\ell}\right) \cap W_{2}^{m}\left((1+\|\zeta\|)^{m} \exp \delta_{k} 1 / k\|\zeta\| ; \Omega_{\ell}\right)
$$

if finer than the topology in $A_{\ell}^{m}$ induced by $S_{\ell}^{m}$, but one easily sees that it is also less fine, hence $A_{\ell}^{m}$ is a DFS*-space. Since $\bar{A}_{i i}$ can also be described by $L^{2}$-norms ([14] G.7) we get

$$
\overline{\mathrm{A}}_{\mathrm{ii}}=\underset{\ell \rightarrow \infty}{\operatorname{ind} \lim _{\ell},}
$$

where $A_{\ell}=\underset{\mathrm{m}}{\mathrm{m} \rightarrow \infty} \underset{\ell}{ } \lim _{\ell}^{\mathrm{m}}$. Although $\overline{\mathrm{A}}_{\mathrm{ii}}$ is an LS-space, this fact is not expressed by the above inductive limit, which is only a weakly compact sequence. Indeed, a neighborhood of zero in $A_{\ell}$ is bounded in $A_{\ell+1}$, hence relatively weakly compact in each $A_{\ell+1}^{m}(m=0,1,2, \ldots)$, since $A_{l}^{m}$ is reflexive ([15] prop. 36.3), and thus relatively weakly compact in $A_{\ell+1}$. Therefore, $\left(A_{i i}\right)_{0}^{\prime}=\operatorname{proj} \underset{\ell \rightarrow \infty}{ } \lim _{\ell}\left(A_{\ell}\right) b_{b}^{\prime}$. However, the projective limit in $A_{\ell}$ has no nice properties, so we are forced to consider the weak ${ }^{*}$ topology in ( $\left.\mathrm{A}_{\mathrm{i}}\right)^{\prime}$ '. The topology of ( $\left.\mathrm{A}_{\mathrm{ii}}\right)^{\prime}$ is also determined by weakly* converging sequences, hence $\left(A_{i i}\right)_{b}^{\prime}=\left(A_{i i}\right)_{\sigma, s}^{\prime}=\operatorname{proj} j_{l \rightarrow \infty} 1 \mathrm{im}\left(A_{\ell}\right)_{\sigma, s}^{\prime}$. Any weakly* converging sequence in $\left(A_{\ell}\right)_{\sigma, s}^{\prime}$ converges weakly* in $\left(A_{l}^{m}\right) \sigma_{\sigma, s}^{\prime}$ for some $m$ and thus it converges weak$1 y^{*}$ in each $A_{l}^{m}(k), k=1,2, \ldots$. Since the embedding map from $A_{l}^{m}(k)$ into $A_{\ell+1}^{m}(k+1)$ is compact according to [14] G.7, the sequence converges strongly in $\left(A_{l-1}^{\mathrm{m}}\right)^{\prime}([14] \mathrm{E} .2)$, thus in ind $\lim _{\mathrm{m} \rightarrow \infty} \operatorname{im}\left(\mathrm{A}_{\ell-1}^{\mathrm{m}}\right)_{b}^{\prime}, \mathrm{s}^{\prime}$. Since $\left(\mathrm{A}_{\ell}^{\mathrm{m}}\right)^{\prime}$ is a Fréchet space, namely an $F S^{*}$-space, $\left(A_{\ell-1}^{m}\right)_{b, s}^{\prime}=\left(A_{\ell-1}^{m}\right)_{b}^{\prime}$, so

$$
\left(\mathrm{A}_{\mathrm{ii}}\right)_{\mathrm{b}}^{\prime}=\underset{\ell \rightarrow \infty}{\operatorname{proj}} \lim \underset{\mathrm{m} \rightarrow \infty}{\operatorname{ind}} \lim _{l}\left(\mathrm{~A}_{\ell}^{\mathrm{m}}\right)_{\mathrm{b}}^{\prime}
$$

Now $A_{l}^{m+1}$ is the kernel of the continuous map

$$
\bar{\partial}_{\mathrm{m}}: s_{\ell}^{\mathrm{m}+1} \rightarrow\left(\mathrm{~s}_{\ell}^{\mathrm{m}}\right)^{\mathrm{n}} .
$$

In virtue of [19] th. 15 and [15] prop. 35.4 we get

$$
\left(A_{\ell}^{\mathrm{m}+1}\right)^{\prime} \cong\left(S_{\ell}^{\mathrm{m}+1}\right)^{\prime} / \frac{}{\operatorname{Im} \bar{\partial}_{\mathrm{m}}^{\mathrm{t}}}
$$

where the closure in $\left(S_{\ell}^{\mathrm{m}+1}\right)^{\prime}$ of $\operatorname{Im} \bar{\partial}_{\mathrm{m}}^{\mathrm{t}}$ equals the weak ${ }^{*}$ closure, since $S_{\ell}^{\mathrm{m}+1}$ is reflexive.

Let

$$
\mathrm{H}_{\ell}^{\mathrm{m}}=\underset{\mathrm{k} \rightarrow \infty}{\operatorname{proj}} \lim \mathrm{H}^{\mathrm{m} *}\left(\mathrm{a}+1 / \ell ; \mathrm{C}_{\mathrm{k}}, \mathrm{k}\right)
$$

and let

$$
\mathrm{H}_{\ell}^{*}=\underset{\mathrm{m} \rightarrow \infty}{\operatorname{ind}} \lim H_{\ell}^{\mathrm{m}} .
$$

Then it follows from [14] th.9.1 (using D. 2 instead of G.5) and G. 3 that

$$
F\left(S_{\ell}^{\mathrm{m}}\right)^{\prime} \subset \mathrm{H}_{\ell}^{\mathrm{m}+\mathrm{n}+1} \text { and } F^{-1} \mathrm{H}_{\ell+1}^{\mathrm{m}} \subset\left(\mathrm{~S}_{\ell}^{\mathrm{m}+2 \mathrm{n}+2}\right)^{\prime} .
$$

As in case (i) problem 6.1 and the following problem imply that

$$
\begin{aligned}
\overline{\mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{m}}} & \subset \bigcap_{\mathrm{k}=1}^{\infty} \overline{\mathrm{W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{m} \star}}\left(\mathrm{a}+1 / \ell ; \mathrm{C}_{\mathrm{k}+1}, \mathrm{k}+1\right) \subset \bigcap_{\mathrm{k}=1}^{\infty} \mathrm{W} \cdot \overrightarrow{\mathrm{H}}^{\mathrm{m}+1+\mathrm{N} *}\left(\mathrm{a}+1 / \ell ; \mathrm{C}_{\mathrm{k}}, \mathrm{k}\right) \subset \\
& \subset \mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{M}},
\end{aligned}
$$

where the closures are taken in the corresponding spaces with $\mathrm{m}+1$ instead of $m$ and where $M>m+1+N$.

PROBLEM 6.2. When a function $\mathrm{f} \in \mathrm{A}\left(\mathbb{R}^{2 \mathrm{n}}+\mathrm{iC}\right)$ in each $\mathbb{R}^{2 \mathrm{n}}+\mathrm{iC}(\mathrm{k})$ can be written as

$$
\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{~g}}_{\mathrm{k}} \quad \text { for some } \overrightarrow{\mathrm{g}}_{\mathrm{k}} \in\left(\mathrm{H}^{\mathrm{m} *}\left(\mathrm{a}+1 / \ell ; \mathrm{C}_{\mathrm{k}}, \mathrm{k}\right)\right)^{\mathrm{n}},
$$

(obviously in that case $\mathrm{f} \in \mathrm{H}_{\ell}^{*}$ ), then f can be written in $\mathbb{R}^{2 n}+\mathrm{i} C$ as

$$
\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{~g}} \quad \text { with } \overrightarrow{\mathrm{g}} \in\left(\mathrm{H}_{\ell}^{*}\right)^{\mathrm{n}} ;
$$

in particular $\overrightarrow{\mathrm{g}} \in\left(\mathrm{H}_{\ell}^{\mathrm{m}+\mathrm{N}}\right)^{\mathrm{n}}$ for some positive N independent of f .
Hence $H_{\ell}^{\mathrm{m}+1} / \overline{\mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{m}}} \rightarrow \mathrm{H}_{\ell}^{M} / \mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{M}} \rightarrow \mathrm{H}_{\ell}^{\mathrm{M}+1} / \overline{\mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{M}}}$ and thus $F$ is an isomorphism between

$$
\begin{equation*}
F: A_{i i}^{\prime} \longleftrightarrow \operatorname{proj} \lim \quad \underset{\ell \rightarrow \infty}{\text { ind }} \lim \mathrm{H}_{\ell}^{\mathrm{m}} / \mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{m}} \tag{6.3}
\end{equation*}
$$

Let $H_{i i}^{*}=\operatorname{pro}_{\mathfrak{j} \rightarrow \infty}^{\operatorname{j}} \lim _{\ell}^{*} \mathrm{H}_{\ell}^{*}$ and let $H_{i}^{*}=\underset{k \rightarrow \infty}{\operatorname{proj}} 1 \mathrm{im} H^{k}$. Then the right hand side in (6.2) is equal to the equivalence classes of elements in $H_{i}^{*}$, when two elements in $H_{i}^{*}$ are equivalent if their difference $f$ can be written in each $\mathbb{R}^{2 n}+\mathrm{iC}(\mathrm{k})$ as $\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{g}}_{\mathrm{k}}$ with $\overrightarrow{\mathrm{g}}_{\mathrm{k}} \in\left(\mathrm{H}^{\mathrm{k}}\right)^{\mathrm{n}}$. The right hand side in (6.3) is equal to the equivalence classes in $H_{i i}^{*}$, when two elements in $H_{i i}^{*}$ are equivalent of their difference $f$ for each $\ell$ can be written in $\mathbb{R}^{2 n}+i C$ as $f=W \cdot \vec{g}_{\ell}$ with $\overrightarrow{\mathrm{g}}_{\ell} \in\left(\mathrm{H}_{\ell}^{*}\right)^{n}$.

Next we consider the set in $\mathbb{R}^{2 n}+i C$ where $W=\left(w_{1}, \ldots, w_{n}\right)$ vanishes, namely

$$
v=\left\{v \mid v \in \mathbb{R}^{2 n}+i C, i v v_{j}^{1}-v_{j}^{2}=0, j=1, \ldots, n\right\}
$$

and

$$
v_{k}=V \cap\left\{\mathbb{R}^{\left.2 n_{+i} C(k)\right\} .}\right.
$$

Since $W \cdot \vec{H}^{k}$ vanishes on $V_{k}$ and $W \cdot \vec{H}_{\ell}^{*}$ vanishes on $V$, we can define the continuous restriction maps $I$

$$
I_{i}: H^{k} /\left.W \cdot \vec{H}^{k} \longrightarrow H^{k}\right|_{V_{k}}
$$

and

$$
\mathrm{I}_{\mathrm{ii}}: \mathrm{H}_{\ell}^{*} /\left._{\mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{*}} \longrightarrow \dot{\mathrm{H}}_{\ell}^{*}\right|_{\mathrm{V}}
$$

by $\operatorname{If}(\nu)=f\left(\nu^{1}, i \nu^{1}\right)$. Here $H^{k} \mid V_{k}$ and $\left.H_{\ell}^{\star}\right|_{V}$ are the spaces of restrictions of functions in $H^{k}$ or $H_{l}^{*}$ to $V_{k}$ or $V$ with the topology induced by $H^{k}$ or $H_{\ell}^{*}$, respectively. Then the maps $I$ are surjective. When we regard $z=\nu^{1}$ as the variable in $V$ there are natural continuous injections $J$

$$
J_{i}:\left.H^{k}\right|_{V_{k}} \longleftrightarrow A_{\infty}(\exp (-\tilde{a}(z)-1 / k\|z\|) ; C(k))
$$

and

Now the topologies of

$$
\left.\underset{\mathrm{k} \rightarrow \infty}{\operatorname{proj} \lim _{\mathrm{H}}}\right|_{V_{\mathrm{k}}}=\left.\mathrm{H}_{\mathrm{i}}^{*}\right|_{V} \quad \text { and }\left.\quad \underset{\ell \rightarrow \infty}{\operatorname{proj}} \lim _{\ell}^{*}\right|_{V}=H_{i i}^{*} \mid V
$$

become extremely simple, as they both coincide with the one induced by $\operatorname{Exp}(a+0 ; C)$, compare (2.14), thus $H_{i}^{*}\left|V=H_{i i}^{*}\right| V$ is an $F \bar{S}-$ space. Finally we have obtained

and
(6.5) $A_{i i}^{!} \stackrel{F}{p} \underset{\ell \rightarrow \infty}{\operatorname{proj}} \lim \underset{\mathrm{~m} \rightarrow \infty}{\operatorname{ind}} \lim \mathrm{H}_{\ell}^{\mathrm{m}} / \mathrm{W} \cdot \overrightarrow{\mathrm{H}}_{\ell}^{\mathrm{m}} \xrightarrow{\mathrm{I}_{\mathrm{ii}}} \underset{\ell \rightarrow \infty}{\operatorname{proj}} \underset{\mathrm{m} \rightarrow \infty}{\lim } \underset{\ell}{\text { ind }} \lim \mathrm{H}_{\ell}^{\mathrm{m}} \mid V$

$$
\xrightarrow{\mathrm{J}_{\mathrm{ii}}} \operatorname{Exp}(\mathrm{a}+0 ; \mathrm{C}) .
$$

In theorem 5.1 we have proved that the map $J_{i} \circ I_{i}$ is surjective, hence $J_{i}$ is surjective. Remark 5.1 is concerned with the question whether $J_{i i} \circ I_{i i}$ is surjective. Using the proof of theorem 5.1 we see that indeed $J_{i i}{ }^{\circ} I_{i i}$ and hence $J_{i i}$ are surjective, if the following problem is solved.

PROBLEM 6.3. Let the function $\mathrm{f}_{\mathrm{k}} \in \mathrm{H}^{\mathrm{m}}\left(\mathrm{a}+1 / \ell ; \mathrm{C}_{\mathrm{k}}, \mathrm{k}\right)$ satisfy for $a l Z$ $k=1,2, \ldots$

$$
\left.\left(f_{k+1}-f_{k}\right)\right|_{V_{k}}=0
$$

Then there exists a function $\mathrm{f} \in \mathrm{H}_{\ell}^{*}$ with for all $\mathrm{k}=1,2, \ldots$

$$
\left.\left(f-f_{k}\right)\right|_{V_{k}}=0
$$

For, in that case (5.20) and (5.21) can be satisfied.
Problem 6.1 says that the map $I_{i}$ is injective. Problem 6.1 as well as problem 6.2 follow from problem 6.4, which says that the map $I_{i i}$ is injective.

PROBLEM 6.4. A function $\mathrm{f} \in \mathrm{H}_{\ell}^{*}$ vanishing on V can be written in $\mathbb{R}^{2 \mathrm{n}}+\mathrm{iC}$ as

$$
\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{~g}} \quad \text { with } \quad \overrightarrow{\mathrm{g}} \in\left(\mathrm{H}_{\ell}^{*}\right)^{\mathrm{n}} .
$$

The next step is to investigate holomorphic functions vanishing on the set $V$, but before doing this we give an intuitive interpretation of the isomorphisms (6.4) and (6.5) in terms of the last section revealing the a priori difference between the spaces (5.11) and (5.10) in terms of this section. In section 5 we have shown that Exp is isomorphic to the dual of the closure (given in (5.11) and here denoted by $\widetilde{\mathbb{A}}_{\mathrm{i}}$ ) of the linear hull of $\left\{e^{i z \cdot \zeta}\right\}_{z \in C}$ in $\bar{A}_{i}$ and in this section $\operatorname{Exp}$ is isomorphic to $\left.H_{i}^{*}\right|_{V}$. Hence $\left(\left.H_{i}^{*}\right|_{V}\right)^{\prime}$ is isomorphic to $\tilde{A}_{i}$ and problem 6.3 implies the same for $\tilde{A}_{i i}$. Indeed, let us examine what elements of ( $H^{*}$ )' yield $\widetilde{A}$ or $\bar{A}$ under $F^{t}$ defined in (5.6) (we do not distinguish between cases (i) and (ii) here). Let $\phi \in \widetilde{A}$, thus

$$
\phi(\zeta)=\sum_{k=1}^{\infty} c_{k} e^{i z_{k} \cdot \zeta}
$$

with $z_{k} \in C$ and with some constants $c_{k}$. If for some $\sigma \in\left(H^{*}\right)^{\prime}$

$$
\left\langle\sigma_{\nu}, \mathrm{e}^{\mathrm{i} \nu^{1} \cdot \xi+\mathrm{i} \nu^{2} \cdot n_{>}}=\phi(\zeta)\right.
$$

then

$$
\begin{aligned}
\sigma_{v} & =\sum_{k=1}^{\infty} c_{k} \delta\left(v^{1}-z_{k}\right) \delta\left(\nu^{2}-i z_{k}\right)=\sum_{k=1}^{\infty} c_{k} \delta\left(\nu^{1}-z_{k}\right) \delta\left(\nu^{2}-i \nu^{1}\right)= \\
& =\delta\left(i \nu^{1}-\nu^{2}\right) \sum_{k=1}^{\infty} c_{k} \delta\left(\nu^{1}-z_{k}\right)
\end{aligned}
$$

thus $\sigma$ acts on the restrictions of functions in $H^{*}$ to $V$, that is $\sigma \in\left(\left.H^{*}\right|_{V}\right)^{\prime}$. Now consider an element $\phi \in \overline{\mathrm{A}}$. If for some $\sigma \in\left(\mathrm{H}^{*}\right)^{\prime}$

$$
\left\langle\sigma_{\nu}, e^{i \nu^{1} \cdot \xi+i \nu^{2} \cdot n_{>}}=\phi(\zeta),\right.
$$

then we only know that

$$
\left\langle\sigma_{\nu}, w_{j} e^{i \nu^{1} \cdot \xi+i \nu^{2} \cdot n_{>}}=0, \quad j=1, \ldots, n\right.
$$

since $\bar{\partial}_{j} \phi=0$. The exponentials are dense in $H^{*}$, so that $\left\langle\sigma, W \cdot \overrightarrow{\mathrm{H}}^{*}\right\rangle=0$, thus

$$
\sigma \in\left(\mathrm{H}^{*} / \mathrm{W} \cdot \overrightarrow{\mathrm{H}}^{\star}\right)^{\prime},
$$

see [15] prop.35.5(b). When we have shown that the map I is injective (problem 6.4), the spaces $\widetilde{A}(5.11)$ and $\bar{A}(5.10)$ coincide and we obtain a theorem similar to theorem 2.1.

V is defined as the simultaneous zero-set of the polynomials $\mathrm{w}_{1}=$ $=i \nu_{1}^{1}-v_{1}^{2}, \ldots, \mathrm{w}_{\mathrm{n}}=\mathrm{i} \nu_{\mathrm{n}}^{1}-\nu_{\mathrm{n}}^{2}$. These polynomials generate a prime ideal in any point of a pseudoconvex set $\Omega \subset \mathbb{C}^{2 n}$. Therefore, according to Hilbert's Nullstellensatz all holomorphic functions $f$ in $\Omega$ vanishing on $V$ can locally, that is in a neighborhood $\omega$ of any point in $\Omega$, be written as

$$
\begin{equation*}
\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{~g}}_{\omega} \quad \text { with } \overrightarrow{\mathrm{g}}_{\omega} \in \mathrm{A}(\omega)^{\mathrm{n}} \tag{6.6}
\end{equation*}
$$

see appendix (A.18). With the aid of Cartan's theorem B (theorem A 14) it is shown in the appendix that $f \in A(\Omega)$ satisfying (6.6), satisfies (6.6) globally, that is $f$ can be written as

$$
\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{~g}} \quad \text { with } \overrightarrow{\mathrm{g}} \in \mathrm{~A}(\Omega)^{\mathrm{n}}
$$

Problem 6.4 asks for functions $g \in H_{\ell}^{*}$, so it is the analogue with estimates of the problem treated in the appendix. By (6.6) we can reformulate problem 6.4:

PROBLEM 6.5. If $\mathrm{f} \in \mathrm{H}_{\ell}^{*}$ can locally (that is in a neighborhood $\omega$ of any point in $\mathrm{T}^{\mathrm{C}}$ ) be written as

$$
\mathrm{f}=\mathrm{W} \cdot \overrightarrow{\mathrm{~g}}_{\omega}, \quad \overrightarrow{\mathrm{g}}_{\omega} \in \mathrm{A}(\omega)^{\mathrm{n}},
$$

then there exists $\vec{g} \in\left(H_{\ell}^{*}\right)^{n}$ with $f=W \cdot \vec{g}$.
In the next section we will solve this problem for general polynomial systems $P$ instead of W . Also in that case, a set V can be so defined (see EHRENPREIS [3]) that a function $f$ vanishing on $V$ can locally be written as
$f=P \cdot \vec{g}$, see theorem A 17. Provided that $J$ is surjective, the isomorphism $I$ in (6.4) and (6.5) is the analogue with estimates of the isomorphism (A 19). Using (6.6) and the above mentioned problem of the appendix (theorem A 15) we can reformulate problem 6.3:

PROBLEM 6.6. Let the functions $f_{k} \in H^{m *}\left(a+1 / l ; C_{k}, k\right)$ satisfy for alZ $\mathrm{k}=1,2, \ldots \mathrm{f}_{\mathrm{k}+1}-\mathrm{f}_{\mathrm{k}}=\mathrm{W} \cdot \overrightarrow{\mathrm{g}}_{\mathrm{k}}$ in $\mathbb{R}^{2 n_{+i}} \mathrm{C}(\mathrm{k}), \overrightarrow{\mathrm{g}}_{\mathrm{k}} \in \mathrm{A}\left(\mathbb{R}^{2 n}+\mathrm{iC}(\mathrm{k})\right)^{\mathrm{n}}$, then there exists a function $\mathrm{f} \in \mathrm{H}_{\ell}^{*}$ with for all $\mathrm{k}=1,2, \ldots \mathrm{f}-\mathrm{f}_{\mathrm{k}}=\mathrm{W} \cdot \overrightarrow{\mathrm{g}}_{\mathrm{k}}$ in $\mathbb{R}^{2 \mathrm{n}}+\mathrm{iC}(\mathrm{k})$ for some $\overrightarrow{\mathrm{g}}_{\mathrm{k}} \in \mathrm{A}\left(\mathbb{R}^{2 \mathrm{n}}+\mathrm{iC}(\mathrm{k})\right)^{\mathrm{n}}$.

Also this problem will be solved in the next section for general polynomial systems $P$ instead of $W$. Therefore, $J_{i i}$ is surjective and we have proved the main theorem of sections 5, 6 and 7, namely

THEOREM 6.1. Let a be a convex function on pr for some open convex cone C in $\mathbb{C}^{\mathrm{n}}$ and let $\Omega$ and $\Omega_{k}$ be the closed convex sets in $\mathbb{C}^{\mathrm{n}}$ determined by (5.4) and (5.5)(ii), respectively. Then the map F from A', the dual of the space $\overline{\mathrm{A}}(6.1)$, into $\operatorname{Exp}(\mathrm{a}+0 ; \mathrm{C})$ given by

$$
F(\mu)(z)=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta}\right\rangle, \quad \mu \in A^{\prime}
$$

is an isornorphism.

We have also shown that Exp is isomorphic to $A_{i}^{\prime}$, hence $\bar{A}_{i}=\bar{A}_{i i}$. Taking into account theorem 5.1 we can conclude that the linear hull of the set $\left\{e^{i z \cdot \zeta}\right\}_{z \in C}$ is dense in $\bar{A}(6.1)$ in both cases (i) and (ii).

Theorem 6.1 is a generalization to non-entire functions of the theorem of EHRENPREIS [2] and MARTINEAU [12] of section 2 which deals with entire functions as Fourier transforms. A particular case of this theorem with (5.5) (i) instead of (5.5) (ii) has already been proved by KAWAI in [9].

## 7. COMPLETION OF THE PROOFS

In this section we solve problems 6.5 and 6.6 . For that purpose cohomology with bounds is introduced. The solution requires estimates in the steps of the proof of a similar statement without bounds in the appendix. We formulate the theorems in a more general way making them useful in other
applications too.
Let $\Omega$ be an open pseudoconvex set in $\mathbb{C}^{\text {n }}$ such that there is an increasing sequence of open pseudoconvex subsets $\Omega_{k}$ with union $\Omega$ and with

$$
\forall \mathrm{k}, \exists \varepsilon=\varepsilon(\mathrm{k}): \quad \Omega_{\mathrm{k}}(\varepsilon) \subset \Omega_{\mathrm{k}+1}
$$

where $\Omega_{k}(\varepsilon)$ is the $\varepsilon$-neighborhood of $\Omega_{k}$. Moreover, in some theorems we require that there is a continuous plurisubharmonic function $\sigma$ in $\Omega$ with
(7.1) $\quad \Omega_{k}=\{z \mid \sigma(z)<k\}$.

This is only a special condition on $\Omega$ ([7], th.2.6.7.ii), if the sets $\Omega_{k}$ are unbounded.

For example, we may take for $\Omega_{k}$ suitable $\varepsilon$-neighborhoods of each other, since the function $d(z)=-\log \delta\left(z, \Omega^{c}\right)$ (here $\delta\left(z, \Omega^{c}\right)$ is the distance between $z \in \Omega$ and the complement of $\Omega$ ) is plurisubharmonic when $\Omega$ is pseudoconvex. We show that in some sense also the sets $\mathbb{R}^{2 n}+i C(k) \subset \mathbb{C}^{2 n}$ of the last section are an example. Therefore, we say that two increasing sequences $\left\{\Omega_{k}\right\}_{k=1}^{\infty}$ and $\left\{\Omega_{k}^{\prime}\right\}_{k=1}^{\infty}$ exhausting $\Omega$ are equivalent if for every $k$ there is an $\ell$ with $\Omega_{k} \subset \Omega_{\ell}^{\prime}$ and $\Omega_{k}^{\prime} \subset \Omega_{\ell}$. Then it is clear that any function on $\Omega$ that is bounded in some norm on all subsets $\Omega_{k}$ is also bounded on the subsets $\Omega_{k}^{\prime}$ and conversely.

LEMMA 7.1. The increasing sequence $\left\{\mathbb{R}^{2 n}+\mathrm{iC}(\mathrm{k})\right\}_{\mathrm{k}=1_{\infty}}^{\infty}$ exhausting $\Omega=\mathbb{R}^{2 \mathrm{n}}+\mathrm{iC} \subset$ $\subset \mathbb{C}^{2 n}$ is equivalent to an increasing sequence $\left\{\Omega_{k}\right\}_{k=1}^{\infty}$ satisfying (7.1). PROOF. Choose a vector $\alpha$ in $C \subset \mathbb{R}^{2 n}$ and a number $c>1$ and consider the hyperplane $H=\{y \mid \alpha \cdot y=c\} \subset \mathbb{R}^{2 n}$. Let for each $y \in C$

$$
\mathrm{y}^{*}=\frac{\mathrm{c}}{\alpha \cdot \mathrm{y}} \mathrm{y}
$$

be the intersection of the vector $y$ with $H$. We define a plurisubharmonic (even convex) function $x$ in $C$ by $x(y)=d\left(y^{*}\right)=-\log \delta\left(y^{*}, C^{c}\right)$. Then the sets $C_{k}^{\prime}=\{y \mid x(y)<k\}, k=1,2, \ldots$, are relatively compact subcones of $C$ exhausting $C$. Now we set for $z=x+i y$

$$
\sigma(z)=\max (d(y), x(y))
$$

which is plurisubharmonic (even convex) and we have

$$
\Omega_{k}=\{z \mid \sigma(z)<k\}=\left\{\mathbb{R}^{2 n}+i C_{k}^{\prime}\right\} \cap\left\{z \in \Omega \mid \delta\left(z, \Omega^{c}\right)>e^{-k}\right\}
$$

Hence the sequence $\left\{\Omega_{k}\right\}_{k=1}^{\infty}$ is obviously equivalent to $\left\{\mathbb{R}^{2 n}+i C(k)\right\}_{k=1}^{\infty}$.
Let $\phi$ be a plurisubharmonic function in $\Omega$. In some theorems $\phi$ will be such that for every $z \in \Omega_{k}$ and $\left\|z^{\prime}-z\right\|<\varepsilon(k)$

$$
\begin{equation*}
\phi\left(z^{\prime}\right)-\phi(z) \leq K_{k}, \tag{7.2}
\end{equation*}
$$

where the constant $K_{k}$ does not depend on $z$ and $z^{\prime}$, but may depend on $k$. For example the function $m \log \left(1+\|z\|^{2}\right)+2 a(y)+2 / \ell\|y\|$ is plurisubharmonic in $\mathbb{R}^{2 n}+i C$ and satisfies (7.2) for every sequence $\left\{\Omega_{k}\right\}_{k=1}^{\infty}$ equivalent to $\left\{\mathbb{R}^{2 n}+i C(k)\right\}_{k=1}^{\infty}$. Finally let $P=\left(P_{j k}\right), j=1, \ldots, p, k=1, \ldots, q$, be a matrix of polynomials.

Then problems 6.5 and 6.6 follow from lemma 7.1 and the next two theorems, theorem 7.1 (a) and theorem 7.2, respectively. In theorem 7.1 we formulate a part (b) with uniform bounds, which we do not need here, but which may be useful in other purposes. Part (b) is derived in the same way as part (a).

THEOREM 7.1. If $\mathrm{f} \in \mathrm{A}(\Omega)^{\mathrm{p}}$ can locally (that is in a neighborhood $w$ of any point in $\Omega$ ) be written as

$$
\mathrm{f}=\mathrm{Pg}_{\omega}, \quad \mathrm{g}_{\omega} \in \mathrm{A}(\omega)^{\mathrm{q}}
$$

then there is a number $N$, such that
(a) there is a function $\mathrm{v} \in \mathrm{A}(\Omega)^{\mathrm{q}}$ with $\mathrm{f}=\mathrm{Pv}$ and with

$$
\int_{\Omega_{k}}|v(z)|^{2} \frac{\exp ^{-\phi(z)}}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z)<\infty, \quad k=1,2, \ldots
$$

when $\Omega=\bigcup_{k=1}^{\infty} \Omega_{k}$ satisfies (7.1) and $\phi$ is a plurisubharmonic function in
$\Omega$ such that

$$
\int_{\Omega_{k}}|f(z)|^{2} \exp -\phi(z) d \lambda(z)<\infty, \quad k=1,2, \ldots .
$$

(b) For all $k=1,2, \ldots$ there are constants $K_{k}$, integers $l_{k} \geq k$ and functions $\mathrm{v}_{\mathrm{k}} \in \mathrm{A}\left(\Omega_{\mathrm{k}}\right)^{\mathrm{q}}$ with $\mathrm{f}=\mathrm{Pv}_{\mathrm{k}}$ in $\Omega_{\mathrm{k}}$ and with

$$
\int_{\Omega_{k}}\left|v_{k}(z)\right|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z) \leq K_{k} \int_{\Omega_{\ell}}|f(z)|^{2} \exp -\phi(z) d \lambda(z),
$$

when the right hand side is finite for some plurisubharmonic function $\phi$.
In part (b) the pseudoconvex subsets $\Omega_{k}$ of $\Omega$ do not have to satisfy (7.1). THEOREM 7.2. If $\mathrm{f}_{\mathrm{k}} \in \mathrm{A}\left(\Omega_{\mathrm{k}}\right)^{\mathrm{p}}, \mathrm{k}=1,2, \ldots$, are functions with $\mathrm{f}_{\mathrm{k}+1}-\mathrm{f}_{\mathrm{k}}=\mathrm{Pg}_{\mathrm{k}}$ in $\Omega_{k}, g_{k} \in A\left(\Omega_{k}\right)^{q}$, then there are a number $N$ and a function $f \in A(\Omega)^{p}$ with $\mathrm{f}-\mathrm{f}_{\mathrm{k}}=\mathrm{P} \tilde{\mathrm{g}}_{\mathrm{k}}$ in $\Omega_{\mathrm{k}}, \tilde{\mathrm{g}}_{\mathrm{k}} \in \mathrm{A}\left(\Omega_{\mathrm{k}}\right) \mathrm{q}$, and with

$$
\int_{\Omega_{k}}|f(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z)<\infty, \quad k=1,2, \ldots,
$$

when $\Omega=\bigcup_{k=1}^{\infty} \Omega_{k}$ satisfies (7.1) and $\phi$ is a plurisubharmonic function satisfying (7.2) such that

$$
\int_{\Omega_{k}}\left|f_{k}(z)\right|^{2} \exp -\phi(z) d \lambda(z)<\infty, \quad k=1,2, \ldots
$$

Here $|f(z)|^{2}$ means $\left|f_{1}(z)\right|^{2}+\ldots+\left|f_{p}(z)\right|^{2}$ when $f=\left(f_{1}, \ldots, f_{p}\right) \in A(\Omega)^{p}$ and $\lambda(z)$ denotes the Lebesgue measure in $\mathbb{C}^{\mathrm{n}}$.

First we need similar theorems as theorem A13 and Castan's theorem B, theorem A 14, but now with estimates. Let $U^{(\lambda)}=\left\{U_{i}^{(\lambda)}\right\}_{i \in I_{\lambda}}, \lambda=0,1,2, \ldots$ be the coverings of $\Omega$ given in the appendix section $V$ satisfying properties (A15)(i),(ii),(iii), (iv), (v) and (vi) and let for every $k u_{k}^{(\lambda)}=$ $=\left\{U_{i}^{(\lambda)} \cap \Omega_{k}\right\}_{i \in I_{\lambda}}$ be the corresponding coverings of $\Omega_{k}$. When $F$ is an analytic sheaf on $\Omega$, we denote by $C^{\mathrm{p}}\left[U^{(\lambda)}, F, \phi\right]$ the set of alternating cochains $c=\left\{c_{s}\right\}$ in $\Omega, s \in I_{\lambda}^{p+1}, c_{s} \in \Gamma\left(U_{s}^{(\lambda)}, F\right)$, satisfying for all $k$

$$
\|c\|_{\phi, k}^{2} \stackrel{\text { def }}{=} \sum_{|s|=p+1} \int_{U_{S}^{(\lambda)}{ }_{n \Omega_{k}}}\left|c_{s}(z)\right|^{2} \exp -\phi(z) d \lambda(z)<\infty
$$

and by $C^{p}\left(U_{k}^{(\lambda)}, F, \phi\right)$ the set of all alternating cochains $c$ in $\Omega_{k}$ with $\|c\|_{\phi, k}^{2}<\infty$. By $\phi_{N}$ we will mean the plurisubharmonic function $\phi_{N}(z)=\phi(z)+$ $+N \log \left(1+\|z\|^{2}\right)$.

Lemma 7.2 will be obtained in the same way as theorem A 13, only we write down explicitely the construction of the map $\delta^{*}$ (A10), so that we can bring estimates in the statements involving $\delta^{*}$. We do not work with the sheaf $E$ of germs of $C^{\infty}$-functions, but rather with the sheaf $L$ of germs of locally square integrable functions. Then we may use theorem Al0 instead of theorem A 9. So let $L_{q}$ be the sheaf of germs of ( $0, q$ )-forms with locally square integrable coefficients and 1 et $Z_{q}$ be the subsheaf of $\bar{\partial}-c l o s e d$ forms of type ( $0, q$ ). Again we have a part (a) with globally defined functions on $\Omega$ and a part (b) with functions in $\Omega_{k}$ and uniform bounds.

LEMMA 7.2.
(a) Every cochain $c$ in $C^{p}\left[U^{(\lambda)}, A, \phi\right], p \geq 1$, with $\delta c=0$ can be written as $c=\delta c^{\prime}$ for $a c^{\prime} \in C^{p^{-1}}\left[U^{(\lambda)}, A, \phi_{2 m}\right]$, where $m=\min (p, n)$, when $\left\{\Omega_{k}\right\}_{k=1}^{\infty}$ satisfies (7.1).
(b) For all $k$ every cochain $c$ in $C^{P^{P}}\left[U^{(\lambda)}, A, \phi\right], p \geq 1$, with $\delta c=0$ can be written in $\Omega_{k}$ as $c=\delta c_{k}^{\prime}$ for a $c_{k}^{\prime} \in C^{p-1}\left(U_{k}^{(\lambda)}, A, \phi_{2 m}\right)$ such that for some constants $K_{k}$

$$
\left\|c_{k}^{\prime}\right\|_{\phi_{2 m}}, k \leq K_{k}\|c\|_{\phi, k},
$$

where $m=\min (p, n)$. Also for fixed kevery cocycze $c \in C^{p}\left(U_{k}^{(\lambda)}, A, \phi\right)$ satisfies the above property (b) for this $k$.

PROOF. A section $c \in \Gamma\left(\Omega, L_{0}\right)$ with $\bar{\partial} c=0$ determines a holomorphic function $c \in A(\Omega)$ (this follows by repeated use of lemma 4.2 .4 in [7]). For $c \in$ $\in C^{p}\left[U^{(\lambda)}, Z_{q}, \phi\right]$ with $\delta c=0$ we want to find a $c^{\prime} \in C^{p-1}\left[U^{(\lambda)}, Z_{q}, \phi{ }_{2 m}\right]$ such that $\delta c^{\prime}=c$, when $q=0$ or in part (b) cochains $c_{k}^{\prime} \in C^{p-1}\left(U_{k}^{(\lambda)}, Z_{q}, \phi_{2 m}\right)$ such that $\delta c_{k}^{\prime}=c$ in $\Omega_{k}$. We assume that this has already been proved for smaller values of $p$ and all $q$, when $p>1, m=p$ and when the constants $K_{k}$ in part (b) depend on $p$.

First we give estimates in the construction of $g$ in theorem A12. For each $k$ we choose $\ell=\ell(k)$ such that, when $U_{i}^{(\lambda)} \cap \Omega_{k} \neq \emptyset, U_{i}^{(\lambda)} \subset \Omega_{\ell}$ according
to property (A15)(iii). Since also all sets in $U^{(\lambda)}$ contained in $\Omega_{k}$ have a minimum size (say they contain a ball with radius $\varepsilon_{k}(\lambda)$ ), we can construct a partition $\left\{\phi_{\nu}\right\}_{\nu=1}^{\infty}$ of unity subordinate to the covering $U^{(\lambda)}$ of $\Omega$ ( $\phi_{\nu}$ has its support in $\mathrm{U}_{\mathrm{i}_{v}}^{(\lambda)}$ ), such that for all k

$$
\begin{equation*}
\max _{z}\left|\bar{\partial}_{\nu}(z)\right|^{2} \leq C_{k} \tag{7.3}
\end{equation*}
$$

for those $v$ with $U_{i_{v}}^{(\lambda)} \cap \Omega_{k} \neq \emptyset$. Here

$$
|\bar{\partial} \phi(z)|^{2}=\sum_{j=1}^{n}\left|\bar{\partial} \bar{j}_{j} \phi(z)\right|^{2} .
$$

For example let for each $v \in I_{\lambda+1} X_{\nu}$ be a $C^{\infty}$-function equal to one in $U_{\nu}^{(\lambda+1)}$ and to zero outside the $\varepsilon_{\ell\left(k_{\nu}\right)}(\lambda+1)$-neighborhood of $U_{\nu}^{(\lambda+1)}$ (which

figure 7.1.
certainly is contained in $U_{\rho}^{(\lambda)}$ ( $\nu$ with $\rho=\rho_{\lambda, \lambda+1}$, because of property (A15) (v)), where $k_{\nu}$ is the smallest integer $k$ with $\left.U_{\rho}(\lambda){ }^{\rho}\right) \Omega_{k} \neq \emptyset$, see figure 7.1. Then for those $v$ with $U_{\rho}^{(\lambda)} \cap \Omega_{k} \neq \emptyset \max _{z}\left|\frac{\rho}{\partial} x_{\nu}(z)\right|$ depends on $\varepsilon_{\ell(1)}(\lambda+1), \ldots, \varepsilon_{\ell(k)}^{(\lambda+1)}$. Since $u^{(\lambda+1)}$ is a covering of $\Omega$,

$$
\phi_{\nu}(z)=\frac{\chi_{\nu}(z)}{\sum_{\mu=1}^{\infty} \chi_{\mu}(z)}, \quad v=1,2, \ldots
$$

is a partition of unity subordinate to the covering $U^{(\lambda)}$ with $\mathbf{i}_{\nu}=\rho(\nu)$, that satisfies (7.3). Note that for each $z$ not more than $M$ terms in the denominator differ from zero because of property (A15)(ii).

As in the proof of theorem Al2 we set for $s \in I_{\lambda}^{p}$

$$
g_{s}=\sum_{v} \phi_{\nu} c_{i_{V}},
$$

when $c \in C^{p}\left[U^{(\lambda)}, Z_{q}, \phi\right]$. Then as in theorem A12 $g \in C^{p-1}\left(U^{(\lambda)}, L_{q}\right)$ and $\delta g=c$, if $\delta c=0$. Furthermore writing $\phi_{\nu}=\sqrt{\phi_{\nu}} \cdot \sqrt{\phi_{\nu}}$ and using $\sum_{\nu} \phi_{\nu}=1$ we find

$$
\begin{aligned}
& \int_{U_{S}^{(\lambda)}}{ }_{n \Omega_{k}}\left|g_{S}(z)\right|^{2} \exp -\psi(z) d \lambda(z) \stackrel{\text { not }}{=}\left\|g_{S}\right\|_{\psi, k}^{2} \leq \\
& \quad \leq \sum_{V} \int_{U} \int_{S}^{(\lambda)}{ }_{n \Omega_{k}} \phi_{V}(z)\left|c_{i_{V}}(z)\right|^{2} \exp -\psi(z) d \lambda(z) \leq \\
& \quad \leq \sum_{V}\left\|c_{i_{V}}\right\|_{\psi, k}^{2}
\end{aligned}
$$

for all plurisubharmonic functions $\psi$ for which the right hand side is finite. Since not more than $M_{\lambda, \lambda+1}(k)$ different $v^{\prime} s$ are mapped by $\rho$ onto the same $i$, when $U_{i}^{(\lambda)} \cap \Omega_{k} \neq \emptyset$, (property (A15) (vi), we get by summing up

$$
\|g\|_{\psi, k}^{2} \leq M_{\lambda, \lambda+1}(k)\|c\|_{\psi, k}^{2}
$$

Let $\bar{\partial} g=f$ be the cochain in $C^{p-1}\left(U^{(\lambda)}, Z_{q+1}\right)$ defined by

$$
\mathrm{f}_{\mathrm{s}}=\bar{\partial} \mathrm{g}_{\mathrm{s}}=\sum_{\nu} \partial \phi_{\nu} \wedge c_{i_{\nu}}
$$

Then

$$
\left\|f_{s}\right\|_{\phi, k}^{2} \leq\left\{\sum_{\nu}^{\|} \partial \phi_{\nu} \wedge c_{i_{\nu}} s_{\phi, k}\right\}^{2} \leq N_{k}^{(\lambda)} \sum_{\nu} \| \bar{\partial}_{\partial}{ }_{\nu} \wedge c_{i_{\nu}} s^{\|}{ }_{\phi, k}^{2}
$$

where at most $N_{k}^{(\lambda)}$ terms in the sum are different from zero, when $U^{(\lambda)} \cap$ $\cap \Omega_{k} \neq \emptyset$ according to property (A15)(iv). If $U_{s}^{(\lambda)} \cap \Omega_{k} \neq \emptyset$, then $U_{s}^{(\lambda)} \subset \Omega_{\ell(k)}$, so that for all $i \in I_{\lambda}$ with $U_{i}^{(\lambda)} \cap U_{S}^{(\lambda)} \neq \emptyset$

$$
\mathrm{U}_{\mathrm{i}}^{(\lambda)} \cap \Omega_{\ell(k)} \neq \emptyset .
$$

Hence using (7.3) in the above estimate we get with $K_{k}^{\prime}=C_{\ell(k)} N_{k}^{(\lambda)} M_{\lambda, \lambda+1}(k)$

$$
\|\mathrm{f}\|_{\phi, \mathrm{k}}^{2} \leq \mathrm{K}_{\mathrm{k}}^{\prime}\|\mathrm{c}\|_{\phi, \mathrm{k}}^{2}
$$

Now $\delta f=\bar{\partial} \delta g=\bar{\partial} c=0$. If $p>1$, by the inductive hypothesis of case (a) we can find a cochain $f^{\prime} \in C^{p-2}\left[U^{(\lambda)}, Z_{q+1}, \phi_{2 p-2}\right]$ with $\delta f^{\prime}=f$ and by the inductive hypothesis of case (b) we can find cochains $f_{k}^{\prime} \epsilon$ $\in C^{p^{-2}}\left(U_{k}^{(\lambda)}, Z_{q+1}, \phi_{2 p-2}\right)$ with $\delta f_{k}^{\prime}=f$ in $\Omega_{k}$ and with

$$
\left\|f_{k}^{\prime}\right\|_{\phi}^{2}, K_{k-2}^{\prime \prime \| f} \|_{\phi, k}^{2}
$$

for some constants $K_{k}^{\prime \prime}$ depending on $k$. By theorem A10 and property (A15) (i) for every $s \in I_{\lambda}^{p-1}$ we can find $\left(g^{\prime}\right)_{s} \in \Gamma\left(U_{s}^{(\lambda)}, L_{q}\right)$ so that $\bar{\partial}\left(g^{\prime}\right)_{s}=\left(f^{\prime}\right)_{s}$ in $\mathrm{U}_{\mathrm{S}}^{(\lambda)}$ and

$$
\left\|\left(g^{\prime}\right)_{s}^{\|_{\phi_{2 p}}^{2}} \leq\right\|\left(f^{\prime}\right)_{S}^{\|_{\phi_{2 p-2}}^{2}}
$$

and since the sets $\Omega_{k}$ are pseudoconvex, theorem A10 yields $\left(g_{k}^{\prime}\right)_{s}$ $\in \Gamma\left(U_{s}^{(\lambda)} \cap \Omega_{k}, L_{q}\right)$, such that $\bar{\partial}\left(g_{k}^{\prime}\right)_{s}=\left(f_{k}^{\prime}\right)_{s}$ in $U_{s}^{(\lambda)} \cap \Omega_{k}$ and

$$
\left\|\left(g_{k}^{\prime}\right)_{s}\right\|_{\phi_{2 p}}^{2} \leq\left\|\left(f_{k}^{\prime}\right)_{s}\right\|_{\phi_{2 p-2}}^{2}
$$

Hence $\left\{\left(g^{\prime}\right) s \mid s \in I_{\lambda}^{p-1}\right\}=g^{\prime} \in C^{p-2}\left[U^{(\lambda)}, L_{q}, \phi_{2 p}\right]$ and $\left\{\left(g_{k}^{\prime}\right)_{s} \mid s \in I_{\lambda}^{p-1}\right\}=$ $=g_{k}^{\prime} \in C^{p^{-2}}\left(U_{k}^{(\lambda)}, L_{q}, \phi_{2 p}\right)$.

Finally put $c^{\prime}=g-\delta g^{\prime}$ and $c_{k}^{\prime}=g-\delta g_{k}^{\prime}$, then for all $k$

$$
\begin{aligned}
\left\|c^{\prime}\right\|_{\phi_{2 p}, k}^{2} & \leq\|g\|_{\phi_{2 p}, k}^{2}+\mathrm{pN}_{k}^{(\lambda)}\left\|g^{\prime}\right\|_{\phi_{2 p}, k}^{2} \leq \\
& \leq M_{\lambda, \lambda+1}(k)\|c\|_{\phi_{2 p}, k}^{2}+{ }_{p N_{k}}^{(\lambda)}{\left\|f f^{\prime}\right\|_{\phi_{2 p-2}}^{2}, k}_{2}<\infty
\end{aligned}
$$

and for some constants $\mathrm{K}_{\mathrm{k}}$

$$
\left\|c_{k}^{\prime \prime \|} \phi_{2 p}^{2}, k \leq M_{\lambda, \lambda+1}(k)\right\| c\left\|_{\phi_{2 p}, k}^{2}+\mathrm{pN}_{k}^{(\lambda)_{\| f}^{\prime} \|_{\phi}^{2 p-2}} 2, k \leq K_{k}\right\| c \|_{\phi, k}^{2} .
$$

Furthermore $\delta c^{\prime}=\delta g=c$ and $\bar{\partial} c^{\prime}=f-\delta \bar{\partial} g^{\prime}=f-\delta f^{\prime}=f-f=0$, hence $c^{\prime} \in C^{p-1}\left[u^{(\lambda)}, Z_{q}, \phi_{2 p}\right]$ and also $\delta c_{k}^{\prime}=\delta g=c$ and $\bar{\partial} c_{k}^{\prime}=f-\delta \bar{\partial} g_{k}^{\prime}=f-\delta f_{k}^{\prime}=$ $=f-f=0$ in $\Omega_{k}$, hence $c_{k}^{\prime} \in C^{p-1}\left(u_{k}^{(\lambda)}, Z_{q}, \phi_{2 p}\right)$.

It remains to consider the case $p=1$. The fact that $\delta f=0$ then means that $f$ defines uniquely a ( $0, q+1$ )-form $f$ in $\Omega$ with $\bar{\partial} f=0$. In case (a) we cannot immediately apply theorem AlO, but we need a modification, where the integrals are performed in the sets $\Omega_{k}$. Assume that this may be done. Then we can find $\tilde{g} \in \Gamma\left(\Omega, L_{q}\right)$ with $\bar{\partial} \tilde{g}=f$ and for all $k$

$$
\int_{\Omega \mathrm{k}}|\tilde{g}(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{2}} \mathrm{~d} \lambda(z)<\infty .
$$

In case (b) we use theorem A10 and obtain ( $0, q$ )-forms $\tilde{g}_{k} \in \Gamma\left(\Omega_{k}, L_{q}\right)$ in $\Omega_{k}$ with $\bar{\partial} \tilde{g}_{k}=\left.f\right|_{\Omega_{k}}$ and

$$
\int_{\Omega_{k}}\left|\tilde{\mathrm{~g}}_{\mathrm{k}}(z)\right|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{2}} \mathrm{~d} \lambda(z) \leq \int_{\Omega_{k}}|f(z)|^{2} \exp -\phi(z) \mathrm{d} \lambda(z) .
$$

Putting

$$
\left(c^{\prime}\right)_{s}=g_{s}-\left.\tilde{g}\right|_{U_{s}} ^{(\lambda)} \quad \text { and } \quad\left(c_{k}^{\prime}\right)_{s}=g_{s}-\left.\tilde{g}_{k}\right|_{U}(\lambda)_{n \Omega_{k}}
$$

for $s \in I_{\lambda}$, we obtain cochains with the required properties (using (Al5)(ii) in the estimates for $\tilde{g}$ ).

In fact we only have $n$ induction steps, since all ( $0, n$ )-forms $g$ satisfy $\bar{\partial} g=0$. Therefore, the estimates hold already when $p$ is replaced by $\min (\mathrm{p}, \mathrm{n})$ and the constants $\mathrm{K}_{\mathrm{k}}$ in part (b) may be taken independent of p .

We only have to prove the modification of theorem AlO.

LEMMA 7.3. Let $\Omega$ be an open pseudoconvex set, let $\left\{\Omega_{k}\right\}_{k=1}^{\infty}$ be an increasing sequence of subsets of $\Omega$ satisfying (7.1) and let $\phi$ be a plurisubharmonic function in $\Omega$. For every ( $0, q+1$ )-form $g$ with locally square integrable coefficients and with $\bar{\partial} \mathrm{g}=0$, there is a ( $0, \mathrm{q}$ )-form u in $\Omega$ with locally square integrable coefficients, such that $\bar{\partial} \mathrm{u}=\mathrm{g}$ and for all k

$$
\int_{\Omega_{k}}|u(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{2}} d \lambda(z)<\infty
$$

provided that for each $k$

$$
\int_{\Omega_{k}}|g(z)|^{2} \exp -\phi(z) d \lambda(z)<\infty .
$$

PROOF. Let $x$ be a convex majorant of the nonnegative function $\tilde{\chi}$

$$
\tilde{x}(a)= \begin{cases}0 & \text { for } a<1, \\ \max \left\{0, \log \left[2^{k} \int_{\Omega_{k+1} \backslash \Omega_{k}}|g(z)|^{2} \exp -\phi(z) d \lambda(z)\right]\right\} \text { for } k \leq a<k+1,\end{cases}
$$

$k=1,2, \ldots$. Then $\psi(z) \xlongequal{\text { def }} \chi(\sigma(z))$ is plurisubharmonic in $\Omega$, so that we may apply theorem Al0 in the domain $\Omega$ with the plurisubharmonic function $\phi+\psi$. This yields a ( $0, q$ )-form $u$ in $\Omega$ with $\bar{\partial} u=g$ and with for each $k$

$$
\begin{aligned}
& \int_{\Omega_{k}}|u(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{2}} d \lambda(z) \leq e^{\chi(k)} \int_{\Omega_{k}}|u(z)|^{2} \frac{\exp (-\phi(z)-\psi(z))}{\left(1+\|z\|^{2}\right)^{2}} d \lambda(z) \leq \\
& \leq e^{\chi(k)} \int_{\Omega}|u(z)|^{2} \frac{\exp (-\phi(z)-\psi(z))}{\left(1+\|z\|^{2}\right)^{2}} d \lambda(z) \leq e^{\chi(k)} \int_{\Omega}|g(z)|^{2} \exp (-\phi(z)-\psi(z)) d \lambda(z)= \\
& =e^{\chi(k)}\left(\int_{\Omega_{1}}+\sum_{k=1}^{\infty} \int_{\Omega_{k+1} \backslash \Omega_{k}}\right)|g(z)|^{2} \exp (-\phi(z)-\psi(z)) d \lambda(z) \leq \\
& \leq e^{\chi(k)}\left[\int_{\Omega_{1}}|g(z)|^{2}{ }^{\left.\exp -\phi(z) d \lambda(z)+\sum_{k=1}^{\infty} \frac{1}{2^{k}}\right]<\infty .}\right.
\end{aligned}
$$

REMARK 7.1. In general lemma 7.3 is not true, if we consider different weightfunctions $\phi_{k}$ in the sets $\Omega_{k}$, or in the same set $\Omega$. For example assume
that $\bar{\partial} \mathrm{g}=0$ and that for every k
(7.4)

$$
\int_{\Omega_{k}}|g(z)|^{2} \exp (-\phi(z)-1 / k\|z\|) d \lambda(z)<\infty,
$$

where $\Omega_{k} \subset \Omega_{k+1} \subset \Omega_{\text {or }} \Omega_{k}=\Omega$ for all $k$. Then it is not true that there is a form $u$ in $\Omega$ with $\bar{\partial} u=g$ and with for all $k$

$$
\begin{equation*}
\int_{\Omega_{k}}|u(z)|^{2} \frac{\exp (-\phi(z)-1 / k\|z\|)}{\left(1+\|z\|^{2}\right)^{2}} \mathrm{~d} \lambda(z)<\infty . \tag{7.5}
\end{equation*}
$$

For if this were true, using theorem 4.4 .2 of [7] as in section 5, we could extend the entire function

$$
f(z)=\oint e^{1 / \zeta} e^{i z \cdot \zeta} d \zeta
$$

in $\mathbb{C}^{1}$ satisfying

$$
|f(z)| \leq 2 \varepsilon e^{1 / \varepsilon} e^{\varepsilon\|z\|} \quad \text { for all } \varepsilon>0
$$

to an entire function $F$ in $\mathbb{C}^{2}$ satisfying

$$
\begin{aligned}
& F(z, i z)=f(z) \\
& \left|F\left(\nu_{1}, \nu_{2}\right)\right| \leq K_{\varepsilon}(1+\|v\|)^{m} e^{\varepsilon\|\operatorname{Im}\| \|} \quad \text { for all } \varepsilon>0 .
\end{aligned}
$$

But then according to VLADIMIROV [16] 29.1 F is a polynomial, hence f would be a polynomial, that is

$$
f(z)=\sum_{j=0}^{k} \oint \frac{a^{j}}{\zeta^{j}} e^{i z \cdot \zeta} d \zeta
$$

for some $k$ and constants $a_{j}$ contradicting the definition of $f$.
In [9] KAWAI has shown that for each ( $0, q+1$ )-form $g$ with $\bar{\partial} g=0$
satisfying (7.4) there does exist a ( $0, q$ )-form $u$ with $\bar{\partial} u=g$ in $\Omega$ satisfying (7.5), when $\Omega$ satisfies

```
sup|Imz|
z\in\Omega
```

for some constant K .
Next we derive Cartan's theorem B with bounds. Let $F$ be either the sheaf of relations of $P$ on $\Omega$, thus $F=R_{P}$ or the image under $P$ of the sheaf $A^{q}$, thus $F=P A^{q}$, see (A5) and (A6).

LEMMA 7.4. Let the plurisubharmonic function $\phi$ in the pseudoconvex open set $\Omega$ satisfy (7.2). There is a positive integer $N$ (depending on $P$ ), such that for all $\lambda$ there is a $\mu>\lambda$ with the following properties
(a) when moreover the subsets $\Omega_{k}$ of $\Omega$ satisfy (7.1), every cochain $\mathrm{f} \in \mathrm{C}^{\mathrm{P}}\left[\mathrm{U}^{(\lambda)}, \mathrm{F}, \phi\right]$ with $\delta \mathrm{f}=0, \mathrm{p} \geq 1$, can be written as $\delta \mathrm{f}^{\prime}=\rho_{\lambda, \mu^{*}}^{\mathrm{f}}$ for some $\mathrm{f}^{\prime} \in \mathrm{C}^{\mathrm{p}-1}\left[U^{(\mu)}, F, \phi_{\mathrm{N}}\right]$;
(b) for all $k$ there are integers $l_{k}>k$ and constants $K_{\lambda, k}$, such that every cochain $\mathrm{f} \in \mathrm{C}^{\mathrm{p}}\left[U^{(\lambda)}, F, \phi\right]$ with $\delta \mathrm{f}=0, \mathrm{p} \geq 1$, for all k can be written as $\delta f_{k}^{\prime}=\rho_{\lambda, \mu}^{*} \mathrm{f}$ in $\Omega_{k}$ with $f_{k}^{\prime} \in C^{p-1}\left(u_{k}^{(\mu)}, F, \phi_{N}\right)$ and with

$$
\left\|f_{k}^{\prime}\right\|_{\phi_{N}, k} \leq K_{\lambda, k}\|f\|_{\phi, l_{k}}
$$

PROOF. First we change theorem A16 into a formulation with $\mathrm{L}^{2}$-estimates. Let $K=\omega+z$ be so that $U \cap \Omega_{k} \neq \emptyset$ and $V=(t+1) \omega+z \subset \Omega_{\ell}$ for some $\ell \geq k$, where $t \omega$ is the enlargement of $\omega$ by a factor $t$ with respect to some center in $\omega$. Then $V$ contains some $\varepsilon$-neighborhood of $t \omega+z$, where $\varepsilon$ depends on the size of w. The condition $\mathrm{HS}_{1}$ ([14] G.7 with $M_{p}^{2}={\exp -\phi_{N+m}, \Omega_{p}=\Omega_{m}=U \text { and } M_{m}^{2}=}^{m}$ $\left.=\exp -\phi_{N+m+(n+1) / 2}\right)$ and by (7.2) the condition $\mathrm{HS}_{2}\left([14] \mathrm{G} .7\right.$ with $M_{p}^{2}=$ $=\exp -\phi_{m}, \Omega_{p}=V$ and $M_{m}^{2}=\exp -\phi_{m}, \Omega_{m}=t \omega+z$ and with $d_{z}=\varepsilon$ ) are satisfied. Hence instead of (A14) we get

$$
\begin{align*}
& \int_{U}|v(w)|^{2} \frac{\exp -\phi(w)}{\left(1+\|w\|^{2}\right)^{N+m+(n+1) / 2} d \lambda(w) \leq C_{1} \sup _{w \in U}|v(w)| \frac{\exp -\frac{1}{2} \phi(w)}{(1+\|w\|)^{N+m}} \leq} \\
& \leq C_{1} C \sup _{w \in t \omega+z}|P(w) u(w)| \frac{\exp -\frac{1}{2} \phi(w)}{(1+\| w)^{m}}\left(\frac{1+\|z\|}{1+\|w\|}\right)^{N} \leq  \tag{7.6}\\
& \leq C_{1} C\left(1+\sup _{\zeta \in t \omega}\| \|\right)^{N} \sup _{w \in t \omega+z}|P(w) u(w)| \frac{\exp -\frac{1}{2} \phi(w)}{(1+\|w\|)^{m}} \leq \\
& \leq C_{2} \int_{V}|P(w) u(w)| \frac{\exp -\phi(w)}{\left(1+\|w\|^{2}\right)^{m}} d \lambda(w),
\end{align*}
$$

since

$$
-\phi_{N+m}(w)+\phi_{N+m}(v) \leq \text { some constant for } w \in U \text { and } v \in t \omega+z
$$

follows by repeated use of (7.2) and since the estimates with $\left(1+\left\|_{W}\right\|^{2}\right)^{m}$ and $(1+\|w\|)^{2 m}$ are equivalent and $\left|v_{1}(w)\right|+\ldots+\left|v_{q}(w)\right| \leq \sqrt{q}|v(w)|$ when $v=\left(v_{1}, \ldots, v_{q}\right) \in A(U)^{q}$. The constants $C_{1}, C$ and $C_{2}$ do not depend on $z$, $C$ depends on the size of $\omega$ and $C_{2}$ depends moreover on $\varepsilon$, in fact $C_{2} \sim \varepsilon^{-n}$ (see [18], proof of $\mathrm{HS}_{2}$ ), but $\varepsilon$ depends on the size of $\omega_{\text {. }}$

For $p \geq M$ (see (A15)(ii)) the theorem is true, since there are no non-zero cochains $f \in C^{M}\left[U^{(\lambda)}, F, \phi\right]$. Thus assume that the theorem has been proved for all $P$ when $p$ is replaced by $p+1$ and when the constants $N$ and $\mu$ and in part (b) the constants $\ell$ and $K$ depend on $p$.

In case $F=R_{P}$ there is a polynomial matrix $Q$, such that $F=Q A^{r}$ by (A6) and according to theorem All we can write $f \in C^{P}\left[U^{(\mu)}, F, \phi\right]$ as $f{ }_{s}=Q g_{s}$ where $g \in C^{p}\left(U^{(\mu)}, A^{r}\right)$. In case $F=P A^{q}$ we write $Q=P$ and $r=q$, then also $f=Q g$ with $g \in C^{P}\left(U^{(\mu)}, A^{r}\right)$ according to theorem A15. Let $v \geq \mu+{ }^{2} \log (t+1)$, then for every $i \in I_{\nu}(t+1)$ times $U_{i}^{(\nu)}$ is contained in $U_{\rho_{\mu, \nu}}^{(\mu)}(i)$, where $t$ is such that (7.6) may be applied with $U=U_{i}^{(\nu)}$ and $V=U_{\rho_{\mu, \nu}}^{(\mu)}(i)$. From theorem A16 and (7.6) we obtain a cochain $\tilde{g} \in C^{p}\left(U^{(\nu)}, A^{r}\right)$ with $Q \tilde{g}_{s}=\mathcal{V g}_{s},=f_{s}$, where $s^{\prime}=\rho_{\mu, \nu}(s)$, hence $\rho_{\mu, \nu}^{*} f=Q \tilde{g}$. When $U_{S}^{(\nu)} \cap \Omega_{k}=\emptyset$, then $U_{s}\left(\mu_{i}\right) \subset \Omega_{\ell(k)}$ for some $\ell(k)$ (property (A15)(iii)), so that (7.6) yields

$$
\int_{U_{s}^{(\nu)}}^{n \Omega_{k}}\left|\tilde{g}_{s}\right|^{2} \exp ^{\left(\phi_{N_{1}}+m\right.} d \lambda \leq C_{k, \mu} \int_{U^{\prime}} \int_{(\mu)}\left|f_{s^{\prime}}\right|^{2} \exp -\phi_{m} d \lambda
$$

for some $N_{1}$ and all $m$. The constant $C_{k, \mu}$ depends on the smallest and the largest size of the sets $U_{S}^{(\nu)}$ with $U_{S}^{(\nu)} \cap \Omega_{k} \neq \emptyset$ and this depends on $k$ and $\nu$, but $\nu$ depends on $\mu ; C_{k, \mu}$ does not depend on $s$. Since not more than a finite number of different $s$ are mapped by $\rho_{\mu, \nu}$ onto the same $s^{\prime}$ (property (A15)(vi)), we get by summing up

$$
\begin{equation*}
\|\tilde{g}\|_{\phi_{N_{1}+m}}, k \leq C_{k, \mu}^{\prime}\|f\|_{\phi_{m}}, \ell(k), \quad k=1,2, \ldots . \tag{7.7}
\end{equation*}
$$

Thus $\tilde{g} \in C^{P}\left[U^{(\nu)}, A^{r}, \phi_{N_{1}+m}\right]$. When $\delta f=0, \delta Q \tilde{g}=Q \delta \tilde{g}=0$, whence together with (A15) (iv) it follows that $\delta \tilde{g}=c$ is a cocycle in $C^{p+1}\left[U^{(\nu)}, R_{Q}, \phi_{N_{1}+m}\right]$.

By the inductive hypothesis of case (a) we can find $\mu^{\prime}>\nu, N_{2}$ and a cochain $c^{\prime} \in C^{P^{2}}\left[U^{\left(\mu^{\prime}\right)}, R_{Q}, \phi_{N_{2}+N_{1}+m}\right]$ with $\delta c^{\prime}=\rho_{\nu, \mu}^{*}, c$ in $\Omega$ and by the inductive hypothesis of case (b) we can find moreover constants $\ell_{k}>k, K_{v, k}^{\prime \prime}$ and cochains $c_{k}^{\prime} \in C^{p}\left(U_{k}^{\left(\mu^{\prime}\right)}, R_{Q}, \phi_{N_{2}+N_{1}+m}\right)$ with $\delta c_{k}^{\prime}=\rho_{\nu, \mu}^{*}, c$ in $\Omega_{k}$ and with

$$
\left\|c_{k}^{\prime}\right\| \phi_{N_{2}+N_{1}+m}, k \leq K_{v, k}^{\prime \prime}\|c\|_{\phi_{N_{1}+m}}, l_{k}
$$

We put $g_{0}=\rho_{\nu, \mu^{\prime}}^{*} \tilde{g}-c^{\prime} \in C^{P^{p}}\left[U^{\left(\mu^{\prime}\right)}, A^{r}, \phi_{N_{2}+N_{1}+m}\right]$, so that $\delta g_{0}=$ $=\rho_{\nu, \mu^{\prime}}^{*} c-\rho_{\nu, \mu^{\prime}}^{*} c=0$, and $g_{k}=\rho_{\nu, \mu}^{*}, \tilde{g}-c_{k}^{\prime} \in C^{P}\left(U_{k}^{\left(\mu^{\prime}\right)}, A^{r}, \phi_{N_{2}+N_{1}+m}\right)$ so that $\delta g_{k}=\rho_{\nu, \mu}^{*}, c=0$ in $\Omega_{k}$. According to lemma 7.2 (a) and (b) there are $g^{\prime} \in C^{P^{-1}}\left[U^{\left(\mu^{\prime}\right)}, A^{r}, \phi_{N}\right]$ with $\delta g^{\prime}=g_{0}$ and $g_{k}^{\prime} \in C^{p-1}\left(U_{k}^{\left(\mu^{\prime}\right)}, A^{r}, \phi_{N}\right)$ with $\delta g_{k}^{\prime}=$ $=g_{k}$ in $\Omega_{k}$, respectively, where $N=N_{2}+N_{1}+m+2 \min (n, p)$ and with moreover

$$
\left\|g_{k}^{\prime \|} \phi_{N}, k \leq K_{k}\right\| g_{k} \| \phi_{N_{2}+N_{1}+m}, k^{\bullet}
$$

Finally we put $f^{\prime}=Q g^{\prime} \in C^{p-1}\left[U^{\left(\mu^{\prime}\right)}, F, \phi_{N_{3}+N}\right]$ and $f_{k}^{\prime}=Q g_{k}^{\prime} \epsilon$ $\in C^{P-1}\left(U_{k}^{\left(\mu^{\prime}\right)}, F, \phi_{N_{3}+N}\right)$ with $N_{3}$ depending on $Q$. Then

$$
\delta f^{\prime}=Q \delta g^{\prime}=Q g_{0}=\rho_{\nu, \mu}^{*}, Q \tilde{g}=\rho_{\nu, \mu^{\prime}}^{*} \rho_{\mu, \nu}^{*} f=\rho_{\mu, \mu}^{*}, f
$$

in $\Omega$ and similarly $\delta f_{k}^{\prime}=\rho_{\mu, \mu}^{*}, f$ in $\Omega_{k}$. Furthermore for all $m$ and $\mu$ we get

$$
\begin{aligned}
& \left\|f_{k}^{\prime}\right\| \phi_{N_{3}+N}, k \leq K_{k}^{\prime} \| g_{k}^{\|} \phi_{N_{2}+N_{1}+m}, k \leq K_{k}^{\prime}\left\{M_{v, \mu^{\prime}}(k)^{p+1}\|\tilde{g}\|_{\phi_{N_{2}}+N_{1}+m}, k\right. \\
& \left.+K_{v, k}^{\prime \prime}\|c\| \phi_{N_{1}+m}, \ell_{k}\right\} \leq \\
& \leq K_{k}^{\prime}\left\{M_{v, \mu},(k)^{p+1}+K_{v, k}^{\prime \prime}(p+2) N_{l_{k}}^{(v)}\right\}\|\tilde{g}\|_{\phi_{N_{1}+m}, l_{k}} \leq K_{\mu, k}\|f\|_{\phi_{m}, \ell},
\end{aligned}
$$

where $\ell=\ell\left(\ell_{k}\right)$ depends on $\ell_{k}$ according to (7.7) and where $K_{\mu, k}$ is a constant depending on $k, \nu$ and $\mu^{\prime}$, but $\mu^{\prime}$ depends on $\nu, \ell_{k}$ depends on $k$ and $\nu$ depends on $\mu ; N_{3}$ depends on $Q, N_{2}$ on $p$ by the inductive hypothesis and $N_{1}$ on $P$, but
$Q$ depends on $P$.
Hence the lemma is proved when $N, \mu$ and in part (b) moreover $\ell$ and $K$ depend on $p$. But there are only finitely many induction steps, so that we can take the largest $N, \mu$, $\ell$ and $K$. We start the induction when $p=M$, $\mu=\lambda$ and $m=0$. Therefore, the lemma is true for all p with constants $N$ (depending on $P$ ), $\mu$ (depending on $\lambda$ ) and in part (b) \& (depending on $k$ ) and K (depending on $\lambda$ and $k$ ).

## Now we are able to prove theorems 7.1 and 7.2.

PROOF OF THEOREM 7.1. It follows from theorem A15 that for all s $\in I_{0}$ we can take $f=\mathrm{Pg}_{S}$ in $U_{S}^{(0)} \in U^{(0)}$ with $g \in A\left(U_{S}^{(0)}\right)^{q}$. As in the proof of lemma 7.4 we set $\nu \geq{ }^{2} \log (t+1)$, so that $(t+1)$ times $U_{S}^{(\nu)}$ is contained in $U_{S}^{(0)}$ for all $s \in I_{v}$, where $s^{\prime}=\rho_{0, v} s \in I_{0}$. As in (7.7) we can find $\tilde{g} \epsilon$
 of (7.7)
where $f$ is regarded as a cocycle in $C^{0}\left[U^{(0)}, A^{p}, \phi\right]$. Consider the differences c of the functions $\tilde{g}_{S}$ in the overlaps of the sets $U_{S}^{(\nu)}$ for $s \in I_{V}$, that is $c=\delta \tilde{g}$. Since not more than a finite number of different $s$ are mapped by $\rho_{0, v}$ onto the same $s^{\prime}$, there are constants $C_{m}^{\prime}$ with

$$
\|c\|_{\phi_{N_{1}}, m} \leq C_{m}^{\prime} \int_{\Omega_{\ell(m)}}|f(z)|^{2} \exp -\phi(z) d \lambda(z)
$$

Then $P c=P \delta \tilde{g}=\delta f=0$ and also $\delta c=0$, hence $c$ is a cocycle in $C^{1}\left[U^{(\nu)}, R_{P}, \phi_{N_{1}}\right]$. According to lemma $7.4(a)$ there are $\nu>\mu, N_{2}$ and

$$
\begin{equation*}
c^{\prime} \in C^{0}\left[U^{(\mu)}, R_{P}, \phi_{N}\right], \tag{7.9}
\end{equation*}
$$

where $N=N_{1}+N_{2}$, with $\delta c^{\prime}=\rho_{\nu, \mu}^{*} c$ in $\Omega$ and according to lemma 7.4(b) there are moreover constants $K_{k}$ (also depending on $v$ ), integers $m>k$ (depending
on $k$ ) and cochains $c_{k}^{\prime} \in C^{0}\left(U_{k}^{(\mu)}, R_{P}, \phi_{N}\right)$ with $\delta c_{k}^{\prime}=\rho_{\nu, \mu}^{*} c$ in $\Omega_{k}$ and with

$$
\begin{equation*}
\left\|c_{k}^{\prime}\right\|_{\phi_{N}}, k \leq K_{k}\|c\|_{\phi_{N_{1}}}, m \leq K_{k} C_{m}^{\prime} \int_{\Omega_{\ell}}|f(z)|^{2} \exp -\phi(z) d \lambda(z), \tag{7.1}
\end{equation*}
$$

where $\ell>m>k$ depends on $k$.
Finally for all $s \in I_{\mu}$ we put $v_{s}(z)=\tilde{g}_{s},(z)-c_{s}^{\prime}(z)$ for $z \in U_{s}^{(\mu)}$ with $s^{\prime}=\rho_{\nu, \mu} s$ which by (A9) defines a function $v \in A(\Omega){ }^{q}$ because $\left\{v_{s} \mid s \in I_{\mu}\right\} \in$ $\epsilon C^{0}\left(U^{(\mu)}, A^{q}\right)$ and $\delta v=\rho_{\nu, \mu}^{*} \delta \tilde{g}-\rho_{\nu, \mu}^{*} c=0$. Furthermore for all $k$

$$
\int_{\Omega_{k}}|v(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z) \leq\|v\|_{\phi_{N}}, k<\infty
$$

by (7.7) and (7.9). Similarly, for each $k\left(v_{k}\right)_{s}(z)=\tilde{g}_{s}(z)-\left(c_{k}^{\prime}\right)_{s}(z)$ for $z \in U_{s}^{(\mu)} \cap \Omega_{k}$ defines a function $v_{k} \in A\left(\Omega_{k}\right)^{q}$ and there are constants $K_{k}$ and $\ell_{k}>k$ with

$$
\int_{\Omega_{k}}\left|v_{k}(z)\right|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z) \leq K_{k} \int_{\Omega_{\ell_{k}}}|f(z)|^{2} \exp -\phi(z) d \lambda(z)
$$

by (7.8) and (7.19). Moreover, for all $s \in I_{\mu}$ in $U_{S}^{(\mu)}$ we have

$$
P v=P v_{s}=P \tilde{g}_{s},-P c_{s}^{\prime}=f,
$$

so that $P v=f$ in $\Omega$ and similarly $\mathrm{Pv}_{\mathrm{k}}=\mathrm{f}$ in $\Omega_{\mathrm{k}}$.
PROOF OF THEOREM 7.2. Let $F$ be the sheaf $P A^{q}$. We construct a cochain $h \in$ $\in C^{0}\left[U^{(0)}, A^{p}, \phi\right]$ as follows: for all $s \in I_{0}$, when $U_{s}^{(0)} \subset \Omega_{1}$ we define $h_{s}(z)=$ $=f_{1}(z)$ for $z \in U_{S}^{(0)}$; for $k=1,2, \ldots$ successively, when $U_{s}^{(0)} \cap \Omega_{k} \neq \emptyset$, $\mathrm{U}_{\mathrm{s}}^{(0)} \notin \Omega_{\mathrm{k}}$, let $\ell$ be the smallest integer with $\mathrm{U}_{\mathrm{s}}^{(0)} \subset \Omega_{\ell}$ or when $\mathrm{U}_{\mathrm{s}}^{(0)} \subset$ $\subset \Omega_{k+1} \cap \Omega_{k}^{c}$, let $\ell$ be $\ell=k+1$, then we define $h_{s}(z)=f_{\ell}(z)$ for $z \in U_{s}^{(0)}$. By (A15)(ii) we obtain for all $k$

$$
\|h\|_{\phi, k} \leq M \max _{1 \leq j \leq \ell} \int_{\Omega_{j}}\left|f_{j}(z)\right|^{2} \exp -\phi(z) d \lambda(z)<\infty,
$$

where $\ell=\ell(k)$ depends on $k$ according to (A15)(iii).

Since $f_{k+m}-f_{k}=P\left(g_{k+m-1}+\ldots+g_{k}\right)$ in $\Omega_{k}$ for all $m \geq 1$, the differences of the functions $h_{s}$ in the overlaps $U_{S_{1} s_{2}}^{(0)}$ of the sets $U_{S}^{(0)}$ are either zero or $\mathrm{Pg}_{\mathrm{s}_{1} \mathrm{~s}_{2}}$ for some $\mathrm{g}_{\mathrm{s}_{1} \mathrm{~s}_{2}} \in \mathrm{~A}\left(\mathrm{U}_{\mathrm{s}_{1} \mathrm{~s}_{2}}^{(0)} \mathrm{q}^{\mathrm{q}}\right.$. Hence $\delta \mathrm{h} \in \mathrm{C}^{1}\left(U^{(0)}, F\right)$.

Now theorem 7.2 follows from the next theorem and theorem A15.
THEOREM 7.3. Let $F$ be the sheaf $P A^{q}$ in the pseudoconvex set $\Omega$, where $\Omega$ is the union of the subsets $\Omega_{k}$ satisfying (7.1) and let $\phi$ be a plurisubharmonic function in $\Omega$ satisfying (7.2). If for some $\lambda h \in C^{0}\left[U^{(\lambda)}, A^{p}, \phi\right]$ with $\delta \mathrm{h} \in \mathrm{C}^{1}\left(U^{(\lambda)}, F\right)$, then there is a constant N and a function $\mathrm{v} \in \mathrm{A}(\Omega)^{\mathrm{p}}$ with for all $s \in I_{\lambda} v(z)-h_{s}(z)=P(z) g_{s}(z)$ for $z \in U_{S}^{(\lambda)}$ and for some $g \epsilon$ $\in C^{0}\left(U^{(\lambda)}, A^{q}\right)$ and $w i t h$

$$
\int_{\Omega_{k}}|v(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z)<\infty \quad \text { for all } k=1,2, \ldots .
$$

PROOF. We can estimate the cocycle $\mathrm{f}=\delta \mathrm{h} \in \mathrm{C}^{1}\left(U^{(\lambda)}, F\right)$ in terms of $h$ by use of (A15) (iv), hence $f \in C^{1}\left[U^{(\lambda)}, F, \phi\right]$ and $\delta f=0$. According to 1emma 7.4(a) there is a cochain $f^{\prime} \in C^{0}\left[U^{(\mu)}, F, \phi_{N}\right]$ with $\delta f^{\prime}=\rho_{\lambda, \mu}^{*} f$ in $\Omega$ for some integer $N$ and $\mu>\lambda$.

Let for all $i \in I_{\mu}$ and $z \in U_{i}^{(\mu)}$

$$
v_{i}(z)=h_{s^{\prime}}(z)-f_{i}^{\prime}(z)
$$

where $s^{\prime}=\rho_{\lambda, \mu}(i)$. Then $\delta v=\rho_{\lambda, \mu}^{*} \delta \mathrm{~h}-\delta \mathrm{f}^{\prime}=\rho_{\lambda, \mu}^{*} \mathrm{f}-\rho_{\lambda, \mu}^{*} \mathrm{f}=0$ in $\Omega$, thus $\left\{v_{i} \mid i \in I_{\mu}\right\}$ determines a function $v \in(\Omega)^{p}$. Moreover, using (A15)(vi) we obtain for all k

$$
\int_{\Omega_{k}}|v(z)|^{2} \frac{\exp -\phi(z)}{\left(1+\|z\|^{2}\right)^{N}} d \lambda(z) \leq\|v\|_{\phi_{N}}, k \leq M_{\lambda, \mu}(k)\|h\|_{\phi, k}+\|f\|_{\phi_{N}}, k<\infty .
$$

For $s \in I_{\lambda}$ let $I^{\prime}(s) \subset I_{\mu}$ be the set of those $i \in I_{\mu}$ with $V_{i} \xlongequal{\text { def }} U_{i}^{(\mu)} n$ $\cap \mathrm{U}_{\mathrm{s}}^{(\lambda)} \neq \emptyset$. For all i $\in I^{\prime}(\mathrm{s})$ and $\mathrm{z} \in \mathrm{V}_{\mathrm{i}}$ we have

$$
v(z)-h_{s}(z)=h_{s^{\prime}}(z)-f_{i}^{\prime}(z)-h_{s}(z)
$$

Since $h_{S^{\prime}}{ }^{-h} h_{s} \in \Gamma\left(U_{S^{\prime}}^{(\lambda)} \cap U_{s}^{(\lambda)}, F\right)$ and also $f_{i}^{\prime} \in \Gamma\left(U_{i}^{(\mu)}, F\right)$, we obtain

$$
v-\left.h_{s}\right|_{V_{i}} \in \Gamma\left(V_{i}, F\right)
$$

As $\mathrm{V}_{\mathrm{i}}$ is pseudoconvex, theorem A15 yields

$$
\left.v^{v-h_{s}}\right|_{V_{i}} \in P \Gamma\left(v_{i}, A^{q}\right)
$$

and again by theorem $A 15 \mathrm{v}-\mathrm{h}_{\mathrm{S}}=\mathrm{Pg}_{\mathrm{S}}$ in $\mathrm{U}_{\mathrm{S}}^{(\lambda)}$ for some $\mathrm{g}_{\mathrm{S}} \in \Gamma_{\mathrm{S}}\left(\mathrm{U}_{\mathrm{S}}^{(\lambda)}, \mathrm{A}^{\mathrm{q}}\right)$, because also $U_{S}^{(\lambda)}$ is pseudoconvex (property (A15)(i)).
8. COROLLARIES AND EXAMPLES

In this section some corollaries and particular cases are given of the theorem on Fourier transforms in section 6, theorem 6.1. We can survey this theorem by: let

$$
\begin{equation*}
\operatorname{Exp}(a+0 ; C) \xlongequal{\text { def }} \operatorname{proj} \underset{k \rightarrow \infty}{ } \lim A_{\infty}(\exp (-a ̃(z)-1 / k\|z\|) ; C(k)) \tag{8.1}
\end{equation*}
$$

and

$$
\begin{equation*}
A(a+0 ; C) \stackrel{\text { def }}{=} \underset{k \rightarrow \infty}{\operatorname{ind}} \lim A_{\infty}(\exp 1 / k\|\zeta\| ; \Omega(a+1 / k ; C)) \tag{8.2}
\end{equation*}
$$

then

$$
\begin{equation*}
F A(a+0 ; C)^{\prime}=\operatorname{Exp}(a+0 ; C) \tag{8.3}
\end{equation*}
$$

where $A(a+0 ; C)^{\prime}$ is the strong dual of $A(a+0 ; C)$ and $F$ is an isomorphism; $\Omega$ is given by formula (5.4). Here we have used the fact that the sequences of weightfunctions

$$
\left\{\exp \delta_{k} / k\|\zeta\|\right\}_{k=1}^{\infty} \quad \text { and } \quad\{\exp 1 / k\|\zeta\|\}_{k=1}^{\infty}
$$

induce the same topology on the space $A(a+0 ; C)$.
Let $w \in \operatorname{pr} C$, then $w \in \operatorname{pr} C_{k}$ for some $k$; since $\Omega(a ; C) \cap\left(C_{k+1}^{*}\right)^{c}$ is
bounded and since for $\zeta \in C_{k+1}^{*}$

$$
\delta\|\zeta\| \leq \operatorname{Im} \mathrm{W} \cdot \zeta \leq\|\zeta\|
$$

for some $\delta>0$, also the system
(8.4) $\quad\left\{\exp \frac{1}{\mathrm{k}} \operatorname{Im} \mathrm{w} \cdot \zeta\right\}_{\mathrm{k}=1}^{\infty}$
induces the same topology on $A(a+0 ; C)$. These weightfunctions $M_{k}$ satisfy $M_{k}=\exp -\phi_{k}$, where $\phi_{k}(\zeta)=-1 / k \operatorname{Im} w \cdot \zeta$ is a $p 1 u r i s u b h a r m o n i c$ function. Therefore, the theorems of the appendix and of section 7 may be applied to the space $A(a+0 ; C)$, because all the $L^{P}$-norms are equivalent, $p=1,2, \ldots, \infty$ (the space $A(a+0 ; C)$ is nuclear (see [14] G.7)).

It follows from (8.3) and (8.4) that any $f \in \operatorname{Exp}(a+0 ; C)$ satisfies for all $\varepsilon>0$ and $\delta>0$

$$
\begin{aligned}
|f(z)| & =\left|<\mu_{\zeta}, e^{i z \cdot \zeta}>\left|\leq M_{\varepsilon, \delta} \sup _{\zeta \in \Omega(a+\delta ; C)}\right| e^{i(z-\varepsilon w) \cdot \zeta}\right| \leq \\
& \leq M_{\varepsilon, \delta} e^{\tilde{a}(z-\varepsilon w)+\delta\|z-\varepsilon w\|}
\end{aligned}
$$

when $z \in \varepsilon w+C$. Now let a be bounded on pr $C$, then a can be continued as a continuous function to $\overline{\mathrm{prC}}$ and thus a is uniformly continuous on pr C. That means that for all $\delta>0$, there is a $\varepsilon(\delta)>0$ with for $\varepsilon \leq \varepsilon(\delta)$

$$
|a(\overparen{z-\varepsilon W})-a ̃(z)|<\delta .
$$

Hence

$$
\tilde{a}(z-\varepsilon w) \leq a(\tilde{z})\|z-\varepsilon w\|+\delta\left\|_{z-\varepsilon w \|} \leq \delta \varepsilon(\delta)+\varepsilon(\delta) \sup _{z \in \operatorname{prC}}|a(z)|+\tilde{a}(z)+\delta\right\|_{z \|},
$$

so that $f$ satisfies for $z \in \varepsilon W+C$, all $\varepsilon>0$ and $\delta>0$

$$
|f(z)| \leq M_{\varepsilon, \delta}^{\prime} e^{\tilde{a}(z)+2 \delta\|z\|}
$$

We can choose $w \in \operatorname{prC}$ so, that the sets $\{\varepsilon w+C\}{ }_{\varepsilon>0}$ are just the subsets of $C$ consisting of all the points of $C$ with distance larger than $\eta$ to the boundary of $C, \eta>0$ and $\eta \rightarrow 0$ if $\varepsilon \rightarrow 0$. Thus we have found that as sets

$$
\begin{equation*}
\operatorname{Exp}(a+0 ; C)=\underset{k \rightarrow \infty}{\operatorname{proj}} \lim _{\infty}(\exp (-a ̃(z)-1 / k\|z\|) ; 1 / k w+C) \tag{8.5}
\end{equation*}
$$

when a is bounded on pr C. Since the topology defined by (8.1) is obviously weaker than the one defined by (8.5) and since both topologies turn Exp (a+0;C) into an F S-space, both topologies coincide (see [15] corollary 2 to th.7.1), so that (8.5) also holds for the topologies. A similar property holds for the spaces $H(a ; C)$ and $H^{*}(a ; C)$ of [14] provided that then $a$ is uniformly continuous on $C$, which is true when a is not vertical at the boundary of pr C (see section 4), for example when a is constant. This surprising property of functions of exponential type in cones is difficult to establish without Fourier transformation.

Another surprising corollary is that, as topological spaces, $\overline{\mathrm{A}}_{\mathrm{i}}=\overline{\mathrm{A}}_{\mathrm{ii}}$, as we have already seen. It means that it does not make a difference if we use $\Omega(a+1 / k ; C)$ or $\Omega\left(a+1 / k ; C_{k}\right)$ in the space $A(a+0 ; C)$ : any function $\phi$ holomorphic in int $\Omega(a+1 / k ; C)$ with

$$
|\phi(\zeta)| \leq M_{1} \exp -1 / k\|\zeta\|, \quad \zeta \in \text { int } \Omega(a+1 / k ; C)
$$

is holomorphic in some larger set int $\Omega\left(a+1 / m ; C_{m}\right) \cup$ int $\Omega(a+1 / k ; C)$ and satisfies there

$$
|\phi(\zeta)| \leq M_{2} \exp -1 / \ell\|\zeta\|
$$

for some $\ell \geq k$ depending on $k$ and $C_{m}$ and some $M_{2}$ depending on $M_{1}$, $k$ and $C_{m}$, but not on $\phi$.

Now we imagine an open convex set $\Omega$ in $\mathbb{C}^{n}$ being given or equivalently an open convex cone $C$ in $\mathbb{C}^{n}$ and a convex homogeneous function a in $G$ ( $\Omega$ such that it does not contain a straight line, whence the cone $C$ is open in $\mathbb{C}^{n}$ ). Let $\left\{\Omega_{\mathrm{m}}\right\}_{\mathrm{m}=1}^{\infty}$ be an increasing sequence of closed convex sets with union $\Omega$ and such that the points of $\Omega_{\mathrm{m}}$ are those points in $\Omega$ with distance larger than
$\varepsilon_{\mathrm{m}}$ from $\partial \Omega$ (see section 4). The sets $\Omega_{\mathrm{m}}$ determine convex homogeneous functions $\tilde{a}_{\mathrm{m}}$ on C with for some $\eta_{m} \geq \varepsilon_{m}$

$$
a(z)-\eta_{m} \leq a_{m}(z) \leq a(z)-\varepsilon_{m}, \quad z \in \operatorname{prC}
$$

$\varepsilon_{m}>\varepsilon_{m+1}>0, \varepsilon_{m} \rightarrow 0$ and $\eta_{m}>\eta_{m+1}>0, \eta_{m} \rightarrow 0$ for $m \rightarrow \infty$. We define
(8.6). $\tilde{\operatorname{Exp}}(a ; C)=\underset{m \rightarrow \infty}{\operatorname{ind}} \lim \operatorname{Exp}\left(a_{m}+0 ; C\right)$,
where we may use (8.5) instead of (8.1) when a is bounded. An equivalent definition is

$$
\tilde{\operatorname{Exp}}(\mathrm{a} ; \mathrm{C})=\underset{\mathrm{m} \rightarrow \infty}{\operatorname{ind}} \lim \underset{\mathrm{k} \rightarrow \infty}{\operatorname{proj}} \lim A_{\infty}\left(\exp -a_{m}(z) ; C(k)\right) .
$$

We also define

$$
\tilde{A}^{\prime}(a ; c)=\underset{m \rightarrow \infty}{\operatorname{ind}} \lim A\left(a_{m}+0 ; C\right)^{\prime}
$$

or equivalently

$$
\tilde{A}^{\prime}(a ; C)=\underset{m \rightarrow \infty}{\operatorname{ind} \lim }\left[\underset{\mathrm{k} \rightarrow \infty}{\operatorname{ind} \lim A_{\infty}\left(\exp 1 / k\|\zeta\| ; \Omega\left(a_{m} ; C\right)\right)}\right]^{\prime} .
$$

It easily follows that $F$ is an isomorphism:
(8.7) $\quad \operatorname{Exp}(a ; C)=F \tilde{A}^{\prime}(a ; C)$.

Exp and $\tilde{A}^{\prime}$ are inductive limits of nuclear Fréchet spaces, so they are nuclear themselves.

In particular we may take for the cone in $\mathbb{C}^{n}$ a tubular radial domain $T^{C} \subset C^{n}$, where $T^{C}=\mathbb{R}^{n}+i C$ with now $C$ an open convex cone in $\mathbb{R}^{n}$. A relatively compact subcone of this domain is $\mathbb{R}^{n}+i C_{k}$ with $C_{k} \subset \subset C$ and the domains C(k) become

$$
\left\{\mathbb{R}^{\mathrm{n}}+\mathrm{i} C_{k}\right\} \cap\{\mathrm{z} \mid\|\operatorname{Im} \mathrm{z}\|>1 / \mathrm{k}\}
$$

see section 4. Let $a_{C}(y, x)$ be a convex homogeneous function on $T^{C}$ which is bounded on each pr $T^{C_{k}}$. Then $a(0, x)$ exists and is finite; so $A=$ $=\max _{\| \|=1} a(0, x)<\infty$. Then the domain $\Omega\left(a ; T^{C}\right) \xlongequal{\text { not }} \Omega(a ; C)$ is bounded in the imaginary direction, that is $\Omega(a ; C) \subset \mathbb{R}^{n}+i B_{A}$, where $B_{A}$ is the ball with radius $A$ in $\mathbb{R}^{n}$. This case will be used in the next section, where the Newton interpolation series will be derived.

We can consider boundary values of functions $f$ holomorphic in $\mathbb{R}^{n}+i C$. When these are finite order distributions the function $f$ satisfies

$$
\begin{equation*}
|f(z)| \leq M_{k}\left(1+\|y\|^{-m}\right), \quad y \in C_{k},\|y\| \leq 1 \tag{8.8}
\end{equation*}
$$

for some $m$ depending on $f$. When moreover $f \in \tilde{E x p}\left(a ; T^{C}\right), f$ is the Fourier transform of an analytic functional in $Z^{\prime}$ (see [14] H.4) carried by $\Omega(a ; C)$. Indeed, in the same way as theorem 6.1 was obtained, using polynomials as weightfunctions instead of (8.4) one can show

$$
\begin{equation*}
D_{F}^{\prime}(a ; C)=F Z^{\prime}(a ; C) \tag{8.9}
\end{equation*}
$$

with $Z^{\prime}(a ; C)$ the dual of

$$
Z(a ; C)=\underset{m \rightarrow \infty}{\operatorname{proj}} \lim A_{\infty}\left((1+\|\zeta\|)^{\mathrm{m}} ; \Omega\left(\mathrm{a}_{\mathrm{m}} ; \mathrm{C}\right)\right)
$$

and with

$$
\mathcal{D}_{\mathrm{F}}^{\prime}(\mathrm{a} ; \mathrm{C})=\underset{\mathrm{m} \rightarrow \infty}{\operatorname{ind} \lim } \underset{\mathrm{proj}}{\lim } \lim _{\infty}\left(\frac{\exp -\mathrm{a}_{\mathrm{m}}(\mathrm{z})}{1+\| \|^{-m}} ; \mathbb{R}^{\mathrm{n}}+i C_{k}\right),
$$

where $\mathbb{R}^{n}+i C_{k}$ may be replaced by $\left\{\mathbb{R}^{n}+\mathrm{i} C_{k}\right\} \cup\left\{\mathbb{R}^{\mathrm{n}}+\mathrm{i}\left(1 / \mathrm{k} \mathrm{y}_{0}+\mathrm{C}\right)\right\}$ with $\mathrm{y}_{0} \in \operatorname{pr} \mathrm{C}$, when a is bounded on pr C. $Z^{\prime}(a ; C)$ and hence also $D_{F}^{\prime}(a ; C)$ is a nuclear LSspace, so that for example they are reflexive (compare the spaces in section 6 [14]).

Finally we give three examples with $\Omega$ contained in $\mathbb{R}^{\mathrm{n}}$ illustrating the differences between analytic functionals and distributions. For simplicity we assume that the function a is constant on prC , so that (8.4) holds.

Firstly, let $f$ be holomorphic in $\mathbb{R}^{n}+i C$ and satisfy for $a l l \varepsilon>0$ and $k$

$$
|f(z)| \leq K(\varepsilon, k) e^{(a+\varepsilon)\|y\|+\varepsilon\|x\|}, \quad y \in C_{k},\|y\|>\varepsilon
$$

then f satisfies these inequalities (with other constants $\mathrm{K}(\varepsilon)$ ) also for $y \in \varepsilon y_{0}+C, y_{0} \in \operatorname{prC}$ and $f=F(\sigma)$ with $\sigma$ carried by

$$
\Omega=\{\zeta \mid n=0,-y \cdot \xi \leq a\|y\|, y \in C\} \subset \mathbb{C}^{\mathrm{n}}
$$

such that for all $\varepsilon>0 \quad \sigma$ can be represented as a measure $\sigma_{\varepsilon}$ in

$$
\Omega_{\varepsilon}=\{\zeta \mid\|\eta\| \leq \varepsilon,-y \cdot \xi \leq(a+\varepsilon)\|y\|, y \in C\}
$$

with

$$
\int_{\Omega} \exp -\varepsilon\|\zeta\| \quad\left|d \sigma_{\varepsilon}(\zeta)\right|<\infty
$$

Secondly, let $f$ satisfy for all $\varepsilon>0$ and $k$ and some $m$

$$
|f(z)| \leq K(\varepsilon, k) e^{(a+\varepsilon)\|y\|+\varepsilon\|x\|}\left(a+\|y\|^{-m}\right), \quad y \in C_{k},
$$

then f satisfies these inequalities (with other constants $\mathrm{K}(\varepsilon, \mathrm{k})$ ) also for $y \in C_{k} \cup\left\{\varepsilon y_{0}+C\right\}$ and $f=F(\mu)$ with $\mu \in Z^{\prime}$ carried by $\Omega$, that is for all $\varepsilon>0 \mu$ can be represented as a measure $\mu_{\varepsilon}$ in $\Omega_{\varepsilon}$ with

$$
\int_{\Omega}(1+\|\zeta\|)^{-\ell}\left|d \mu_{\varepsilon}(\zeta)\right|<\infty
$$

for some $\ell>m$ (actually $\ell=m+n+2$, see [14] (6.10)).
Finally, let $f$ satisfy for $a l l \varepsilon>0$ and $k$ and some $m$

$$
|f(z)| \leq K(\varepsilon, k) e^{(a+\varepsilon)\|y\|}(1+\|x\|)^{m}\left(1+\|y\|^{-m}\right), \quad y \in C_{k},
$$

then $f$ satisfies these inequalities (with other constants $K(\varepsilon, k)$ and $(1+\|x\|)^{m}$ replaced by $(1+\|x\|)^{\ell}$ for some $\left.\ell>m\right)$ also for $y \in C_{k} \cup\left\{\varepsilon y_{0}+C\right\}$ and $\mathrm{f}=\mathrm{F}(\mathrm{g})$ with $\mathrm{g} \in \mathrm{S}^{\prime}$ having its support contained in $0=\{\xi \mid-y \cdot \xi \leq a l l y$, $y \in C\}$, so that $g$ can be represented as a finite combination of derivations
of measures $\mathrm{g}_{\mathrm{j}}$ in 0 with

$$
\int_{0}(1+\|\xi\|)^{-l}\left|\operatorname{dg}_{\mathrm{j}}(\xi)\right|<\infty
$$

for $j=0,1, \ldots, \ell(\ell=m+n+2$, see [14] (6.10)).
As in section 2, in the first two examples we have when $b>a \geq 0$

$$
\operatorname{Exp}\left(a+0 ; \mathrm{T}^{\mathrm{C}}\right) \xrightarrow{\text { dense }} \operatorname{Exp}\left(\mathrm{b}+0 ; \mathrm{T}^{\mathrm{C}}\right),
$$

while in the third example

$$
\mathrm{H}(\mathrm{a} ; \mathrm{C}) \xrightarrow{\text { closed linear subspace }} \mathrm{H}(\mathrm{~b} ; \mathrm{C}) \text {. }
$$

Even, since the restriction map from $A(b+0 ; C)$ into $S(a ; C)$ is injective, we have

$$
\mathrm{H}(\mathrm{a} ; \mathrm{C}) \xrightarrow{\text { dense }} \operatorname{Exp}\left(\mathrm{b}+0 ; \mathrm{T}^{\mathrm{C}}\right) .
$$

9. NEWTON SERIES FOR FUNCTIONS HOLOMORPHIC IN TUBULAR RADIAL DOMAINS

In this section we derive the Newton interpolation series for functions in $\tilde{\operatorname{Exp}}\left(a ; \mathbb{T}^{\mathrm{C}}\right)$. We give the most general class of holomorphic functions for which the Newton series is valid for $h$ in a convex cone $C$ in $\mathbb{R}^{n}$. However, since the detailed description becomes quite complicated, we discuss a particular case, namely a class of holomorphic functions of constant exponential type and we give a uniform bound on the length of $h$. The bound for $\|\mathrm{h}\|$ will not be the best possible, but still this case gives a good idea of the generalization of the validity of the Newton series discussed in this paper. Finally we make some general remarks on the validity of the Newton series.

In [10] KIOUSTELIDIS derived the Newton interpolation series (and similar series) with the aid of Fourier transformation. The adventage of this method against the classical one (Cauchy's integral formula, NÖRLUND
[13], GELFOND [5]) is that it treats the case of several variables as well. However, his treatment is valid only for entire functions. This is not a restriction due to the method, for, as we have shown here, one has to extend the method (namely the formalism of Fourier transformation) to non-entire functions. Then we are able to derive the Newton series (and the similar series of KIOUSTELIDIS [10] or [14] remark 10.1) in several variables for non-entire functions as well. Moreover, in some way we obtain the largest possible class, for which the formalism is valid, since we use the domain of convergence completely (that is we do not cut off a compact subset of this domain as it is done in [10]) and since outside this domain the formalism is not valid, see Satz 5 in [10].

As we have seen in [14], section 5, we have to restrict the vector $h$ to a real open convex cone $C$ in $\mathbb{R}^{n}$ in order to get the Newton series for non-entire functions. Moreover, let $\|h\|$ be bounded by a positive number $b$. Let the convex (unbounded) open set $\Omega$ in $\mathbb{C}^{n}$ be the interior of one of the components of

$$
\left\{\zeta\left|\zeta \in \mathbb{C}^{\mathrm{n}},\left|\mathrm{e}^{-\mathrm{h} \cdot \zeta}-1\right|<1, \forall \mathrm{~h} \in \mathrm{C} \text { with }\|\mathrm{h}\| \leq \mathrm{b}\right\}\right.
$$

see figure 4.1 of [14]. Then $\Omega$ is bounded in the imaginary direction, because $|h \cdot n| \leq\left(2 k+\frac{1}{2}\right) \pi$ for some $k$ and for $a l l h \in C$ with $\|h\| \leq b$ and also $\Omega$ is contained in

$$
\begin{equation*}
\{\zeta \mid-h \cdot \xi \leq \log 2, \forall h \in C \text { with }\|\mathrm{h}\| \leq \mathrm{b}\} . \tag{9.1}
\end{equation*}
$$

Hence $\Omega$ determines the convex cone $\mathbb{R}^{n}+i C$ in $\mathbb{C}^{n}$ and the convex homogeneous function $H_{\Omega}$ on $\mathbb{R}^{n}+i C$ by

$$
\begin{equation*}
\mathrm{H}_{\Omega}(z)=\sup _{\zeta \in \Omega}-\operatorname{Im} z \cdot \zeta \cdot \tag{9.2}
\end{equation*}
$$

$H_{\Omega}(z)$ is continuous up to $y=0$, that is $H_{\Omega}(x)$ exists for $x \in \mathbb{R}^{n}$ and it follows from (9.1) that $H_{\Omega}(\tilde{z})$ is bounded by ( $\left.\log 2\right) / b+B$, where $B$ is a bound for $\|\eta\|$, thus (8.5) may be applied. Also we have, see (4.2), (4.3), (4.4)

$$
\Omega=\operatorname{int}\left\{\zeta \mid-\operatorname{Im} z \cdot \zeta \leq H_{\Omega}(z), z \in \mathbb{R}^{\mathrm{n}}+\mathrm{iC} C .\right.
$$

Let $\Omega_{\mathrm{m}}$ be an increasing sequence of convex closed subsets of $\Omega$ such that some $\varepsilon$-neighborhood $\Omega_{\mathrm{m}, \varepsilon}$ of $\Omega_{\mathrm{m}}$ is contained in $\Omega$ and $\Omega=\bigcup_{\mathrm{m}}^{\infty}{ }_{1}^{\infty} \Omega_{\mathrm{m}}$. Let $H_{m}$ be the functions $H_{\Omega_{m}}, m=1,2, \ldots$. For $y \in C, h \in C$ let $s \in \mathbb{C}$ be such that

$$
z+i s h \in \mathbb{R}^{\mathrm{n}}+\mathrm{i} C,
$$

so that Res $\geq-\alpha$ for some non-negative number $\alpha$ depending on $y$ and $h$.
LEMMA 9.1. Let $z \in \mathbb{R}^{n}+i C, h \in C$ and $s \in \mathbb{C}$ as above. Then the sequence

$$
\phi_{N, z}(\zeta)=e^{i z \cdot \zeta} \sum_{k=0}^{N}\binom{\mathrm{~s}}{k}\left(e^{-h \cdot \zeta}-1\right)^{k}
$$

converges for $N \rightarrow \infty$ to exp $i(z+i s h) \cdot$ in all the spaces $A\left(H_{m}+0 ; C\right)$, $\mathrm{m}=1,2, \ldots$.

PROOF. The space $A\left(H_{m}+0 ; C\right)$ is defined by (8.2) so that according to (5.9) $\exp i z \cdot \zeta \in A\left(H_{m}+0 ; C\right)$ when $z \in \mathbb{R}^{n}+i C$, hence $\exp i(z+i s h) \cdot \zeta$ and $\exp (i z \cdot \zeta-k h \cdot \zeta), k=0,1,2, \ldots$, belong to $A\left(H_{m}+0 ; C\right)$ for all m. Since $A\left(H_{m}+0 ; C\right)$ is a Fréchet space, we have to show that for some $\varepsilon>0$ small enough

$$
\begin{equation*}
\sup _{\zeta \in \Omega_{\mathrm{m}, \varepsilon}}\left|\phi_{\mathrm{N}, \mathrm{z}}(\zeta)\right| \exp \varepsilon\|\zeta\| \leq \mathrm{K}, \tag{9.3}
\end{equation*}
$$

where $K$ is independent of $N$.
In section 5 of [14] we defined subsets $\Omega(\varepsilon) \stackrel{\text { not }}{=} \Omega_{h}(\varepsilon)$ of

$$
\Omega_{h}=\left\{\zeta| | e^{-h \cdot \zeta}-1 \mid<1\right\}
$$

by

$$
\Omega_{h}(\varepsilon) \stackrel{\operatorname{def}}{=}\{\zeta \mid-h \cdot \xi<\log (2 \cos h \cdot n-\varepsilon)\}
$$

and we showed that for all $\varepsilon>0$ there is a $\varepsilon_{1}$ (here $\varepsilon_{1}=\varepsilon /(6 b)$ ) such that the $\varepsilon_{1}$-neighborhood of $\Omega_{h}(\varepsilon)$ is contained in $\Omega_{h}\left(\frac{1}{2} \varepsilon\right) \subset \Omega_{h}$. On the other hand, we will show that for all $\varepsilon>0$ there is a $\varepsilon_{2}$ such that the boundary of the $\varepsilon$-neighborhood of

$$
\begin{equation*}
\Omega\left(\varepsilon_{2}\right) \stackrel{\text { def }}{=} \bigcap_{\substack{h \in C \\\|h\| \leq b}} \Omega_{h}\left(\varepsilon_{2}\right) \tag{9.4}
\end{equation*}
$$

is contained in $\Omega^{c}$.
First let us remark that $\Omega_{h} \subset \Omega_{\beta h}$ when $\beta \leq 1$, since

$$
\beta \log (2 \cos x) \leq \log (2 \cos \beta x), \quad|x|<\frac{1}{2} \pi
$$

Then $\zeta \epsilon \partial \Omega\left(\varepsilon_{2}\right)$ means, that there is an $h \in C$, depending on $\zeta$ and $\delta$, with $\|h\|=b$ and with

$$
-h \cdot \xi \geq \log \left(2 \cos h \cdot n-\varepsilon_{2}\right)-\delta
$$

Now we choose $\varepsilon_{2}=\min \left(b^{2} \varepsilon^{2} / 16,1 / 171\right), \delta=\frac{1}{4} \varepsilon_{2}$ and

$$
\zeta_{0}=\zeta+i \operatorname{sign}(\sin h \cdot n) \frac{4}{b} \sqrt{\varepsilon_{2}} \tilde{h}
$$

where $\operatorname{sign}(0)=1$. Then $\left|\zeta-\zeta_{0}\right| \leq \varepsilon$ and for some integer $k$

$$
\frac{1}{2} \pi+4 \sqrt{\varepsilon_{2}}>\left|\operatorname{Im~h} \cdot \zeta_{0}+2 k \pi\right| \geq 4 \sqrt{\varepsilon_{2}}
$$

so that when $\left|\operatorname{Im~} h \cdot \zeta_{0}+2 k \pi\right| \xlongequal{\text { not }}|x|<\frac{1}{2} \pi$
$-\operatorname{Re} h \cdot \zeta_{0} \geq \log \left(2 \cos h \cdot n-\varepsilon_{2}\right)-\frac{1}{4} \varepsilon_{2}=$
$=\log \left\{2 \cos \left(|x|-4 \sqrt{\varepsilon_{2}}\right)-\varepsilon_{2}\right\}-\frac{1}{4} \varepsilon_{2} \geq \log \left(2 \cos \operatorname{Im} h \cdot \zeta_{0}\right)$,
for, $\varepsilon_{2} \leq 1 / 171$ implies $\sin 4 \sqrt{\varepsilon_{2}} \leq 63 / 16 \sqrt{\varepsilon_{2}}$, so that the following estimates with $|x|<\frac{1}{2} \pi$ hold

$$
\begin{aligned}
& 2 \cos \left(|x|-4 \sqrt{\varepsilon_{2}}\right)-\varepsilon_{2} \geq 2 \cos x-16 \varepsilon_{2} \cos x+2|\sin x| \sin 4 \sqrt{\varepsilon_{2}}-\varepsilon_{2} \geq \\
& \geq 2 \cos x-16 \varepsilon_{2}\left(1-\frac{1}{3} x^{2}\right)+2 \frac{4}{7} x \frac{63}{16} \sqrt{\varepsilon_{2}}-\varepsilon_{2}
\end{aligned}
$$

and the right hand side is larger than

$$
2 \cos x+\varepsilon_{2} \geq 2 \cos x \exp \frac{1}{4} \varepsilon_{2}
$$

if

$$
\frac{4}{9} \frac{16}{3} \sqrt{\varepsilon_{2}} x^{2}+2 x-8 \sqrt{\varepsilon_{2}} \geq 0
$$

which is true when

$$
x \geq 4 \sqrt{\varepsilon_{2}} \geq\left[-1+\left(1+\frac{64}{27} 8 \varepsilon_{2}\right)^{\frac{1}{2}}\right] /\left(\frac{64}{27} \sqrt{\varepsilon_{2}}\right)
$$

Hence $\zeta_{0} \notin \Omega$ and also when $\frac{1}{2} \pi \leq|x|<\frac{1}{2} \pi+4 \sqrt{\varepsilon_{2}}, \zeta_{0} \notin \Omega$. Thus the sets $\Omega(1 / \mathrm{m})$ defined by (9.4) with $\varepsilon_{2}$ replaced by $1 / \mathrm{m}$ may serve as the sets $\Omega_{\mathrm{m}}$.

From the formula above (5.3) in [14] we get for $\zeta \in \Omega_{\mathrm{m}, \varepsilon}$, $\varepsilon$ sufficiently small such that $\Omega_{\mathrm{m}, \varepsilon} \subset \Omega\left(\varepsilon_{1}\right)$ for some $\varepsilon_{1}$,

$$
\begin{array}{ll}
\left|\phi_{N, z}(\zeta)\right| \leq C_{1}\left(\varepsilon_{1}\right) \exp \alpha h \cdot \zeta \exp -\operatorname{Im} z \cdot \zeta, & \alpha>0 \\
\left|\phi_{N, z}(\zeta)\right| \leq C_{2}\left(\varepsilon_{1}\right)(1+\|\zeta\|) \exp -\operatorname{Im} z \cdot \zeta, & \alpha=0 .
\end{array}
$$

For $\zeta \in \Omega$ outside a compact set and $\varepsilon$ again sufficiently small

$$
(y-\alpha h) \cdot \xi \geq \varepsilon\|\xi\|
$$

hence there is a constant $K$ such that for $\zeta \in \Omega$

$$
-\operatorname{Im} z \cdot \zeta+\alpha h \cdot \xi \leq K-\varepsilon\|\zeta\|,
$$

since $\|n\|$ is bounded in $\Omega$. Now (9.3) follows when $\alpha>0$ and for $\alpha=0$ it follows by replacing $\varepsilon$ by $\frac{1}{2} \varepsilon$.

With the aid of lemma 9.1 and formula (8.3) the Newton series is derived for functions $f$ belonging to $\operatorname{Exp}\left(H_{\Omega}, T^{C}\right)$ given in (8.6), where $H_{\Omega}$ is defined by (9.2):

$$
f(z+i(s+\alpha) h)=\left\langle\mu_{\zeta}, e^{i z \cdot \zeta-\alpha h \cdot \zeta-s h \cdot \zeta}\right\rangle=
$$

$$
\begin{align*}
& =\left\langle e^{-\alpha h \cdot \zeta} \mu_{\zeta}, \lim _{N \rightarrow \infty} e^{i z \cdot \zeta} \sum_{k=0}^{N}\left(\begin{array}{l}
s
\end{array}\right)\left(e^{-h \cdot \zeta}-1\right)^{k}\right\rangle= \\
& =\sum_{k=0}^{\infty}\binom{s}{k}\left\langle e^{-\alpha h \cdot \zeta} \mu_{\zeta}, e^{i z \cdot \zeta}\left(e^{-h \cdot \zeta}-1\right)^{k}\right\rangle=  \tag{9.5}\\
& =\sum_{k=0}^{\infty}\binom{s}{k}\left\langle e^{-\alpha h \cdot \zeta}\left(e^{-h \cdot \zeta}-1\right)^{k}{ }_{\zeta}, e^{i z \cdot \zeta}\right\rangle=\sum_{k=0}^{\infty}\binom{s}{k} \Delta_{i h}^{k} f(z+i \alpha h)
\end{align*}
$$

valid for $z \in \mathbb{R}^{\mathrm{n}}+\mathrm{i} C, h \in \mathrm{C},\|\mathrm{h}\| \leq \mathrm{b}$, Re $\mathrm{s}+\alpha \geq 0$, $\alpha \geq 0$ arbitrary. The sequence

$$
\sum_{k=0}^{N}\binom{s}{k} e^{-\alpha h \cdot \zeta}\left(e^{-h \cdot \zeta}-1\right)^{k} \mu_{\zeta}
$$

converges weakly in $A\left(H_{m}+0 ; C\right)^{\prime}$ for some $m$ depending on $\mu$, which depends on $f$, and since $A\left(H_{m}+0 ; C\right)$ is a Montel space (see (8.2)), this sequence converges strongly in $A\left(H_{m}+0 ; C\right)$ ', hence according to (8.7) the series (9.5) converges in the topology of $\tilde{\operatorname{Exp}}\left(\mathrm{H}_{\Omega}, \mathrm{T}\right.$ ). Thus, reminding (8.5) we get, when f satisfies

$$
\begin{equation*}
\forall k, \forall y \in C_{k} \text { with }\|y\|>1 / k:|f(z)| \leq M_{k} \exp H_{m}(z), \tag{9.6}
\end{equation*}
$$

for Res $\geq-\alpha$ with $\alpha \geq 0, \mathrm{~h} \in \mathrm{C}$ with $\|\mathrm{h}\| \leq \mathrm{b}$ :

$$
\forall \varepsilon>0, \forall l>m, \forall p, \exists N_{0}(\varepsilon, l, p) \geq N_{1}(s), \forall z \in \mathbb{R}^{n}+i(1 / p w+C), \forall N \geq N_{0}
$$

$$
\begin{equation*}
\left|f(z+i(s+\alpha) h)-\sum_{k=0}^{N}\binom{s}{k} \Delta_{i h}^{k} f(z+i \alpha h)\right|<\varepsilon A(s) \exp H_{\ell}(z), \tag{9.7}
\end{equation*}
$$

where $N_{1}(s)$ is determined by (5.1) [14] and $A(s)$ by (5.4) [14].
Replacing $z+i \alpha h$ by $z$ in (9.5) we see that the Newton series

$$
\begin{equation*}
f(z+i s h)=\sum_{k=0}^{\infty}\binom{s}{k} \Delta_{i h}^{k} f(z) \tag{9.8}
\end{equation*}
$$

valid for $y \in n w+C, h \in C,\|h\| \leq b$, when $R e s \geq-\alpha, \alpha>0$ depending on $\eta>0$ and $h$, such that $y-\alpha h \in \delta w+C$ for some $\delta>0$, converges according to

$$
\forall \varepsilon>0, \forall \ell>m, \exists N_{0}(\varepsilon, \ell) \geq N_{1}(s), \forall z \in \mathbb{R}^{n}+i(n w+C), \quad \forall N \geq N_{0}
$$

$$
\begin{equation*}
\left|f(z+i s h)-\sum_{k=0}^{N}\binom{s}{k} \Delta_{i h}^{k} f(z)\right|<\varepsilon A(s) \exp H_{\ell}(z-i \alpha h) . \tag{9.9}
\end{equation*}
$$

We restate the results in
THEOREM 9.1. Let $\mathrm{h} \in \mathrm{C}$ with $\|\mathrm{h}\| \leq \mathrm{b}$ and let f be an element of $\underset{\operatorname{Exp}\left(\mathrm{H}_{\Omega}, \mathrm{T}^{\mathrm{C}}\right)}{ }$ where $H_{\Omega}$ is given in (9.2). If $\alpha>0$ is such that $y-\alpha h \in \delta w+C, \delta>0$, when $\mathrm{y} \in \mathrm{nw}^{+} \mathrm{C}$ for some $\eta>\delta$, then the Newton series (9.8) is valid for this y and $h$, when Re $s \geq-\alpha$. The series (9.8) converges absolutely in one of the norms of $\operatorname{Exp}\left(H_{\Omega}, \mathrm{T}^{\mathrm{C}}\right)$ or, more precisely, it converges according to (9.9). When $\operatorname{Re} s \geq-\alpha$ with $\alpha \geq 0$ arbitrary, the Newton series (9.5) holds for all $y \in \mathrm{C}, \mathrm{h} \in \mathrm{C},\|\mathrm{h}\| \leq \mathrm{b}$; then the series (9.5) converges absolutely in the topology of $\operatorname{Exp}\left(\mathrm{H}_{\Omega}, \mathrm{T}^{\mathrm{C}}\right)$ or, more precisely, it converges according to (9.7) when f satisfies (9.6). In both cases (9.5) and (9.8) converge uniformly in $s$ on compact subsets of $\{s \mid s \in C, \operatorname{Res} \geq-\alpha\}$.

Using (8.8) and (8.9) as in [14] section 7 we can derive the Newton series (9.5) for functions $f$ satisfying

$$
\forall k, \forall y \in C_{k}: \quad|f(z)| \leq M_{k}\left(1+\|y\|^{-m}\right) \exp H_{m}(z) .
$$

This series holds for $z \in \mathbb{R}^{n}+i C, h \in C$ and Re $s+\alpha \geq 0, \alpha \geq 0$ arbitrary and it converges in the topology of $D_{\mathrm{F}}^{\prime}(\mathrm{a} ; \mathrm{C})$, namely according to

$$
\begin{gathered}
\forall \varepsilon>0, \forall \ell>m, \forall p, \exists N_{0}(\varepsilon, \ell, p) \geq N_{1}(s), \forall z \in \mathbb{R}^{n}+i\left\{C_{p} \cup\{1 / p w+C\}\right\}, \\
\forall N \geq N_{0} \\
\left|f(z+i(s+\alpha) h)-\sum_{k=0}^{N}\binom{s}{k} f(z+i \alpha h)\right|<\varepsilon A(s)\left(1+\|y\|^{-t}\right) \exp H_{\ell}(z),
\end{gathered}
$$

where $t=m+n+2$ if $\alpha>0$ or $t=m+n+3$ if $\alpha=0$. This yields the convergence of the series (9.8) similarly to section 7 [14].

Actually, theorem 9.1 gives the condition $f$ should satisfy in order that the Newton series holds when $h$ ranges in a given domain. However, the function $H_{\Omega}(z)$ (formula (9.2)) arising in condition (9.6) is not given explicitely. This would be quite complicated (see [10] for entire functions and $h$ complex). Therefore, we now start with a given class of functions and determine the domain of $h$ the Newton series is valid in. For simplicity we will not give the largest possible domain, but still we get a considerable generalization of theorems 7.1 and 10.1 of [14].

The domain of convergence $\Omega_{h}=\left\{\zeta| | e^{-h \cdot \zeta}-1 \mid<1\right\}$ is determined by

$$
-h \cdot \xi<\log (2 \cos h \cdot n)
$$

for $h \in \mathbb{R}^{n}$, see 4.1 and figure 4.1 of $[14]$, here figure 9.1 when $k=0$.


Figure 9.1 gives the component of $\Omega_{h}$ that contains the origin. We approximate this domain from the inside by

$$
\begin{array}{r}
\left\{\zeta \left\lvert\,-\frac{1}{3} \pi h \cdot \xi \pm \log 2 h \cdot n<\frac{1}{3} \pi \log 2\right. \text { when }-h \cdot \xi>0 \text { and }|h \cdot n|<\frac{1}{3} \pi\right.  \tag{9.10}\\
\text { when }-h \cdot \xi \leq 0\} \subset \Omega_{h} .
\end{array}
$$

Now let a convex homogeneous function $a(z)$ be given on $\mathbb{R}^{n}+i C$ with $C$ a convex open cone in $\mathbb{R}^{n}$, such that $a(0, x)$ exists for $x \in \operatorname{pr} \mathbb{R}^{n}$. This function determines an open set $\Omega$ by

$$
\begin{equation*}
\Omega=\operatorname{int}\left\{\zeta \mid-\operatorname{Im} z \cdot \zeta \leq \tilde{a}(z), z \in \mathbb{R}^{n}+i C\right\} \tag{9.11}
\end{equation*}
$$

Let $\left\{\tilde{a}_{m}(z)\right\}_{m=1}^{\infty}$ be an increasing sequence convex homogeneous function with limit $\tilde{a}(z)$ and with $a_{m}(z)+\varepsilon_{m} \leq a(z), z \in \operatorname{pr} T^{C}$, for some $\varepsilon_{m}>0$. Let $\Omega_{m}$ be the domain determined by the function $a_{m}$. Then from (9.10) and (9.11) it follows that $\bar{\Omega} \subset \bar{\Omega}_{h}$ when $h \in C$ satisfies

$$
\begin{equation*}
\|h\| \leq \min \left\{\frac{\frac{1}{3} \pi \log 2}{\tilde{a}\left(\frac{1}{3} \pi \tilde{h}, \pm \log 2 \tilde{h}\right)}, \frac{\frac{1}{3} \pi}{a(0, \pm \tilde{h})}\right\} \tag{9.12}
\end{equation*}
$$

Hence in that case $\Omega_{\mathrm{m}} \subset \Omega_{\mathrm{h}}$ for all m and we obtain

COROLLARY 9.1. For functions $f \in \operatorname{Exp}\left(a ; \mathrm{T}^{\mathrm{C}}\right)$ the Newton series is valid, when $h \in C$ satisfies (9.12).

However, when a is a rather constant function, a better condition for $\|\mathrm{h}\|$ than (9.12) is obtained by approximating $\Omega_{\mathrm{h}}$ from the inside by

$$
\begin{gathered}
\left\{\zeta \mid(h \cdot \xi)^{2}+(h \cdot n)^{2}<\log ^{2} 2 \text { when }-h \cdot \xi>0 \text { and }|h \cdot n|<\log 2\right. \\
\text { when }-h \cdot \xi \leq 0\} \subset \Omega_{h} .
\end{gathered}
$$

This inclusion follows from $\log ^{2} 2 \leq \log ^{2}(2 \cos v)+v^{2},|v|<\frac{1}{2} \pi$, which is true because

$$
\begin{aligned}
& \log ^{2} 2-1 \log ^{2}(2 \cos \mathrm{v})=(\log 2-\log 2 \cos \mathrm{v})(\log 2+\log 2 \cos \mathrm{v}) \leq \\
& \leq \log 2 \cdot(2-2 \cos \mathrm{v}) \cdot 2 \log 2 \leq \mathrm{v}^{2} 2 \log 2 \leq 0.98 \mathrm{v}^{2} \leq \mathrm{v}^{2}
\end{aligned}
$$

For $\zeta \in \Omega_{\mathrm{m}}$ and $\mathrm{h} \in \mathrm{C}$ such that $-\mathrm{h} \cdot \xi>0$, we get

$$
\begin{aligned}
& (h \cdot \xi)^{2}+(h \cdot n)^{2}=-(-h \cdot \xi) h \cdot \xi-(-h \cdot n) h \cdot n \leq \tilde{a}_{m}((-h \cdot \xi) h,(-h \cdot n) h)= \\
& =\|h\|\left\{(h \cdot \xi)^{2}+(h \cdot n)^{2}\right\}^{\frac{1}{2}} a_{m}((-h \cdot \xi) h,(-h \cdot n) h),
\end{aligned}
$$

hence

$$
\left\{(h \cdot \xi)^{2}+(h \cdot n)^{2}\right\}^{\frac{1}{2}} \leq\|h\| a_{m}(\widetilde{\alpha, \beta h})
$$

for some $\alpha \in \mathbb{R}^{+}$and $\beta \in \mathbb{R}$. This is smaller than $\log 2$ when we require that

$$
\begin{equation*}
\|h\| \leq \min _{\substack{(\alpha, \beta) \in\left(\mathbb{R}^{+}, \mathbb{R}\right) \\ \alpha^{2}+\beta^{2}=1}} \frac{\log 2}{a(\alpha \tilde{h}, \beta \tilde{h})} . \tag{9.13}
\end{equation*}
$$

In case $\zeta \in \Omega_{m}$ and $-h \cdot \xi \leq 0,|h \cdot n| \leq\|h\| a_{m}(0, \pm \hat{n})$, so that $|h \cdot n|<\log 2$ if h satisfies (9.13). Thus for $h \in C$ with (9.13) satisfied, the domain $\Omega$ (9.11) is contained in $\Omega_{h}$.

COROLLARY 9.2. For functions $\mathrm{f} \in \tilde{\operatorname{Exp}}(\mathrm{a} ; \mathrm{T})$ the Newton series is valid, when $h \in C$ satisfies (9.13).

For example, when $a(z)$ is a constant $a$ on $\operatorname{pr}\left(\mathbb{R}^{n}+i C\right)$, the Newton series holds for $h \in C$ with $\|h\| \leq 10 g 2 / a$, if the function $f$ satisfies

$$
\forall k:|f(z)| \leq M_{k} \exp a_{m}\|z\|, \quad y \in C_{k},\|y\|>1 / k, a_{m}<a .
$$

This is a better bound than condition (9.12) since $\frac{1}{3} \pi \log 2\left\{\left(\frac{1}{3} \pi\right)^{2}+\log ^{2} 2\right\}^{-\frac{1}{2}}<$ < $\log 2$. This condition for $\|h\|$ generalizes the one dimensional case of NÖRLUND [13] p. 237.

In sections 7 and 10 of [14] we have seen that the bounds for $\|$ hll were determined by the value of the convex homogeneous function a on $C$ at the point $\tilde{h}$, namely $\|\mathrm{h}\| \leq \log 2 / a(\tilde{\mathrm{~h}})$ when $a(\tilde{\mathrm{~h}})>0$ or $\|\mathrm{h}\|$ arbitrarily large when $a(\tilde{h}) \leq 0$, where the function $f$ was of polynomial growth for $\|x\|$ large. Here the function f is of exponential growth also for $\|\mathrm{x}\|$ large and the bounds for $\|h\|$ are determined by the values of a on

$$
\left\{\beta \tilde{\mathrm{h}}+\mathrm{i} \alpha \tilde{\mathrm{~h}} \mid \alpha \geq 0, \beta \in \mathbb{R}, \alpha^{2}+\beta^{2}=1\right\},
$$

see conditions (9.12) and (9.13). This bound is always positive and finite, except in one case, where the Newton series is valid for $h \in C$ with $\|h\|$ arbitrarily large, namely for functions $f$ of exponential type, holomorphic in $\mathbb{R}^{\mathrm{n}}+\mathrm{i} \mathrm{C}$, satisfying

$$
\forall \varepsilon>0:|f(\beta \tilde{h}+i \alpha \tilde{h})| \leq M_{\varepsilon} \exp \varepsilon\left(\alpha^{2}+\beta^{2}\right)^{\frac{1}{2}}, \quad \alpha>0, \beta \in \mathbb{R} .
$$

This generalizes the case that $a(\tilde{h}) \leq 0$ in sections 7 and 10 of [14].
Finally we consider the case $\alpha<0$ more carefully and we will find that in that case too the Newton series (9.8) is valid for all y such that $y-\alpha h \in C$, even if $y$ does not belong to $C$. But first we have to modify the meaning of all the terms occurring in the series. We assume in the remaining of this section that for $\alpha<0$

$$
\operatorname{Re} s \geq-\alpha>m_{0},
$$

where $m_{0}$ is a non-negative integer. Firstly, we consider the series

$$
\begin{equation*}
\sum_{k=0}^{\infty}\binom{s}{k}\left(e^{-h \cdot \zeta}-1\right)^{k}=\sum_{k=0}^{\infty}\binom{s}{k}\left\{\sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} e^{-m h \cdot \zeta}\right\}=e^{-s h \cdot \zeta} \tag{9.14}
\end{equation*}
$$

for $\zeta \in \Omega_{h}$. We will show that we can rearrange some of the terms in this series. Therefore, we remark that the series

$$
\sum_{m=0}^{m_{0}} \sum_{k=m}^{\infty}\binom{s}{k}\binom{k}{m}(-1)^{k-m} \lambda_{m}
$$

is absolutely convergent for arbitrary numbers $\lambda_{m}, m=0,1, \ldots, m_{0}$, since by (5.1) [14]

$$
\left.\sum_{m=0}^{m_{0}} \sum_{k=N_{1}(s)}^{\infty}\left|\binom{s}{k}\right| \cdot\binom{k}{m} \cdot\left|\lambda_{m}\right|=\sum_{m=0}^{m_{0}}\left|\lambda_{m}\right| \cdot\left|\binom{s}{m}\right| \sum_{k=N_{1}}^{\infty}(s)-m\binom{s-m}{k} \right\rvert\, \leq
$$

$$
\leq \sum_{m=0}^{m_{0}} \frac{2}{|\Gamma(-s+m)|}\left|\lambda_{m}\right| \cdot\left|\binom{s}{m}\right| \sum_{k=1}^{\infty} k^{\alpha+m-1}<\infty,
$$

because $\alpha+m$ < 0 . Hence

$$
\begin{align*}
& \sum_{m=0}^{m_{0}} \sum_{k=m}^{\infty}\binom{s}{k}\binom{k}{m}(-1)^{k-m} \lambda_{m}=\sum_{k=m_{0}+1}^{\infty}\binom{s}{k} \sum_{m=0}^{m_{0}}\binom{k}{m}(-1)^{k-m} \lambda_{m}+  \tag{9.15}\\
& \quad+\sum_{k=0}^{m_{0}\binom{s}{k} \sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} \lambda_{m} \xlongequal{n o t} \Phi\left(s ; \lambda_{1}, \ldots, \lambda_{m_{0}}\right)}
\end{align*}
$$

exists. Now we write (9.14) in the following way

$$
\begin{aligned}
& \left.e^{-s h \cdot \zeta_{-\Phi}(s ; \lambda}, \ldots, \lambda_{m_{0}}\right)=\lim _{N \rightarrow \infty}\left\{\sum_{k=0}^{k}\binom{s}{k}\left[\sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} e^{-m h \cdot \zeta}\right]-\right. \\
& \left.-\sum_{k=m_{0}+1}^{N}\binom{s}{k}\left[\sum_{m=0}^{m}\binom{k}{m}(-1)^{k-m} \lambda_{m}\right]-\sum_{k=0}\binom{s}{k}\left[\sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} \lambda_{m}\right]\right\}= \\
& =\lim _{N \rightarrow \infty} \sum_{k=0}^{N}\left({ }_{k}^{s}\right)\left[\sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} \mu_{m}\right],
\end{aligned}
$$

where $\mu_{m}=\exp -m h \cdot \zeta$ for $m>m_{0}$ and $\mu_{m}=\exp (-m h \cdot \zeta)-\lambda_{m}$ for $m \leq m_{0}$. In order to compute $\Phi\left(s ; \lambda_{1}, \ldots, \lambda_{m_{0}}\right)$ we derive from (9.15) and from Re $s>m$

$$
\Phi\left(s ; \lambda, \ldots, \lambda_{m_{0}}\right)=\sum_{m=0}^{m_{0}} \lambda_{m}\binom{s}{m} \sum_{k=0}^{\infty}\binom{s-m}{k}(-1)^{k}=\sum_{m=0}^{m_{0}} \lambda_{m}\binom{s}{m}(1-1)^{s-m}=0
$$

for any numbers $\lambda_{m}$. Choosing $\lambda_{m}=\exp -\mathrm{mh} \cdot \zeta$ we obtain

$$
\begin{align*}
e^{-s h \cdot \zeta} & =\sum_{k=m_{0}+1}^{\infty}\binom{s}{k} \sum_{m=m_{0}+1}^{k}\binom{k}{m}(-1)^{k-m} e^{-m h \cdot \zeta}=  \tag{9.16}\\
& =\sum_{k=0}^{\infty}\binom{s}{k} \sum_{m=0}^{k}\binom{k}{m}(-1)^{k-m} \mu_{m}
\end{align*}
$$

where $\mu_{m}=\exp -m h \cdot \zeta$ for $m>m_{0}$ and $\mu_{m}$ arbitrary for $m \leq m_{0}$. In fact, we have rearranged the terms in (9.14) so, that first the summation is performed over all the terms with $\exp -m h \cdot \zeta$ for $m \leq m_{0}$ and it turns out that the series is independent of these terms.

Secondly, we give bounds to the functions

$$
\phi_{N}(\zeta)=\sum_{k=0}^{N}\binom{s}{k}\left(e^{-h \cdot \zeta}-1\right)^{k},
$$

when Re $s \geq-\alpha>m_{0}$. From p. 27 [14] we get for $\zeta \in \Omega_{h}(\varepsilon)$

$$
\left|\phi_{\mathrm{N}}(\zeta)\right| \leq 1+\mathrm{B}_{\mathrm{s}}\left\{1+(-1 \log \rho)^{-\alpha_{\alpha}}\right\}
$$

with $\rho=1-\frac{1}{2} \varepsilon \exp -\operatorname{Re} h \cdot \zeta$, whence

$$
\left|\phi_{N}(\zeta)\right| \leq 1+B_{s}+C(s, \alpha) \varepsilon^{-\alpha} e^{\alpha \operatorname{Reh} \cdot \zeta} .
$$

Therefore, we may conclude as in lemma 9.1 that the series $e^{-\alpha h \cdot \zeta} \phi_{N, z}(\zeta)$ converges in every space $A\left(H_{m}+0 ; C\right)$, when $y$ is such that $y-\alpha h \in C, h \in C$ and that for $\mu_{\zeta} \in A\left(H_{m}+0 ; C\right)$ ' the series

$$
\sum_{k=0}^{N}\binom{s}{k} e^{-\alpha h \cdot \zeta}\left(e^{-h \cdot \zeta}-1\right)^{k} \mu_{\zeta}
$$

converges strongly in $A\left(H_{m}+0 ; C\right)^{\prime}$.
Now using (9.16) we derive that for $\mathrm{f} \in \tilde{\operatorname{Exp}}\left(\mathrm{H}_{\Omega}, \mathrm{T}^{\mathrm{C}}\right)$ the following Newton series converges in the topology of $\operatorname{Exp} A\left({ }_{\Omega}, T^{C}\right)$

$$
\begin{aligned}
& f(z+i(s+\alpha) h)=\sum_{k=m_{0}+1}^{\infty}\binom{s}{k}<\sum_{m=m_{0}+1}^{k}\binom{k}{m}(-1)^{k-m} e^{-m h \cdot \zeta} e^{-\alpha h \cdot \zeta} \mu_{\zeta}, e^{i z \cdot \zeta}>= \\
& =\sum_{k=m_{0}+1}^{\infty}\binom{s}{k} \sum_{m=m_{0}+1}^{k}\binom{k}{m}(-1)^{k-m} f(z+i(m+\alpha) h) .
\end{aligned}
$$

Replacing $z+i \alpha h$ by $z$ and using the second part of (9.16) we find that the Newton series

$$
\begin{equation*}
f(z+i s h)=\sum_{k=0}^{\infty}\binom{s}{k} \Delta_{i h}^{k *} f(z), \tag{9.17}
\end{equation*}
$$

where the asterix means that in the points $\left\{z+i m h \mid m=0,1, \ldots, m_{0}\right\}$ where f is singular or undefined we may take zero instead of $\mathrm{f}(\mathrm{z}+\mathrm{imh})$, is valid for all $h \in C,\|h\| \leq b, \operatorname{Re} s \geq-\alpha>m_{0} \geq 0$ and all $y$ such that $y-\alpha h \in C$ and that it converges according to (9.9).

It may happen that $\mathrm{f} \in \operatorname{Exp}\left(\mathrm{H}_{\Omega}, \mathrm{T}^{\mathrm{C}}\right)$ can be continued analytically outside the domain $\mathbb{R}^{n}+i C$, so that $f(z+i m h)$ is defined for all $m$. But in fact, this is not essential and the series (9.17) has a meaning even if $f$ is singular or undefined in some points $z+i m h, m \leq m_{0}$, as long as $R e s>m_{0}$. Obviously, this is the generalization to several variables of the one dimensional case given in NÖRLUND [13] p. 237 in the first example 123.

We conclude with
THEOREM 9.1*. When Re $s \geq-\alpha>m_{0} \geq 0$, theorem 9.1 also holds for all y such that $\mathrm{y}-\alpha \mathrm{h} \in \mathrm{C}$; then the modified Newton series (9.17) converges according to (9.9).

## APPENDIX

## PASSAGE FROM LOCAL TO GLOBAL RELATIONS

In this appendix we discuss some well-known properties in the theory of functions of several complex variables. Except section $I$ all sections are devoted to the problem how to extend local relations between holomorphic functions to global relations. As some readers may not be familiar with the topics used to solve this problem, we will go more into detail than merely copying definitions and theorems from litterature. We give those proofs that show how to use the various concepts (as sheaves and cohomology) in deriving the main result. In fact, since we want a quantitative result in section 7 , we perform the same steps there as in section IV of this appendix, then taking care of estimates. Therefore, we also give the quantitative theorems these steps start from. Almost the same method HORMANDER [7] uses in his book is applied here and we repeatedly refer to this book.

## I. DOMAINS OF HOLOMORPHY

In this section we give some definitions and theorems which are used in section 2 , the case of holomorphic functions on compact sets.

Let $\Omega$ be an open set in $\mathbb{C}^{n}$. We denote by $\mathrm{A}(\Omega)$ the space of all holomorphic functions in $\Omega$ with the topology of uniform convergence on compact subsets $K$ of $\Omega$. All functions holomorphic in a certain domain $\Omega$ in $\mathbb{C}^{n}, n \geq 2$, might be continued analytically into a larger domain. Domains for which this is not possible are called domains of holomorphy. Thus $\Omega$ is a domain of holomorphy if and only if there exists a function $f \in A(\Omega)$ which cannot be continued analytically beyond $\Omega$, that is, it is not possible to find $\Omega_{1}$ and $\Omega_{2}$, with $\emptyset \neq \Omega_{1} \subset \Omega_{2} \cap \Omega$ and with $\Omega_{2}$ connected and not contained in $\Omega$, and $\mathrm{f}_{1} \in \mathrm{~A}\left(\Omega_{2}\right)$ so that $\mathrm{f}=\mathrm{f}_{1}$ in $\Omega_{1}$. One can decide whether a domain $\Omega$ is a domain of holomorphy by other means too. We will discuss some of these means which are most useful in applications.

For a compact set K of an open set $\Omega$ we define the $\mathrm{A}(\Omega)$-hull $\hat{\mathrm{K}}_{\Omega}$ of K by

$$
\begin{equation*}
\hat{K}_{\Omega}=\left\{z\left|z \in \Omega,|f(z)| \leq \sup _{z \in K}\right| f(z) \mid \text { for all } f \in A(\Omega)\right\} . \tag{Al}
\end{equation*}
$$

If we choose $f(z)=\exp z \cdot \zeta$ we find that $\hat{K}_{\Omega}$ is contained in the convex hull $\mathrm{ch}(\mathrm{K})$ of K . Domains of holomorphy can be characterized by the following theorem, th.2.5.5(ii) of [7]:

THEOREM A1. $\Omega$ is a domain of holomorphy if and only if from $\mathrm{K} \subset \Omega$ it follows that $\hat{\mathrm{K}}_{\Omega} \subset \Omega$.

Hence convex open sets in $\mathbb{C}^{\mathfrak{n}}$ are domains of holomorphy. Conversely Bochner's theorem (th.2.5.12 of [7] or 17.5 of [16]) yields:

THEOREM A2. A tube domain $\mathbb{R}^{\mathrm{n}}+\mathrm{i} 0$, where 0 is a domain in $\mathbb{R}^{\mathrm{n}}$, is a domain of holomorphy if and only if 0 is convex.

A more geometrical characterization of domains of holomorphy is obtained by regarding them as pseudoconvex sets. These sets can be defined with the aid of plurisubharmonic functions. Rather than giving a precise definition (2.6.1 of [7]) we state some results. As in (A1) one can define a $\mathrm{P}(\Omega)$-hull $\hat{K}_{\Omega}^{P}$ of $K$ by requiring that $f$ is plurisubharmonic instead of $f \in A(\Omega)$. Then like theorem Al an open set $\Omega$ is pseudoconvex if and only if from $K \subset \Omega$ it follows that $\hat{K}_{\Omega}^{P} c \subset \Omega$. Since the function $|f(z)|$ is plurisubharmonic if $f$ is holomorphic, domains of holomorphy are pseudoconvex. The converse is also true (th.4.2.8 of [7]):

THEOREM A3. An open pseudoconvex set is a domain of hoZomorphy.
Actually, if K is a compact set of an open pseudoconvex set $\Omega$, then $\hat{\mathrm{K}}_{\Omega}$ equals $\hat{\mathrm{K}}_{\Omega}^{\mathrm{P}}$ (th. 3.4.3 of [7]). Therefore, we will not distinguish between the concepts of pseudoconvex open set and of domains of holomorphy and we assume $\Omega$ to be one or the other where necessary.

THEOREM A5. Let $\Omega$ be a pseudoconvex open set in $\mathbb{C}^{\mathrm{n}}$ and K a compact subset of $\Omega$, such that $\hat{K}_{\Omega}=K$. Every function analytic in a neighborhood of K can then be approximated uniformly on K by functions in $\mathrm{A}(\Omega)$.

This is theorem 4.3.2 of [7].
DEFINITION A5. A domain of holomorphy $\Omega \subset \mathbb{C}^{\mathrm{n}}$ is called a Runge domain if polynomials are dense in $A(\Omega)$, that is if every $f \in A(\Omega)$ can be uniformly
approximated on an arbitrary compact set in $\Omega$ by analytic polynomials.

Since polynomials are dense in $A\left(\mathbb{C}^{n}\right)$ we might as well have considered arbitrary entire functions instead of polynomials in definition A5. For a compact set K we define

$$
\tilde{K}=\left\{z\left|z \in \mathbb{C}^{n},|P(z)| \leq \sup _{z \in K}\right| P(z) \mid \text { for all polynomials } P\right\} .
$$

Then $\tilde{K}=\hat{K}_{\mathbb{C}^{n}}$ and compact sets $K$ with $K=\widetilde{K}$ are called polynomially convex. However, we even have (th.2.7.3 of [7]):

THEOREM A6. $\Omega$ is a Runge domain if and only if for every compact set $K$ in $\Omega$ $\widetilde{\mathrm{K}}=\hat{\mathrm{K}}_{\Omega}$ 。

This theorem is a special case (namely when $\Omega_{1}=\Omega$ and $\Omega_{2}=\mathbb{C}^{n}$ ) of the following theorem (th.4.3.3 of [7]):

THEOREM A7. Let $\Omega_{1} \subset \Omega_{2}$ be domains of holomorphy. Then every function in $A\left(\Omega_{1}\right)$ can be approximated by functions in $A\left(\Omega_{2}\right)$ uniformly on every compact subset of $\Omega_{1}$ if and only if for every compact subset $K$ of $\Omega_{1}$ we have $\hat{\mathrm{K}}_{\Omega_{2}}=\hat{\mathrm{K}}_{\Omega_{1}}$.
II. THE $\bar{\partial}$-OPERATOR

In this section we define the $\bar{\partial}$-operator and give some existence theorems.

Let $u$ be a complex valued differentiable function in $\Omega \subset \mathbb{C}^{n}$. We denote $z=x+i y \in \Omega$ also as $z=(y, x)$ with $x \in \mathbb{R}^{n}$ and $y \in \mathbb{R}^{n}$, where now $\Omega$ is regarded as an open set in $\mathbb{R}^{2 n}$ (the reason for not writing $z=(x, y)$ is, that, when we do so for $\zeta=\xi+i n \in \mathbb{C}^{n}, \zeta=(\xi, \eta)$, then $-\operatorname{Re}(i z \cdot \zeta)$ can be written as inproduct between the vectors $(y, x)$ and $(\xi, \eta)$ in $\mathbb{R}^{2 n}$ ). The components of $z$ are denoted by $z_{j}=x_{j}+i y_{j_{-}}$and $\bar{z}_{j}=x_{j}-i y_{j}$. When differentiation takes place, we rather use $z_{j}$ and $\bar{z}_{j}, j=1, \ldots, n$, as coordinates than $(y, x)$, so that

$$
\frac{\partial}{\partial z_{j}}=\frac{1}{2}\left(\frac{\partial}{\partial x_{j}}-i \frac{\partial}{\partial y_{j}}\right) \text { and } \frac{\partial}{\partial \bar{z}_{j}}=\frac{1}{2}\left(\frac{\partial}{\partial x_{j}}+i \frac{\partial}{\partial y_{j}}\right)
$$

Then we get

$$
d u=\sum_{j=1}^{n} \frac{\partial u}{\partial z_{j}} d z_{j}+\sum_{j=1}^{n} \frac{\partial u}{\partial \bar{z}_{j}} d \bar{z}_{j}
$$

and with

$$
\partial u=\sum_{j=1}^{n} \frac{\partial u}{\partial z_{j}} d z_{j} \quad \text { and } \quad \bar{\partial} u=\sum_{j=1}^{n} \frac{\partial u}{\partial \bar{z}_{j}} d \bar{z}_{j}
$$

we may also write

$$
\mathrm{du}=\partial \mathrm{u}+\bar{\partial} \mathrm{u}
$$

When we write $\bar{\partial} u=0$ in $\Omega$, we mean that every component $\bar{\partial}_{j} u=\partial u / \partial \bar{z}_{j}$ must vanish in $\Omega$. These are exactly the Cauchy Riemann equations, so that we get THEOREM A8. A function $u$ in $C^{1}(\Omega)$ is holomorphic in the open set $\Omega$ if and only if $\bar{\partial} \mathrm{u}=0$ in $\Omega$.

In the above $\bar{\partial} \mathrm{u}$ is a $(0,1)$-form. We call g a $(0,1)$-form in $\Omega$ if it can be written as

$$
g(z)=\sum_{k=1}^{n} g_{k}(z) d \bar{z}_{k}, \quad z \in \Omega
$$

where $g_{k}, k=1, \ldots, n$, are functions in $\Omega$. We will give a condition when a $(0,1)$-form $g$ can be written as $\bar{\partial} u$ for some function $u$. A necessary condition on $g$ is $\bar{\partial} g=0$, where we define

$$
\overline{\partial g}=\sum_{m=1}^{n} \sum_{k=1}^{n} \frac{\partial g_{k}}{\partial \bar{z}_{m}} d \bar{z}_{m} \wedge \mathrm{~d} \bar{z}_{\mathrm{k}}
$$

when the functions $g_{k}$ are differentiable. Here we may use the rule

$$
\begin{equation*}
\mathrm{d} \bar{z}_{\mathrm{m}} \wedge \mathrm{~d} \overline{\mathrm{z}}_{\mathrm{k}}=-\mathrm{d} \overline{\mathrm{z}}_{\mathrm{k}} \wedge \mathrm{~d} \overline{\mathrm{z}}_{\mathrm{m}}, \quad \mathrm{k}, \mathrm{~m}=1, \ldots, \mathrm{n} \tag{A2}
\end{equation*}
$$

and $\bar{\partial} g=0$ if the coefficients of all the $d \bar{z}_{\mathrm{m}} \wedge \mathrm{d} \bar{z}_{\mathrm{k}}(\mathrm{m}<\mathrm{k})$ vanish. It is easy to see that for any $u \in C^{2}(\Omega) \bar{\partial} \bar{\partial} u=0$, so that indeed $\bar{\partial} g=0$ is a
necessary condition. When $\Omega$ is pseudoconvex, it is also a sufficient condition. This we state in a theorem, which we give in a more general form, namely for ( $0, q$ )-forms. We say that $g$ is a ( $0, q$ )-form in $\Omega(q=0,1, \ldots, n$ ) if it can be written in the form

$$
g=\sum_{|I|=q} g_{I}(z) d \bar{z}^{-I}, \quad z \in \Omega
$$

where $\mathrm{I}=\left(\mathrm{k}_{1}, \ldots, \mathrm{k}_{\mathrm{q}}\right)$ is a multiindex and $\mathrm{d} \overline{\mathrm{z}}^{\mathrm{I}}=\mathrm{d} \overline{\mathrm{z}}_{\mathrm{k}_{1}} \wedge \ldots \wedge \overline{\mathrm{z}}_{\mathrm{k}_{\mathrm{q}}}$ and where the summatition is performed over all multiindices $I$ with $k_{1}<\mathrm{k}_{2}<\ldots<\mathrm{k}_{\mathrm{q}}$ (for again we may use the rule (A2)). Thus $g$ has $(\underset{q}{n})$ coefficients $g_{I}$. We define

$$
\bar{\partial} g=\sum_{k=1}^{n} \sum_{|I|=q} \frac{\partial g_{I}}{\partial \bar{z}_{k}} d \bar{z}_{k} \wedge d \bar{z}^{I}
$$

where (A2) should be used. It is clear that $\bar{\partial} \bar{\partial} g=0$. Now the following existence theorem for the $\bar{\partial}$-operator holds (cor.4.2.6 of [7]):

THEOREM A9. Let the coefficients $\mathrm{g}_{\mathrm{I}}$ of the $(0, \mathrm{q}+1)$-form g in the pseudoconvex open set $\Omega$ be $C^{\infty}$-functions and let $\bar{\partial} g=0$. Then there exists a $(0, q)$-form $u$ with $C^{\infty}$-coefficients in $\Omega$ such that $\bar{\partial} u=g$.

Next we state a similar theorem, where besides the existence of $u$ also estimates of $u$ in terms of estimates for $g$ are given. We use the measure $e^{-\phi} d \lambda$, where $d \lambda$ is the Lebesgue measure in $\mathbb{C}^{n}$, and $\phi(y, x)$ is a plurisubharmonic function. We do not give the definition of a plurisubharmonic function (see 2.6 .1 of [7]), but we merely state that a convex function $\tilde{a}(y, x)$ is plurisubharmonic, that $\log \left(1+\|z\|^{2}\right)$ is plurisubharmonic and that $\alpha \phi+\beta \psi$ is plurisubharmonic for $\alpha \geq 0, \beta \geq 0$ whenever $\phi$ and $\psi$ are plurisubharmonic. These facts will be sufficient for the applications we make. For a ( $0, q$ )-form f in $\Omega\left(\mathrm{q}=0,1, \ldots, \mathrm{n}\right.$ ), where the coefficients $\mathrm{f}_{\mathrm{I}}$ are locally square integrable functions, we write

$$
|f(z)|^{2}=\sum_{|I|=q}\left|f_{I}(z)\right|^{2}, \quad z \in \Omega
$$

and

$$
\|f\|_{\phi}=\int_{\Omega}|f(z)|^{2} e^{-\phi(z)} d \lambda(z)
$$

We remark that for such an $f$ we must take the weak derivative in $\bar{\partial} f$, thus derivatives in distributional sense. Then we state the following theorem (th.4.4.2 of [7]):

THEOREM A10. Let $\Omega$ be a pseudoconvex open set in $\mathbb{C}^{\mathrm{n}}$ and $\phi$ any plurisubharmonic function in $\Omega$. For every $(0, q+1)$-form $g$ with locally square integrable coefficients, with $\|g\|_{\phi}$ finite and with $\overline{\partial g}=0$, there is a $(0, q)-$ form u in $\Omega$ with locally square integrable coefficients, such that $\bar{\partial} \mathrm{u}=\mathrm{g}$ and

$$
\int_{\Omega}|u(z)|^{2} e^{-\phi(z)}\left(1+\|z\|^{2}\right)^{-2} d \lambda(z) \leq\|g\|_{\phi}^{2} .
$$

Here $u$ depends on $\phi$, when the right hand side is finite for more than one function $\phi$.

## III. ANALYTIC SHEAVES

In this section we discuss some properties of analytic sheaves and we formulate the main problem of this appendix. We do not give a general definition of a sheaf on an open set $\Omega$ in $\mathbb{C}^{n}$, but we just give the properties we need in this paper. A more complete description can be found in [6] or [7].

For $z \in \Omega$ we denote by $A_{z}$ the set of equivalence classes of functions $f$ which are analytic in a neighborhood of $z$, under the equivalence relation $\mathrm{f} \sim \mathrm{g}$ if $\mathrm{f}=\mathrm{g}$ in a neighborhood of z in $\Omega$. The residue class $\mathrm{f}_{\mathrm{z}}$ of f in $A_{z}$ is called the germ of $f$ at $z .{ }^{1)}$ It is clear that $A_{z}$ is a ring. Let

$$
A=\bigcup_{z \in \Omega} A_{z}
$$

1) 

Since an analytic function is determined completely when it is given in an open set, the residue class of $f$ is trivial: it consists of $f$ only. But when we consider the restriction of $f$ to a variety V in $\Omega$, we get a sheaf on $V$ ( $V$ is a simultaneous zero set of holomorphic functions in $\Omega$ ) and the equivalence classes are no longer trivial, see [6] def.IV D.5, p. 143 and see also section VI. Also, when we consider $C^{\infty}$-functions instead of analytic functions, it has sense to define the germ of $f$ at $z$ as a residue class.
where the rings $A_{z}$ are considered as disjoint sets. Furthermore, let the collection of subsets of $A$ of the form

$$
\left\{f_{z} \mid z \in \omega \subset \Omega, \text { where } \omega \text { is open and } f \in A(\omega)\right\}
$$

where $\omega$ runs over the collection of open subsets of $\Omega$ and f runs over the elements of $A(\omega)$, be a basis for the topology of $A$. Then for every open subset $\omega$ of $\Omega$ and every $f \in A(\omega)$ the map $\phi$ from $\omega$ into $A$ with $\phi(z)=f_{z}$ is open and continuous.

Let $\pi$ be the map from $A$ into $\Omega$ which maps $A_{z}$ onto $z$. Then $\pi \phi=$ identity. In general we call the image of a subset $U$ of $\Omega$ under a continuous map $\phi$ : $U \rightarrow A$, with $\pi \phi=$ identity, or the map $\phi$ itself, a section of A over $U$. The set of all sections of $A$ over $U$ is denoted by $\Gamma(U, A)$. In fact an element of $\Gamma(U, A)$ is the restriction to $U$ of a holomorphic function in a neighborhood of $U$ in $\Omega$ or if $U$ itself is open, it is a holomorphic function in $U$.

The space $A$ is an example of a sheaf on $\Omega$. Since $A_{z}$ is a ring for each $z \in \Omega$, we can consider a sheaf $F$ such that $F_{z}$ is an $A_{z}$-module for each $z \in \Omega$ and such that the product of a section in $A$ and a section in $F$ is a section in $F$. Such a sheaf is called an analytic sheaf. In particular we will consider ideals in $A_{z}$ and modules in $A_{z}^{p}$. Since the ring $A_{z}$ is a noetherian ring ([7] th.6.3.3 or [6] th.II.B.9) the ideals in $A_{z}$ and the modules in $A_{z}^{p}$ are finitely generated.

For example, let $U$ be an open subset of $\Omega$ with $\emptyset \neq U \neq \Omega$ and let an analytic sheaf $F$ be given by $F_{z}=A_{z}$ if $z \in U$ and $F_{z}=0$ if $z \in \Omega \backslash U$. A section of this sheaf over a connected open set intersecting $\Omega \backslash U$ must be zero by the uniqueness of analytic continuation. In any point $z \in \Omega, F_{z}$ is finitely generated, but in any neighborhood $\omega$ of a boundary point of $U$ in $\Omega$ $F$ is not finitely generated by the sections over $\omega$.

Thus although $F_{z}$ is finitely generated in any point $z \in \Omega$, we cannot always use the same generators for all $z$ in a neighborhood of any point. However, we consider sheafs where this property indeed is satisfied. Namely, an analytic sheaf $F$ is said to be locally finitely generated if for every given point in $\Omega$ there exists a neighborhood $\omega$ in $\Omega$ and a finite number of sections $f_{1}, \ldots, f_{q} \in \Gamma(\omega, F)$ so that $F_{z}$ is generated by $\left(f_{1}\right)_{z}, \ldots,\left(f_{q}\right)_{z}$ as an $A_{z}$-module for every $z \in \omega$. In particular we will consider locally finitely
generated subsheaves $F$ of $A^{p}$, so that then in the above definition for $k=1, \ldots, q$ $f_{k}$ is a $p$-tuple of analytic functions $f_{k}^{j} \in A(\omega)$ in $\omega$, $j=1, \ldots, p$ with $\left(f_{k}\right)_{z}=\left(f_{k}^{1}(z), \ldots, f_{k}^{p}(z)\right)$.

Let $F$ be a locally finitely generated analytic sheaf, let $f_{1}, \ldots f_{q}$ be sections over an open set $U$ of $\Omega$ and let for any $z \in U$

$$
R_{z}\left(f_{1}, \ldots, f_{q}\right)=\left\{\left(g^{1}, \ldots, g^{q}\right) \in A_{z}^{q} \mid \sum_{k=1}^{q} g^{k}\left(f_{k}\right)_{z}=0\right\}
$$

$R_{z}$ is a submodule of $A_{z}^{q}$, called the module of relations between $f_{1}, \ldots, f_{q}$ at $z$. Then

$$
R\left(f_{1}, \ldots, f_{q}\right)=\bigcup_{z \in U} R_{z}\left(f_{1}, \ldots, f_{q}\right)
$$

is a subsheaf of $A^{q}$ on $U$, called the sheaf of relations between $f_{1}, \ldots, f_{q}$. A locally finitely generated analytic sheaf $F$ is called a coherent analytic sheaf, if $R\left(f_{1}, \ldots, f_{q}\right)$ is locally finitely generated for all $\mathrm{U} \subset \Omega$, all $f_{k} \in \Gamma(U, F), k=1, \ldots, q$ and $a l l q$. When $F$ is a locally finitely generated subsheaf of $A^{p}$, the last condition is always satisfied. For by Oka's theorem ([7] th.6.4.1 and th.7.1.5 or [6] th.IV.C.1 and IV.B.7 and 8) every locally finitely generated subsheaf $F$ of $A^{p}$ is coherent; that is for any point in $U \subset \Omega$ one can find a neighborhood $\omega \subset U$ and finitely many elements $G_{1}, \ldots, G_{r} \in \Gamma\left(\omega, R\left(f_{1}, \ldots, f_{q}\right)\right)$ (thus for $\ell=1, \ldots, r \quad G_{\ell}=\left(g_{\ell}^{1},,,, g_{\ell}^{q}\right) \in$ $\epsilon A(\omega)^{q}$ and for $z \in \omega$

$$
\begin{equation*}
\sum_{\mathrm{k}=1}^{\mathrm{q}} \mathrm{~g}_{\ell}^{\mathrm{k}}(\mathrm{z}) \mathrm{f}_{\mathrm{k}}^{\mathrm{j}}(\mathrm{z})=0, \quad \mathrm{j}=1, \ldots, \mathrm{p}, \tag{A3}
\end{equation*}
$$

$\ell=1, \ldots, r)$, so that $R_{z}$ for every $z \in \omega$ is equal to the $A_{z}$-module generated by $\left(G_{1}\right)_{z}, \ldots,\left(G_{r}\right)_{z}$.

If two of the sheafs of the exact sequence (that is the image of one map is the kernel of the next map)

$$
0 \rightarrow F \rightarrow G \rightarrow H \rightarrow 0
$$

are coherent, then the third sheaf is coherent too, see [6] th.IV.B.13.

From now on we only consider coherent analytic sheaves and we do not state this all the time.

Let now in (A3) the functions $f_{k}^{j} \in A(U)$ be polynomials $P_{j k}=f_{k}^{j}$, $j=1, \ldots, p, k=1, \ldots, q$ and let $U=\Omega$. Then we consider the sheaf homomorphism

$$
P: A^{q} \rightarrow A^{p}
$$

defined by mapping $\left(g^{1}(z), \ldots, g^{q}(z)\right) \in A_{z}^{q}$ to

$$
\left(\sum_{k=1}^{q} P_{1 k}(z) g^{k}(z), \ldots, \sum_{k=1}^{q} P_{p k}(z) g^{k}(z)\right) \in A_{z}^{p}, \quad z \in \Omega
$$

We have seen that the image and the kernel of this map are coherent analytic sheaves and in particular it follows from the proof of the Oka theorem (th.6.4.1 of [7]), that the functions $g_{p}^{k}$ in (A3), $\ell=1, \ldots, r, k=1, \ldots, q$, can be chosen to be polynomials. Thus the kernel $R_{P}$ of this map is generated by the germs of all $q$-tuples $Q=\left(Q_{1}, \ldots, Q_{q}\right)$ with $Q_{k}$ polynomials for $\mathrm{k}=1, \ldots, \mathrm{q}$, such that

$$
\begin{equation*}
\sum_{k=1}^{q} P_{j k}(z) Q_{k}(z)=0, \quad z \in \Omega \tag{A4}
\end{equation*}
$$

Furthermore, since the polynomial ring (over $\mathbb{C}$ ) is noetherian, the module of all $Q=\left(Q_{1}, \ldots, Q_{q}\right)$ with $Q_{q}$ polynomials satisfying (A4) is finitely generated over the polynomial ring. Thus since all the generators, that is in a neighborhood of all the points of $\Omega$, are polynomial q-tuples $Q, R_{P}$ is generated by a finite number of such $Q$, say by $Q_{\ell}=\left(Q_{1 \ell}, \ldots, Q_{1 \ell}\right)$, $\ell=1, \ldots, r$, where $Q_{k \ell}$ is a polynomial. Summarizing we get the exact sequences of sheaf homomorphisms

$$
\begin{equation*}
0 \rightarrow R_{P} \rightarrow A^{q} \xrightarrow{P} F \rightarrow 0, \tag{A5}
\end{equation*}
$$

where $F$ is the image of $P$ and

$$
\begin{equation*}
0 \rightarrow R_{Q} \rightarrow A^{r} \xrightarrow{Q} R_{P} \rightarrow 0 . \tag{A6}
\end{equation*}
$$

A section $f=\left(f_{1}, \ldots, f_{p}\right)$ in $F$ is a p-tuple holomorphic functions in $\Omega$, thus $f_{j} \in A(\Omega), j=1, \ldots, p$, satisfying locally, that is in a neighborhood $\omega_{s}$ of each point in $\Omega$,

$$
\begin{equation*}
f_{j}(z)=\sum_{k=1}^{q} P_{j k}(z) g_{k}^{s}(z), \quad z \in \omega_{s}, g_{k}^{s} \in A\left(\omega_{s}\right), j=1, \ldots, p \tag{A7}
\end{equation*}
$$

In $\omega_{s} \cap \omega_{t}$ the functions $g_{k}^{s} \in A\left(\omega_{s}\right)$ are not necessarily equal to the functions $g_{k}^{t} \in A\left(\omega_{t}\right)$, for they may differ by a section $h^{s t}$ in $R_{p}$ over $\omega_{s} \cap \omega_{t}$. We would like that $h^{s t}=0$, thus that (A7) holds globally, that is we would like to find $g_{k} \in A(\Omega)$ such that (A7) holds for all $z \in \Omega$. The main problem of this appendix is to prove that such functions $g_{k}, k=1, \ldots, q$, exist. We can formulate this as: the problem is to prove that the following sequence of sections is exact

$$
\begin{equation*}
0 \rightarrow \Gamma\left(\Omega, R_{\mathrm{P}}\right) \longrightarrow \Gamma\left(\Omega, A^{\mathrm{q}}\right) \xrightarrow{\mathrm{P}} \Gamma(\Omega, F) \longrightarrow 0 \tag{A8}
\end{equation*}
$$

That the sequence is exact in the first two places is clear, but our attention is paid to the exactness in the last place, thus to prove that the map $P$ is surjective. We will find that (A8) is indeed exact, when $\Omega$ is pseudoconvex. Then starting with (A6) we would at the same time have solved the problem:

THEOREM A11. If the function $\mathrm{f}_{\mathrm{k}} \in \mathrm{A}(\Omega)$ satisfy

$$
\sum_{k=1}^{q} P_{j k}(z) f_{k}(z)=0, \quad z \in \Omega, j=1, \ldots, p
$$

then there are functions $g_{l} \in A(\Omega), \ell=1, \ldots, r$, such that

$$
f_{k}(z)=\sum_{\ell=1}^{r} Q_{k \ell}(z) g_{\ell}(z), \quad z \in \Omega, k=1, \ldots, q,
$$

when $\Omega$ is pseudoconvex.

In this section we define cohomology groups and show how they are used to solve the problem formulated in section III.

We consider the sheaf $F$ as an additive commutative group. Let $U=\left\{I_{i}\right\}_{i \in I}$ be an open covering of the open set $\Omega$ in $\mathbb{C}^{n}$. If $p$ is a nonnegative integer, we denote by $s=\left(s_{o}, \ldots, s_{p}\right)$ any element in $I^{p+1}$ and we set $U_{s}=U_{S_{0}} \cap \ldots \cap U_{S_{p}}$. A map assigning to every $s \in I^{P^{+1}}$ a section $c_{s} \in \Gamma\left(U_{s}, F\right)$ so that $c_{s}$ is an alternating function of $s$ (that is, $c_{s}$ changes sign if two indices in $s$ are permuted) is called a p-cochain of the covering $U$ with values in $F$. Here we define $\Gamma(\emptyset, F)=0$, the abelian group with one element. Then the set $C^{P}(U, F)$ of all $p$-cochains is an abelian group.

A map $\delta$ from $C^{p}(U, F)$ into $C^{p+1}(U, F)$, called the coboundary operator, is defined as follows: if $c \in C^{p}(U, F)$, then for $s \in I^{p+2}$

$$
(\delta c)_{s}=\sum_{j=0}^{p+1}(-1)^{j} c_{s_{0}} \ldots \hat{s}_{j} \ldots s_{p+1}
$$

where the notation $\hat{s}_{j}$ means that the index $\mathbf{s}_{\mathbf{j}}$ should be removed. We introduce the group of p -cocycles

$$
z^{\mathrm{P}}(U, F)=\left\{c \mid c \in C^{\mathrm{P}}(U, F), \delta c=0\right\}
$$

and the group of $p$-coboundaries

$$
{ }_{B}{ }^{\mathrm{p}}(U, F)=\left\{\delta c \mid c \in C^{\mathrm{p}-1}(U, F)\right\}
$$

where $C^{-1}=0$. Since for all $c \in C^{p}(U, F) \delta \delta c=0, B^{p}$ is a subgroup of $Z^{p}$. We can, therefore, define the quotient group

$$
{ }_{\mathrm{H}}{ }^{\mathrm{p}}(U, F)=\mathrm{z}^{\mathrm{p}}(U, F) /{ }_{\mathrm{B}^{\mathrm{P}}(U, F)} \text {, }
$$

which is called the $p^{\text {th }}$ cohomology group of $U$ with values in $F$.
For example, if $c$ is a 0 -cocycle, then $c_{s_{0}}{ }^{-c} s_{s_{1}}=0$ in $U_{s_{0}} \cap U_{s_{1}}$ for all
$s_{0}$ and $s_{1}$, which means that there is a section $f \in \Gamma(\Omega, F)$ with the restriction $c_{s}$ to $U_{s}$ for every s. Hence

$$
\begin{equation*}
H^{0}(U, F) \cong \Gamma(\Omega, F) \tag{A9}
\end{equation*}
$$

Let $V=\left\{V_{j}\right\} j_{j \in J}$ be another covering of $\Omega$, which is a refinement of $U$. This means that there is a map $\rho$ from $J$ into $I$ such that $V_{j} \subset U_{\rho(j)}$ for every $j \in J$. If $c \in C^{P}(U, F)$, we can then define a cochain $\rho c \in C^{P}(V, F)$ by setting ( $\rho c)_{s}$ equal to the restriction of $c_{\rho(s)} \ldots \ldots\left(s_{p}\right)$ to $v_{s}$. One easily sees that $\rho$ commutes with the coboundary operators in $C^{p}(U, F)$ and $C^{p}(V, F)$ and, therefore, it induces a map $\rho^{*}$ from $H^{p}(U, F)$ into $H^{p}(V, F)$. This map $\rho^{*}$ is independent of the choice of $\rho$ (see prop.7.3.1 of [7]).

Let $E$ be the sheaf of germs of $C^{\infty}$-functions on $\Omega$ (see the footnote on page 85 ).

THEOREM A12. Let $F$ be a sheaf of $E$-modules on $\Omega$, then $H^{p}(U, F)=0$ for $p \geq 1$ and every covering $U$ of $\Omega$.

PROOF. Let $\phi_{\nu}$ be a partition of unity subordinate to the covering $U$, that is i) $\phi_{V}$ is a $C^{\infty}$-function with compact support in $U_{i_{V}}$ for a certain index $i_{V}$; ii) all but a finite number of functions $\phi_{V}$ vanish identically on any compact subset of $\Omega$;
iii) $\sum_{V} \phi_{V}=1$ on $\Omega$.

For $c \in Z^{p}(U, F)$ we put, when $s \in I^{p}$,

$$
g_{s}=\sum_{V} \phi_{V} c_{i_{v} s},
$$

which defines a cochain $g$ in $C^{p-1}(U, F)$. Since with $s \in I^{p+1}$

$$
(\delta c)_{i_{v} s}=c_{s}+\sum_{j=0}^{p+1}(-1)^{j+1} c_{i_{\nu} s_{0} \ldots \hat{s}_{j} \ldots s_{p}=0}=0
$$

we get

Thus $c$ is a coboundary.

Let $F, G, H$ be three sheaves of abelian groups on $\Omega$ and let $\phi$ and $\psi$ be sheaf homomorphisms such that the sequence

$$
0 \rightarrow F \xrightarrow{\phi} G \xrightarrow{\psi} H \rightarrow 0
$$

is exact (thus $\phi$ is injective, $\psi$ is surjective, the kernel of $\psi$ is the image of $\phi$ ). This defines exact sequences between sections; thus we get the exact sequence

$$
0 \rightarrow C^{p}(U, F) \rightarrow C^{p}(U, G) \rightarrow C^{p}(U, H)
$$

but the last map is not necessarily surjective. We denote its image by $C_{a}^{p}(U, H)$ and call it the group of liftable cochains. We then have an exact sequence

$$
0 \rightarrow C^{p}(U, F) \rightarrow C^{p}(U, G) \rightarrow C_{a}^{p}(U, H) \rightarrow 0
$$

Since $\delta$ commutes with $\psi, C_{a}^{p}(U, H)$ is mapped by $\delta$ into $C_{a}^{p^{+1}}(U, H)$ and we can define the cohomology groups

$$
\mathrm{H}_{\mathrm{a}}^{\mathrm{p}}(U, H)=\mathrm{z}_{\mathrm{a}}^{\mathrm{p}} /_{\mathrm{B}_{\mathrm{a}}^{\mathrm{p}}}
$$

where $Z_{a}^{p}\left(B_{a}^{p}\right)$ is the group of all liftable p-cocycles (coboundaries of liftable ( $p-1$ )-cochains). Then we have the commutative diagram with exact columns:


Now we construct a map $\delta^{*}$ from $\mathrm{H}_{\mathrm{a}}^{\mathrm{p}}(U, H)$ into $\mathrm{H}^{\mathrm{p}+1}(U, F)$ as follows: If $\mathrm{f} \in \mathrm{Z}_{\mathrm{a}}^{\mathrm{P}}(U, H)$, then $\mathrm{f}=\psi \mathrm{g}$ for some $\mathrm{g} \in \mathrm{C}^{\mathrm{P}}(U, G)$ and $\psi \delta \mathrm{g}=\delta \psi \mathrm{g}=\delta \mathrm{f}=0$, hence $\delta g=\phi c$ for some $c \in C^{p+1}(U, F)$. We put $\delta^{*} f=c$ and we have $\phi \delta c=\delta \phi c=$ $=\delta \delta g=0$, hence $\delta c=0$, that is $c \in Z^{p+1}(U, F)$, since $\phi$ is injective.

Another representative of $f$ in $H_{a}^{p}(U, H)$ differs from $f$ by a coboundary $f_{1} \in B_{a}^{P}(U, H)$. Then $f_{1}=\psi g_{1}$ for some $g_{1} \in C^{P}(U, G)$ and also $f_{1}=\delta f^{\prime}$ for some $f^{\prime} \in C_{a}^{p-1}(U, H)$. Furthermore there is a $g^{\prime} \in C^{p-1}(U, G)$ with $\psi g^{\prime}=f^{\prime}$ and we have $\psi\left(\mathrm{g}_{1}-\delta \mathrm{g}^{\prime}\right)=\psi \mathrm{g}_{1}-\delta \psi \mathrm{g}^{\prime}=\mathrm{f}_{1}-\delta \mathrm{f}^{\prime}=0$, thus $\mathrm{g}_{1}-\delta \mathrm{g}^{\prime}=\phi \mathrm{c}^{\prime}$ for some $\mathrm{c}^{\prime} \epsilon$ $\in C^{\mathrm{P}}(U, F)$. Let $c_{1}=\delta c^{\prime} \in B^{p+1}(U, F)$, then $\delta g_{1}=\delta \phi c^{\prime}+\delta \delta g^{\prime}=\phi \delta c^{\prime}=\phi c c_{1}$, hence $\delta^{*} \mathrm{f}_{1}=c_{1} \in \mathrm{~B}^{\mathrm{p}+1}(U, F)$.

Thus indeed $\delta^{*}$ is a homomorphism between the cohomology groups

$$
\begin{equation*}
\delta^{*}: \mathrm{H}_{\mathrm{a}}^{\mathrm{p}}(U, H) \longrightarrow \mathrm{H}^{\mathrm{p}^{+1}}(U, F) . \tag{A10}
\end{equation*}
$$

The kernel of $\delta^{*}$ consists of those $\mathrm{f} \in \mathrm{Z}_{\mathrm{a}}^{\mathrm{p}}(U, H)$ mapped by $\delta^{*}$ on coboundaries $c \in B^{p^{+1}}(U, F)$. For such an $f$ we have $\delta^{\star} f=c=\delta c^{\prime \prime}$ with $c " \in C^{p}(U, F)$; hence $\psi\left(\mathrm{g}-\phi \mathrm{c}^{\prime \prime}\right)=\psi \mathrm{g}=\mathrm{f}$ and $\delta\left(\mathrm{g}-\phi \mathrm{c}^{\prime \prime}\right)=\phi \mathrm{c}-\phi \delta \mathrm{c}^{\prime \prime}=0$, thus f is the image under $\psi$ of a cocycle in $z^{p}(U, G)$. Conversely, the image $f$ under $\psi$ of a cocycle $g$ in $Z^{\mathrm{P}}(U, G)$ is mapped by $\delta^{*}$ to 0 , since $0=\delta g=\phi \delta^{*} \mathrm{~F}$ and $\phi$ is injective.

The image of $\delta^{*}$ consists of those $c$ in $Z^{p+1}(U, F)$ mapped by $\phi$ into coboundaries of $\mathrm{B}^{\mathrm{p}^{+1}}(l, G)$, because it follows from the construction of $\delta^{*}$ that $\phi \delta^{*} \mathrm{f}=\phi \mathrm{c}=\delta \mathrm{g}$. Conversely, if $\mathrm{c} \in \mathrm{Z}^{\mathrm{p}+1}(U, F)$ is such that $\phi \mathrm{c}=\delta \mathrm{g}$ for some $\mathrm{g} \in \mathrm{C}^{\mathrm{P}}(U, G)$, then $0=\psi \phi \mathrm{c}=\psi \delta \mathrm{g}=\delta \psi \mathrm{g}$, thus $\psi \mathrm{g}=\mathrm{f}$ is a cocycle in $\mathrm{Z}_{\mathrm{a}}^{\mathrm{p}}(U, H)$ with $\delta^{*} f=c$.

Therefore, we have obtained an exact sequence

$$
\begin{array}{r}
0 \longrightarrow H^{0}(U, F) \xrightarrow{\phi^{*}} H^{0}(U, G) \xrightarrow{\psi^{*}} H_{a}^{0}(U, H) \xrightarrow{\delta^{*}} H^{1}(U, F) \xrightarrow{\phi^{*}} H^{1}(U, G)  \tag{A11}\\
\xrightarrow{\psi^{*}} H_{a}^{1}(U, H) \xrightarrow{\delta^{*}} H^{2}(U, F) \longrightarrow \ldots,
\end{array}
$$

where the maps $\phi^{*}$ and $\psi^{*}$ are obtained from $\phi$ and $\psi$ in the obvious way, using the fact that the maps of cochains defined by $\phi$ and $\psi$ commute with the coboundary operators.

We shall now prove that existence theorems for the $\bar{\partial}$-operator are equiv-
alent to statements involving $H^{P}(U, A)$.
THEOREM A13. When $\Omega \in \mathbb{C}^{n}$ is covered by an open covering $U=\left\{U_{i}\right\}$ i ${ }^{n}$, where each $U_{i}$ is pseudoconvex, then for $p \geq 1 \mathrm{H}^{\mathrm{P}}(\mathrm{U}, \mathrm{A})$ is isomorphic to the quotient space
$\left\{\mathrm{f} \mid \mathrm{f}\right.$ is a $(0, \mathrm{p})$-form with $\mathrm{C}^{\infty}-$
coefficients in $\Omega$ and with
$\bar{\partial} \mathrm{f}=0\}$

PROOF. Denote by $E_{q}$ the sheaf of germs of ( $0, q$ )-forms with $C^{\infty}$-coefficients and by $Z_{q}$ the sheaf of germs of $(0, q)$-forms $f$ with $C^{\infty}$-coefficients and with $\bar{\partial} f=0$. Then it follows from theorem A9 that the sequence

$$
0 \rightarrow Z_{\mathrm{q}} \longrightarrow E_{\mathrm{q}} \stackrel{\bar{\partial}}{\longrightarrow} Z_{\mathrm{q}+1} \longrightarrow 0
$$

is exact and that this also holds for the sequence of sections

$$
0 \rightarrow C^{p}\left(U, z_{q}\right) \rightarrow C^{p}\left(U, E_{q}\right) \rightarrow C^{p}\left(U, Z_{q+1}\right) \rightarrow 0
$$

since intersections of pseudoconvex sets are pseudoconvex. Thus $C_{a}^{p}\left(U, Z_{q+1}\right)=$ $=C^{p}\left(U, Z_{q+1}\right)$ and using (A11) and (A9) we get the exact sequence

$$
\begin{aligned}
& 0 \longrightarrow \Gamma\left(\Omega, Z_{q}\right) \longrightarrow \Gamma\left(\Omega, E_{\mathrm{q}}\right) \longrightarrow \Gamma\left(\Omega, Z_{\mathrm{q}+1}\right) \longrightarrow H^{1}\left(U, Z_{\mathrm{q}}\right) \longrightarrow H^{1}\left(U, E_{\mathrm{q}}\right) \longrightarrow \\
& \longrightarrow H^{1}\left(U, Z_{\mathrm{q}+1}\right) \longrightarrow \mathrm{H}^{2}\left(U, Z_{\mathrm{q}}\right) \longrightarrow \mathrm{H}^{2}\left(U, E_{\mathrm{q}}\right) \longrightarrow \ldots
\end{aligned}
$$

Theorem Al2 yields $H^{p}\left(U, E_{q}\right)=0$ for $p \geq 1$ and, therefore, we get

$$
\mathrm{H}^{\mathrm{p}}\left(U, Z_{\mathrm{q}+1}\right) \cong \mathrm{H}^{\mathrm{p}+1}\left(U, Z_{\mathrm{q}}\right), \quad \mathrm{p} \geq 1
$$

and

$$
H^{1}\left(U, Z_{q}\right) \cong \Gamma\left(\Omega, Z_{q+1}\right) / \Gamma\left(\Omega, E_{q}\right)
$$

So using theorem A8 we get for $\mathrm{p} \geq 1$

$$
\begin{aligned}
& \mathrm{H}^{\mathrm{p}}(u, z)=\mathrm{H}^{\mathrm{p}}\left(u, Z_{0}\right) \cong \mathrm{H}^{\mathrm{p}-1}\left(u, Z_{1}\right) \cong \ldots \cong \mathrm{H}^{1}\left(U, Z_{p-1}\right) \cong \\
& \cong \Gamma\left(\Omega, Z_{\mathrm{p}}\right) / \\
& \Gamma\left(\Omega, E_{p-1}\right)
\end{aligned}
$$

In particular, it follows from theorem A9, that if $\Omega$ itself is pseudoconvex

$$
\begin{equation*}
\mathrm{H}^{\mathrm{p}}(\mathrm{U}, \mathrm{~A})=0, \quad \mathrm{p} \geq 1 \tag{A12}
\end{equation*}
$$

for all open coverings $U$ of $\Omega$ consisting of pseudoconvex sets.
This result holds more generally for all coherent analytic sheaves $F$, which is Cartan's theorem B (th.7.4.3 of [7] or th.VIII.A. 14 of [6]). We will prove this only for subsheaves $F$ of $A^{p}$ that are finitely generated by polynomial vectors in $\mathrm{A}(\Omega)^{\mathrm{p}} 1$ ), which is all we need in this paper. For the general case we only indicate where the proof follows the same pattern, which will be sufficient to show why $F$ should be coherent. Moreover, we assume that the covering $U$ is such that more than $M$ distinct sets $U_{i} \in U$ have empty intersection, although this requirement is not necessary (see the footnote on page 98).

THEOREM A14. Let $\Omega$ be an open pseudoconvex set in $\mathbb{C}^{\mathrm{n}}$, let $U$ be an open covering of $\Omega$ consisting of pseudoconvex sets such that the intersection of more than any $M$ elements of $U$ is empty and let $F$ be a subsheaf of $A^{P}$ on $\Omega$ finitely generated by polynomial vectors. Then

$$
\mathrm{H}^{\mathrm{p}}(U, F)=0 \quad \text { for } \mathrm{p} \geq 1
$$

1) The fact that a coherent analytic sheaf $F$ is generated in each point by sections over $\Omega$ is Cartan's theorem A (th.7.2.8 of [7] or th.VIII.A. 13 of [6]). Here we only assume that there is a finite number of sections generating $F$ in all points of $\Omega$ and that these sections consist of polynomials.

PROOF: Let $F$ be generated by $h_{1}, \ldots, h_{q} \in \Gamma(\Omega, F)$, thus each $h_{k}=\left(h_{k}^{1}, \ldots, h_{k}^{r}\right) \in$ $\epsilon A(\Omega)^{r}, k=1, \ldots, q$ and each $h_{k}^{j} \xlongequal[=]{\text { not }} P_{j k}$ is a polynomial. Let us suppose that the problem of section III is solved, that is the sequence (A8) is exact when $\Omega$ is pseudoconvex. This means that the cochains in $C^{P}(U, F)$ are liftable, hence from the exact sequence

$$
0 \longrightarrow R_{P} \longrightarrow A^{q} \xrightarrow{P} F \longrightarrow 0
$$

we get the exact sequence

$$
0 \rightarrow c^{p}\left(U, R_{P}\right) \rightarrow c^{p}\left(U, A^{q}\right) \longrightarrow c^{p}(U, F) \longrightarrow 0
$$

As in (All) we obtain the exact sequence

$$
H^{p}\left(U, A^{q}\right) \longrightarrow H^{p}(U, F) \longrightarrow H^{p+1}\left(U, R_{p}\right) \longrightarrow H^{p+1}\left(U, A^{q}\right)
$$

From (A12) it follows that the right and left hand terms are zero for $p \geq 1$, thus

$$
\mathrm{H}^{\mathrm{p}}(U, F) \cong \mathrm{H}^{\mathrm{p}^{+1}}\left(U, R_{\mathrm{P}}\right)
$$

From (A6) it follows that also $R_{p}$ is a sheaf which is finitely generated by polynomial vectors. Thus if we have shown that $H^{t+1}(U, G)=0$ for every sheaf $G$ finitely generated by polynomial vectors, it follows that $H^{t}(U, H)=0$ for every sheaf $H$ finitely generated by polynomial vectors, $t \geq p$, in particular that $H^{p}(U, F)=0$. But $H^{M}(U, G)=0$, hence the theorem is proved.

The above proof is based on the fact that when $F$ is a sheaf which is finitely generated by sections $h_{1}, \ldots, h_{q}$, then also $R\left(h_{1}, \ldots, h_{q}\right)$ is a sheaf with this property. For that reason we had to require that the vector $h_{k}$ consists of polynomials (see (A6)). In the general case, $F$ is just a coherent analytic sheaf. Then it follows from Cartan's theorem A (see footnote on page 95) and the Heine-Borel theorem that $F$ is finitely generated by sections $h_{1}, \ldots, h_{q} \in \Gamma(\Omega, F)$ in the interior $\Omega^{\prime}$ of any compact subset of $\Omega$. Let $U^{\prime}$ be the covering of $\Omega^{\prime}$ consisting of the sets $U_{i}^{\prime}=U_{i} \cap \Omega$. Since $F$ is coherent,
also the sheaf $R\left(h_{1}, \ldots, h_{q}\right)$ is finitely generated in $\Omega^{\prime}$ by sections over $\Omega$, and the above shows that ${ }_{H}{ }^{p}\left(U^{\prime}, F\right)=0$ for $p \geq 1$. For the passage from $U^{\prime}$ to $U$ see parts (a) and (b) of the proof of theorem 7.4.3. in [7].

We still have to prove that (A8) is exact. Briefly we can say that by definition all cochains are locally liftable and that by theorem Al4 locally liftable cochains are globally liftable in $\Omega$, when $\Omega$ is pseudoconvex. Let us investigate this statement more precisely.

We assume that either $\Omega$ is an open pseudoconvex set whose closure is compact in the open pseudoconvex set $\Omega^{\prime}$ and that $F$ is a coherent analytic sheaf on $\Omega^{\prime}$, or that $\Omega$ is an open pseudoconvex set and that $F$ is a coherent analytic sheaf on $\Omega$ such that $F$ is generated in any point of $\Omega$ by finitely many sections $H^{j}, j=1, \ldots, q$, over $\Omega$, such that $R_{H}$ is generated in any point by finitely many sections $S_{1}^{j}=\left(S_{1}^{j 1}, \ldots, S_{1}^{j q}\right), j=1, \ldots, r$, over $\Omega, \ldots$, such that $R_{S_{m-1}}$ is generated by finitely many sections $\mathrm{S}_{\mathrm{m}}^{\mathrm{j}}=$ $=\left(S_{m}^{j 1}, \ldots, S_{m}^{j r_{m-1}}\right), j=1, \ldots, r_{m}$, over $\Omega, m=2,3, \ldots$. For example, when $F$ is generated by polynomial vectors, we deal with the last case. In both cases we can find for any $z_{j} \in \Omega$, any $m$ and any $f \in \Gamma(\Omega, F),\left(c^{k}\right)_{z_{j}} \in\left(R_{S_{k}}\right)_{z_{j}}$, $k=0,1, \ldots, m-1, S_{0}=H$, an open neighborhood $\omega_{j}^{m}$ of $z_{j}$ in $\Omega$, such that in the following sequence

$$
\begin{align*}
& A\left(\omega_{j}^{m}\right)^{r_{m}} \xrightarrow{S_{m}} A\left(\omega_{j}^{m}\right)^{r_{m-1}} \xrightarrow{S_{m-1}} \ldots \quad . \quad A\left(\omega_{j}^{m}\right)^{r_{1}} \xrightarrow{S_{1}}  \tag{A13}\\
& A\left(\omega_{j}^{m}\right)^{q} \xrightarrow{H} \Gamma\left(\omega_{j}^{m}, F\right)
\end{align*}
$$

$\left.\mathrm{f}\right|_{\omega \mathrm{j}}$ belongs to the image of $H$ and $c^{k} \in A\left(\omega_{j}^{m}\right)^{r} k$ belongs to the image of $S_{k+1}$ for $k \stackrel{\omega_{j}}{=}{ }_{0,1}, \ldots, m-1 \quad\left(r_{0}=q\right) \cdot \omega_{j}^{m}$ depends moreover on $f$ and $c^{k}, k=0, \ldots, m-1$ and it is clear that the above property also holds with $\omega_{j}^{m}$ replaced by an open subset of $\omega_{j}^{m}$, hence $\omega_{j}^{m+1} \subset \omega_{j}^{m}$. Now $U^{(k)}=\left\{\omega_{j}^{k} \mid z_{j} \in \Omega\right\}$ is an open covering of $\Omega$ and $U^{(l)}$ is an open refinement of $U^{(k)}$ when $\ell>k$; we denote the restriction map from $C^{p}\left(U^{(k)}, G\right)$ into $C^{p}\left(U^{(\ell)}, G\right)$ induced by the map from $U^{(l)}$ into $U^{(k)}$ by $\rho_{k, \ell}(G$ is any sheaf on $\Omega)$.

Actually we will show that there is an open refinement $V$ of $U^{(0)}$ such that in the exact sequence (All) $f$ is liftable and that $\delta^{*}$ maps $f$ onto a coboundary of $B^{1}\left(V, R_{H}\right)$, that is $\delta^{*} f=0$, so that $H$ is surjective. The proof is in fact the same as that of theorem Al4, only we develop the sequence
(All) explicitely using (A10).
Let $f \in \Gamma(\Omega, F)$, then $f=H_{j}^{0}$ in $\omega_{j_{0}}^{0}$ for some $g_{j}^{0} \in A\left(\omega_{j}^{0}\right)^{r} 0$, and we regard $f$ as a cocycle in $C^{0}\left(U^{(0)}, F\right)$. The set $g^{0}=\left\{g_{j}^{0} \mid z_{j} \in \Omega\right\}$ determines a cochain in $C^{0}\left(U^{(0)}, A^{r}\right)$. Let $c^{0}=\delta g^{0}$, then $\mathrm{Hc}^{0}=\delta H^{0}=\delta \mathrm{f}=0$, hence $c^{0}$ is a cocycle in $C^{1}\left(U^{(0)}, R_{H}\right)$ (in fact $c^{0}=\delta^{*} f$ by (A10) and (A5) with $P=H$ ). According to (A13) there is a $g^{1} \in C^{1}\left(U^{(1)}, A^{r}\right)$ with $\rho_{0,1} C^{0}=S_{1} g^{1}$. Let $c^{1}=\delta g^{1}$, then $S_{1} c^{1}=\delta S_{1} g^{1}=\rho_{0,1} \delta c^{0}=0$, hence $c^{1} \in C^{2}\left(U^{(1)}, R_{S_{1}}\right)$ (in fact $c^{1}=\delta^{*} \rho_{0,1} c^{0}$ by (A10) and (A6) with $P=H$ and $Q=S_{1}$ ).

Generally we find cochains $g^{k} \in C^{k}\left(U^{(k)}, A^{r_{k}}\right)$ and cocycles

$$
c^{k}=\delta g^{k} \in C^{k+1}\left(U^{(k)}, R_{S_{k}}\right), \quad k=0,1, \ldots, m
$$

since $S_{k} c^{k}=\delta S_{k} g^{k}=\rho_{k-1, k} \delta c^{k-1}=0$, so that

$$
\rho_{k, k+1} c^{k}=s_{k+1} g^{k+1}, \quad k=0,1, \ldots, m-1
$$

In the next section we show that any open covering of $\Omega$ has a refinement consisting of pseudoconvex open sets such that the intersection of more than $M$ of these sets is empty. Let $m=M-1$ and let $V$ be such a refinement of $U^{(m)}$; we denote the restriction map from $C^{p}\left(U^{(k)}, G\right)$ into $C^{p}(V, G)$ by $\rho_{k}$. Now $c^{M-1}=0^{1)}$, so certainly we may write $\rho_{M-1} c^{M-1}=\delta \widetilde{c}^{M-1}$ with $\tilde{c}^{M-1} \in C^{M-1}\left(V, R_{S_{M-1}}\right)$. Assume that for $k \leq M-1$

$$
\rho_{k} c^{k}=\delta \tilde{c}^{k}, \quad \tilde{c}^{k} \in C^{k}\left(V, R_{S_{k}}\right)
$$

Let $\mathrm{g}^{\mathrm{k}}=\rho_{\mathrm{k}} \mathrm{g}^{\mathrm{k}}-\tilde{c}^{\mathrm{k}}$, then $\delta \tilde{\mathrm{g}}^{\mathrm{k}}=\rho_{\mathrm{k}} \mathrm{c}^{\mathrm{k}}-\rho_{\mathrm{k}} \mathrm{c}^{\mathrm{k}}=0$. Since $\Omega$ is pseudoconvex, by (A12) there is a cochain $f^{k-1} \in C^{k-1}\left(V, A^{r}\right)$ with $\mathfrak{g}^{k}=\delta f^{k-1}$. Then we define $\tilde{c}^{k-1}=S_{k} f^{k-1}$, so that $\tilde{c}^{k-1} \in C^{k-1}\left(V, R_{S_{k-1}}\right)$ and

$$
\delta \tilde{c}^{k-1}=S_{k} \delta f^{k-1}=S_{k} \tilde{g}^{k}=\rho_{k} S_{k} g^{k}-S_{k} \tilde{c}^{k}=\rho_{k} \rho_{k-1, k} c^{k-1}=\rho_{k-1} c^{k-1} .
$$

1) The Hilbert syzygy theorem says that $R_{S_{n}}=0$, hence $c^{n}=0$, see [6] IV.C.th.4. So, neither here nor in theorem Al4 we have to require that more than $M$ sets of the covering have empty intersection. However, the Hilbert syzygy theorem is not proved here.

Thus this holds for all k , in particular for $\mathrm{k}=0$ :

$$
\rho_{0} c^{0}=\delta \tilde{c}^{0}, \quad \tilde{c}^{0} \in C^{0}\left(V, R_{H}\right)
$$

(that is $\rho_{0} c^{0}$ is a coboundary, thus $\delta^{*} \rho_{0} f=0$ ). Hence we have

$$
\begin{array}{lll}
f=H\left(g_{k}^{0}-\tilde{c}_{j}^{0}\right) & \text { in } & v_{j} \in V, k=\rho_{0}(j) \\
f=H\left(g_{l}^{0}-c_{i}^{0}\right) & \text { in } & v_{i} \in V, \ell=\rho_{0}(i),
\end{array}
$$

while

$$
\left(\delta \tilde{c}^{0}\right)_{j i}=\tilde{c}_{i}^{0} \tilde{c}_{j}^{0}=\left(\rho_{0} c^{0}\right)_{j i}=\left(\delta g^{0}\right)_{k \ell}=g_{\ell}^{0}-g_{k}^{0}
$$

yields

$$
\mathrm{g}_{\mathrm{k}}^{0}-\tilde{c}_{j}^{0}=\mathrm{g}_{\ell}^{0}-\tilde{c}_{\mathrm{i}}^{0} \quad \text { in } \quad \mathrm{v}_{\mathrm{j}} \cap \mathrm{v}_{\mathrm{i}} .
$$

Thus there is a holomorphic vector function $g \in A(\Omega){ }^{q}$ with $f=H g$ in $\Omega$, namely for all j

$$
g=g_{\rho_{0}(j)}^{0}-\tilde{c}_{j}^{0} \quad \text { in } \quad V_{j} \in V
$$

So we have solved the main problem of this appendix:

THEOREM A15. When $\Omega$ is pseudoconvex, the following sequence is exact

$$
\begin{equation*}
0 \longrightarrow \Gamma\left(\Omega, R_{P}\right) \longrightarrow \Gamma\left(\Omega, A^{q}\right) \xrightarrow{P} \Gamma(\Omega, F) \longrightarrow 0 \tag{A8}
\end{equation*}
$$

We have proved theorem All too, so that the sequence (Al3) is exact for any open pseudoconvex set $\omega_{j}^{m}$. In the general case when $F$ is not generated by polynomial vectors, theorem A of CARTAN (see footnote on page 95) and consequently theorem $B$, as we have shown, follow from the next result due to CARTAN ([7] th.7.2.1.(ii)):
Let $\Omega$ be pseudoconvex and K a compact subset of $\Omega$ with $\mathrm{K}=\hat{\mathrm{K}}_{\Omega}$ (see (A1)) and let $h_{1}, \ldots, h_{q}$ be sections over a neighborhood of $K$ of a coherent analytic sheaf $F$ on a neighborhood of $K$, which generate $F$ there. If $f$ is an arbitrary
section of $F$ over a neighborhood of $K$, then there are $g_{1}, \ldots, g_{q}$ analytic in a neighborhood of $K$ so that $f=\sum_{k=1}^{q} h_{k} g_{k}$ there.

In section 7 we give a quantitative version of theorems A14 and A15, when $H$ is a polynomial matrix. For that purpose we need a quantitative version of the above semilocal result. This is the following modification of th.7.6.5 in [7], which is actually proved there (or th. III 3.4.(3) when $\mathrm{p}=\mathrm{q}=1$ and the general case is contained in th. III 3.6 in [3]):

THEOREM A16. For any polynomial matrix $P$, there is an integer $t>1$, such that for any neighborhood $\omega$ of 0 and every $u \in A(t \omega+z)^{q}, z \in \mathbb{C}^{n}$, there is $a \mathrm{v} \in \mathrm{A}(\omega+z)^{\mathrm{q}}$ with $\mathrm{Pv}=\mathrm{Pu}$ and

$$
\begin{equation*}
\sup _{w \in \omega+z}|v(w)| \leq C(1+\|z\|)^{N} \sup _{w \in t \omega+z}|P(w) u(w)|, \tag{A14}
\end{equation*}
$$

where C is a constant depending on P and $\omega$ (the smaller $\omega$ the larger C ) and where N only depends on P. Here $\mathrm{t} \omega+\mathrm{z}$ denotes $\{\mathrm{w} \mid \mathrm{w}=\mathrm{t} \zeta+\mathrm{z}, \zeta \in \omega\}$ and $|v(w)|^{2}=\sum_{k=1}^{q}\left|v_{k}(w)\right|^{2}$.

In section 7 we perform all the steps of the proofs in this section again, then taking care of the bounds. The sets $\omega+z_{j}$ and $t \omega+z_{j}$ in theorem A16 in fact will be just the sets $\omega_{j}^{m}$ and $\omega_{j}^{m-1}$, respectively, in a quantitative semilocal version of (A13).

## V. SPECIAL COVERINGS

In this section we show that any open covering of the open set $\Omega$ has a refinement that satisfies properties (A15)(i) and (ii) below. This is based on a theorem of dimension theory, th. $3, \S 2, \mathrm{Ch} .7$, p. 278 [4]. Moreover, we construct a special covering of $\Omega$ with refinements satisfying some additional properties needed in section 7. The essential idea for this construction has already been used by WHITNEY in [17], whcih can be found in [8] too.

Let $\Omega$ be an open set in $\mathbb{R}^{n}$ and let $0=\left\{0_{\alpha}\right\}_{\alpha \in \mathrm{A}}$ be an open covering of $\Omega$. Each point in $\Omega$ has a bounded open neighborhood whose closure in $\mathbb{R}^{n}$ is contained in some open set $O_{\alpha}$. Hence there exists an open refinement of $O$ consisting of open sets whose closure is compact and contained in $\Omega$.

Since $\Omega$ is paracompact, we can find a locally finite open refinement $V=\left\{V_{j}\right\}_{j=1}^{\infty}$, where each $\bar{v}_{j}$ is compact and contained in $\Omega$ (such a refinement is necessarily countable, because $\Omega$ is separable). According to [4] 7.2.th. 3 and 7.3.th. 13 V has an open shrinking $W=\left\{W_{j}\right\}_{j=1}^{\infty}$ (which, therefore, is also locally finite and $W_{j}$ cc $\Omega$ ), such that more than $n+1$ distinct sets $W_{j}$ have empty intersection (that is the order of $W$ is at most $n$ ). Since $\Omega$ is normal, lemma 1 to th. $3, \S 1, \mathrm{Ch} .5$ [4] yields an open shrinking $W^{\prime}=\left\{W_{j}!\right\}_{j=1}^{\infty}$ of $W$ such that the closure with respect to $\Omega$ of each $W_{j}^{\prime}$ is contained in $W_{j}$, but since $\bar{W}_{j} \subset \Omega$, this yields $\bar{W}_{j}^{\prime} \subset W_{j}$ for all $j$. Of course $W^{\prime}$ is a locally finite open covering of order at most $n$.

For each $\mathrm{j} \bar{W}_{j}^{\prime}$ is compact and contained in $W_{j}$ and, therefore, we can find finitely many open convex sets $U_{j}, k, k=1, \ldots, m_{j}$ with $U_{j, k} \subset W_{j}$ and $W_{j}^{\prime} \subset U_{k=1}^{m} U_{j, k}$, such that more than $M^{\prime}{ }^{j}$ distinct sets $U_{j, k}$ have empty intersection, where $M^{\prime}$ is a positive integer independent of $j$. For example, this can be done by covering $W_{j}^{\prime}$ by sufficiently small closed hypercubes in $W_{j}$ (so, that the vertices form a rectangular lattice) and by taking sufficiently small convex open neighborhoods of these cubes. Then we get $M^{\prime}=2^{n}$, but it is also possible (by choosing sufficiently small convex open neighborhoods of some cubes and sufficiently large convex open sets contained in the other cubes) to obtain $M^{\prime}=n+1$.

Since each point in $\Omega$ has a neighborhood that intersects a finite number of the sets $W_{j}$, this neighborhood also intersects a finite number of the sets $U_{j, k}$. Furthermore, each point in $\Omega$ is contained in at least one set $W_{j}^{\prime}$ and in at most $n+1$ sets $W_{j}$, thus in at least one and at most $M=(n+1) M^{\prime}$ sets $U_{j, k}$. Therefore, the covering $U=\left\{U_{j, k}\right\}_{j=1, ~}^{\infty}{ }_{k=1}^{m}$ is a locally finite open refinement of $O$ consisting of convex open sets, such that more than $M$ distinct sets of $U$ have empty intersection.

Since convex sets in $\mathbb{C}^{\mathrm{n}}$ are pseudoconvex, we have obtained the
COROLLARY. Let $\Omega$ be an open set in $\mathbb{C}^{\mathrm{n}}$ and let 0 be an open covering of $\Omega$. Then there exists a locally finite open refinement $U=\left\{U_{i}\right\}_{i=1}^{\infty}$ of 0 with the properties
(i) for every $i U_{i}$ is pseudoconvex and $U_{i} \subset \subset \Omega$,
(ii) there is an integer $M$ such that more than $M$ distinct sets in $U$ have empty intersection.

Now we construct coverings of $\Omega$ that satisfy some additional properties. Let $\Omega$ be an open set in $\mathbb{C}^{\mathrm{n}}$ and let $\left\{\Omega_{\mathrm{k}}\right\}_{\mathrm{k}=1}^{\infty}$ be an increasing sequence of open subsets with union $\Omega$ and with

$$
\forall \mathrm{k}, \exists \varepsilon=\varepsilon(\mathrm{k}): \Omega_{\mathrm{k}}(\varepsilon) \subset \Omega_{\mathrm{k}+1}
$$

where $\Omega_{k}(\varepsilon)$ is the $\varepsilon$-neighborhood of $\Omega_{k}$.
Choose positive integers $m_{k}$ with $m_{k+1}>m_{k}$ such that a cube with side $1 / m_{k}$ is contained in the ball with radius $\varepsilon(k)$, for $k=1,2, \ldots$ and let $m_{0}=1$. Divide $\mathbb{C}^{\mathrm{n}}$ into a collection $U^{\prime}$ of closed cubes with side 1 (such that the vertices form a rectangular lattice) and select those cubes contained in $\Omega_{1}$. Call the collection of these cubes $U_{0}^{\prime}$. Divide the remaining cubes and parts of cubes of $U^{\prime}$ into a collection of cubes with side $1 / m_{1}$ and let $U_{1}^{\prime}$ be the collection of those cubes that are contained in $\Omega_{2}$. Generally when we have defined sets $U_{0}^{\prime}, \ldots, U_{k-1}^{\prime}$ of cubes, we define the set $U_{k}^{\prime}$ of cubes obtained by dividing the remaining cubes and parts of cubes of $U$ into a collection of cubes with side $1 / m_{k}$ and by selecting those cubes that are contained in $\Omega_{k+1}$.

Then the union of $U_{0}^{\prime}, U_{1}^{\prime}, \ldots$ covers $\Omega$, since $\Omega_{k}$ is covered by the union of $U_{0}^{\prime}, U_{1}^{\prime}, \ldots, U_{k}^{\prime}$. For, a point $x \in \Omega_{k}$ either belongs to one cube of $U_{0}^{\prime}$ or $\ldots$ or $U_{k-1}^{\prime}$, since these cubes are all contained in $\Omega_{k}$, or it belongs to some cube of $U_{k}^{\prime}$, since any cube with side $1 / m_{k}$ containing $x$ is contained in $\Omega_{k+1}$. Hence any cube in $U_{k}^{\prime}$ can intersect only cubes of $U_{\ell}^{\prime}$ for $\ell=k-1$, $k$ or $k+1$. Furthermore, the intersection of more than $2^{2 n}$ distinct cubes is empty.

Now we will define sufficiently small open neighborhoods of the cubes of $U_{0}^{\prime}, U_{1}^{\prime}, \ldots$, so that we get an open covering. Define the map $\alpha$ by mapping a cube $K^{\prime} \in U_{k}^{\prime}$ to the enlargement of the interior of $K^{\prime}$ by a factor $1+m_{k} / m_{k+1}$, the center kept fixed. Then $\alpha K^{\prime}=K$ is an open cube. Let for each $k U_{k}^{(0)}$ be the set $U_{k}^{(0)}=\left\{\alpha K^{\prime} \mid K^{\prime} \in U_{k}^{\prime}\right\}$ and let

$$
u^{(0)}=\bigcup_{k \geq 0} u_{k}^{(0)}
$$

Then $U^{(0)}$ is a covering of $\Omega$ that satisfies besides properties (A15)(i) and (ii) the following properties for $\lambda=0$
(A15) (iii) all the sets in the covering $U^{(\lambda)}$ intersecting $\Omega_{k}$ have a minimum size and are contained in $\Omega_{\ell(k)}$ with $\ell(k)=k+3$;
(iv) when a set in $U^{(\lambda)}$ intersects $\Omega_{k}$, it intersects not more than $\mathrm{N}_{\mathrm{k}}^{(\lambda)}$ elements of the covering $U_{\mathrm{k}}^{(\lambda)}$, where $\mathrm{N}_{\mathrm{k}}^{(\lambda)}$ is some number depending only on $k$.
The proof follows from the fact that two cubes $K_{j}$ and $K_{i}$ in $U^{(0)}$ have a non-empty intersection if and only if $\alpha^{-1} K_{j}=K_{j}^{\prime}$ intersects $\alpha^{-1} K_{i}=K_{i}^{\prime}$. To prove this, assume that $K_{j}^{\prime} \in U_{k}^{\prime}, K_{i}^{\prime} \in U_{l}^{\prime}, \ell \geq k$, thus $\ell-k=m \geq 0$ and that $K_{j}^{\prime} \cap K_{i}^{\prime}=\emptyset$. Since cubes in $U_{p}^{\prime}$ can intersect cubes in $U_{q}^{\prime}$ on $1 y$ when $q=p-1$, $p$ or $p+1$, the distance between $K_{j}^{!}$and $K_{i}^{\prime}$ is at least

$$
\frac{1}{m_{k+1}}+\frac{1}{m_{k+2}}+\ldots+\frac{1}{m_{k+m-1}} \geq \frac{1}{m_{k+1}}
$$

when $m \geq 2$, or

$$
\frac{1}{m_{k+1}}
$$

when $m$ equals zero or one. The distance from the boundary of $K_{j}$ to $K_{j}$ is by definition of $\alpha$

$$
\begin{equation*}
\frac{1}{2}\left[\frac{1}{m_{k}}\left(1+\frac{m_{k}}{m_{k+1}}\right)-\frac{1}{m_{k}}\right]=\frac{1}{2 m_{k+1}} \tag{A16}
\end{equation*}
$$

and the distance from the boundary of $K_{i}$ to $K_{i}^{\prime}$ is $1 /\left(2 \mathrm{~m}_{\ell+1}\right)$, so that the distance between $K_{j}$ and $K_{i}$ is at least

$$
\frac{1}{m_{k+1}}-\frac{1}{2 m_{k+1}}-\frac{1}{2 m_{k+m+1}}>\frac{1}{m_{k+1}}-\frac{1}{2 m_{k+1}}-\frac{1}{2 m_{k+1}}=0
$$

when $m \geq 1$, or

$$
\frac{1}{m_{k+1}}-\frac{1}{2 m_{k+1}}-\frac{1}{2 m_{k+1}}=0
$$

when $m=0$. Only in this case the boundaries of $K_{j}$ and $K_{i}$ might touch each other, but since $K_{j}$ and $K_{i}$ are open, $K_{j} \cap K_{i}=\emptyset$.

Now property (A15)(ii) follows from the same property for the cubes of
$U_{0}^{\prime}, U_{1}^{\prime}, \ldots$. Let the cube $K$ in $U^{(0)}$ intersect $\Omega_{k}$. If $\alpha^{-1} K$ does not intersect $\Omega_{k+1}, \alpha^{-1} K$ does not intersect the elements of $U_{0}^{\prime}, \ldots, U_{k}^{\prime}$, hence $K$ does not intersect the elements of $U_{0}^{(0)}, \ldots U_{k}^{(0)}$, the union of which contains $\Omega_{k}$. Thus $\alpha^{-1} \mathrm{~K}$ intersects $\Omega_{k+1}$, hence $\alpha^{-1} \mathrm{~K}$ is contained in $\Omega_{k+2}$, so that $K$ is contained in $\Omega_{k+3}$. Thus $K$ has a minimum size, namely the size of $\alpha^{-1} K$ is at least $1 / \mathrm{m}_{\mathrm{k}+2}$. Property (A15)(iv) follows from property (iii) and the same property for $U_{0}^{\prime}, U_{1}^{\prime}, \ldots$.

Finally we construct open refinements $U^{(\lambda)}$ of the covering $U^{(0)}$ satisfying besides the properties (A15)(i), (ii), (iii) and (iv) the following properties
(A15)(v) for each $\lambda U^{(\lambda+1)}$ is a refinement of $U^{(\lambda)}$ and moreover each open cube $K_{j}^{(\lambda)} \epsilon U^{(\lambda)}$ enlarged $2^{\lambda-\mu}$ times with the center kept fixed is contained in some $K_{i}^{(\mu)} \in U^{(\mu)}$ for every $\mu=0,1, \ldots, \lambda-1$; we denote the map $\rho$ between the index sets of $U^{(\lambda)}$ and $U^{(\mu)}$ with $\rho(i)=i_{\mu}$ by $\rho_{\mu, \lambda}$;
(vi) when $K_{j}^{(\lambda)} \in U^{(\lambda)}$ intersects $\Omega_{k}$, there are at most $M_{\lambda, \mu}(k)$ indices $i_{p}$ with $\rho_{\lambda, \mu}\left(i_{p}\right)=j, p=1, \ldots, M_{\lambda, \mu}(k) \quad(\mu>\lambda)$.
Eventually by taking larger integers $m_{k}$, we may assume that each $m_{k+1}=$ $=p_{k+1} m_{k}$ for some integer $p_{k+1} \geq 2, k=0,1, \ldots$. Let $m_{k}^{(0)}=m_{k}$ and $1 e t$ $m_{k}^{(\lambda)}=2^{\lambda} m_{k+\lambda}$ for $\lambda=1,2, \ldots$, then $m_{k}^{(\lambda)}=2 p_{k+\lambda} m_{k}^{(\lambda-1)}$ and $m_{k}^{(\lambda)^{k}}=2 m_{k+1}^{(\lambda-1)}$; $\left(\mathrm{m}_{\mathrm{k}}^{\mathrm{k}}(\lambda)\right)^{-1}, \mathrm{k}=0,1, \ldots$, will be the length of the sides of the closed cubes the covering $U^{(\lambda)}$ is constructed from similarly to the construction of $U^{(0)}$. Namely, let $K_{\lambda}^{\prime}$ be a closed cube with side $\left(\mathrm{m}_{\mathrm{k}}^{(\lambda)}\right)^{-1}$, then the enlargement with a factor $\left(1+m_{k}^{(\lambda)} / m_{k+1}^{(\lambda)}\right.$ ) of the interior of $K_{\lambda}^{\prime}$ will be a cube $K_{\lambda}$ of $U^{(\lambda)}$ as in the construction of $U^{(0)}$. Then $U^{(\lambda)}$ satisfies the same properties on $u^{(0)}$. So, let us assume that the coverings $U^{(0)}, \ldots, U^{(\lambda-1)}$ with the desired properties have been constructed in the same way as $U^{(0)}$ have been constructed from closed cubes.

We divide the closed cubes $K_{\lambda-1}^{\prime}$, with side $\left(\mathrm{m}_{\mathrm{k}}^{(\lambda-1)}\right)^{-1}$, the sets $K_{\lambda-1} \epsilon$ $\epsilon U^{(\lambda-1)}$ are constructed from into $\left(2 p_{k+\lambda}\right)^{2 n}$ closed cubes $K_{\lambda}^{\prime}$ with side $\left(m_{k}^{(\lambda)}\right)^{-1}$ and the covering $U^{(\lambda)}$ is defined as the set of open cubes $K_{\lambda}$ being the enlargement of the interior of the cubes $K_{\lambda}^{\prime}$ by the above factor, $k=0,1, \ldots$. Then the difference of two times half the side of $K_{\lambda}$ and half the side of $K_{\lambda}^{\prime}$ satisfies

$$
2\left[\frac{1}{2}\left(1+\frac{m_{k}^{(\lambda)}}{m_{k+1}^{(\lambda)}}\right) \frac{1}{m_{k}^{(\lambda)}}\right]-\frac{1}{2} \frac{1}{m_{k}^{(\lambda)}} \leq \frac{1}{m_{k}^{(\lambda)}}=\frac{1}{2 m_{k+1}^{(\lambda-1)}}
$$

where the right hand side equals the distance from the boundary of $K_{\lambda-1}$ to $K_{\lambda-1}^{\prime}$ according to (A16). Hence two times $K_{\lambda}$, with the center kept fixed, is contained in $K_{\lambda-1}$, so that property (A15)(v) follows. Furthermore, $K_{\lambda-1}^{\prime}$ contains $\left(2 p_{k+\lambda}\right)^{2 n}$ cubes $K_{\lambda}^{\prime}$, hence $\rho_{\lambda-1, \lambda}$ maps not more than $\left(2 p_{k+\lambda}\right)^{2 n}$ sets $K_{\lambda}$ onto the same $K_{\lambda-1}$. From this and from property (A15) (iii) the above property (Al5)(vi) follows.

REMARK. A1though we use property (A15) (iv) in section 7, this could be avoided. However, the coverings $U^{(\lambda)}$ satisfy (A15) (iv) anyhow.

## VI. NULLSTELLENSATZ AND FUNDAMENTAL PRINCIPLE

In this section we discuss Hilbert's Nullstellensatz, Ehrenpreis' generalization and fundamental principle.

Consider an ideal $I_{z}^{\prime}$ in $A_{z}$ generated by the germs $\left(h_{1}\right)_{z}, \ldots,\left(h_{q}\right)_{z}$ at $z$ of functions $h_{1}, \ldots, h_{q}$ holomorphic in some neighborhood $U$ of $z$. We define the set

$$
V=\left\{w \mid h_{1}(w)=0, \ldots, h_{q}(w)=0, w \in U\right\}
$$

and let $V_{z}$ be the equivalence class of $V$ under the equivalence relation $V \sim W$ if there is a neighborhood $\omega$ of $z$ with $V \cap \omega=W \cap \omega . V_{z}$ is called the germ at $z$ of $V$. It is clear, that the ideal $I_{z}^{\prime}$ is not trivial only if $h_{1}(z)=\ldots=h_{q}(z)=0$. When $f_{z} \in I_{z}^{\prime}$ we denote by $f$ a holomorphic function in a neighborhood of $z$ such that $f_{z}$ is the germ of $f$ at $z$. Then for any $f_{z} \in I_{z}^{\prime}, z \in V$, there is a neighborhood $\omega$ of $z$ with

$$
f(w)=0, \quad w \in V \cap \omega
$$

Conversely, let us consider the ideal $I_{z}$ in $A_{z}$ of all the germs at $z$ of holomorphic functions vanishing on $V_{z}$, that is
(A17) $\quad I_{z}=\left\{f_{z} \mid\right.$ there is a neighborhood $\omega$ of $z$ and $f \in A(\omega)$ with $f(w)=0$ for $w \in V \cap \omega\}$.

It is clear that $I_{z}$ is an ideal and that $I_{z}^{\prime} \subset I_{z}$.
Hilbert's Nullstellensatz says that, if $f_{z} \in I_{z}$, there is a positive integer $m$ with $f_{z}^{m} \in I_{z}^{\prime}$ or

$$
I_{z}=\operatorname{rad} I_{z}^{\prime}=\left\{f_{z} \mid f_{z}^{m} \in I_{z}^{\prime} \text { for some } m \text { depending on } f_{z}\right\}
$$

see [6] II.E.th.20. Obviously, when $I_{z}^{\prime}$ is a prime ideal this yields ([6] III.A.7)

$$
\begin{equation*}
I_{z}^{\prime}=I_{z} \tag{A18}
\end{equation*}
$$

that is, $f_{z} \in I_{z}$ can be written as $f(w)=\sum_{k=1}^{q} g_{k}(w) h_{k}(w)$ for $w$ in some neighborhood $\omega$ of $z$ and for some $g_{k} \in A(\omega), k=1, \ldots, q$.

EHRENPREIS has generalized this result in the following way (see [3] chapter II): let the functions $h_{1}, \ldots, h_{q}$ be polynomials, let for example

$$
\frac{\partial h_{1}}{\partial z_{1}}(z)=0
$$

(of course, also $h_{1}(z) \ldots=h_{q}(z)=0$ ) and let $V_{z}^{\prime}$ be the germ at $z$ of

$$
\mathrm{V}^{\prime}=\left\{\mathrm{w} \left\lvert\, \frac{\partial \mathrm{h}_{1}}{\partial \mathrm{z}_{1}}(\mathrm{w})=0\right., \mathrm{w} \in \mathrm{U}\right\}
$$

Then we require that $f_{z} \in I_{z}$ moreover satisfies in some neighborhood $\omega$ of $z$

$$
\frac{\partial f}{\partial z_{1}}(w)=0, \quad w \in V^{\prime} \cap \omega_{\bullet}
$$

Now let $W_{z}$ be defined as $\left(V_{z}, V_{z}^{\prime}\right)$, where this should be understood in the following way: a function $f$ holomorphic in a neighborhood $\omega$ of $z$ vanishes on $W_{z}$ if $f$ vanishes on $V \cap \omega$ and $\partial f / \partial z_{1}$ vanishes $V^{\prime} \cap \omega$.

The same can be done for higher order derivatives and the other polynomials $h_{k}$. The characterization of $W_{z}$ is not immediately clear from the polynomials $h_{1}, \ldots, h_{q}$ (see example $4, I I .2$ in [3]). Anyhow, $W_{z}$ can be defined in such a way that, if $I_{z}$ is the ideal in $A_{z}$ of germs of functions vanishing on $W_{z}$, we always have (A18), that is $I_{z}$ is the ideal in $A_{z}$ generated by $\left(h_{1}\right)_{z}, \ldots,\left(h_{q}\right)_{z}(t h . I I 2.4$ of [3]).
$V_{z}$ in the Nullstellensatz is called the germ at $z$ of a variety and $W_{z}$ in Ehrenpreis' formulation is called the germ at $z$ of a multiplicity variety. In case of modules in $A_{z}^{P}$ instead of ideals, it is possible to define $p$ germs $\left(W_{1}\right)_{z}, \ldots,\left(W_{p}\right)_{z}$ of multiplicity varieties and so we get the germ $\vec{W}_{z}=\left(\left(W_{1}\right)_{z}, \ldots,\left(W_{p}\right)_{z}\right)$ of a vector multiplicity variety. This can be done in such a way, that the analogue of (Al8) holds, namely (th.II 2.6 of [3]):

THEOREM A17. Let $P_{j k}$ be polynomials, $j=1, \ldots, p, k=1, \ldots, q$. Then it is possible for each $z$ to define the germ $\vec{W}_{z}$ at $z$ of a vector multiplicity variety, such that each p-tuple of functions $f_{j}, j=1, \ldots, p$, holomorphic in a neighborhood of $z$, whose germ at $z$ vanishes on $\vec{W}_{z}$, can be written as

$$
f_{j}(w)=\sum_{k=1}^{q} P_{j k}(w) g_{k}(w), \quad j=1, \ldots, p
$$

for $w$ in some neighborhood $\omega$ of $z$ and for some functions $g_{k} \in A(\omega)$, $\mathrm{k}=1, \ldots, \mathrm{q}$.

Next we consider a sheaf $I^{\prime}$ of ideals generated in each point of an open pseudoconvex set $\Omega$ by polynomials $h_{1}, \ldots, h_{q}$, thus $p=1$. Their simultaneous zero-set defines a variety $V=\underset{z \in \Omega}{\bigcup_{Z}} V_{z}$ in $\Omega$ (at points $z$ where some $h_{k}(z) \neq 0 V_{z}$ is empty). Similarly we can define a multiplicity variety $W$ in $\Omega$ (see [3]). We will consider sheafs of functions on $V$; the same can be done for a multiplicity variety $W$. Let $I$ be the sheaf on $\Omega$

$$
I=\bigcup_{z \in \Omega} I_{z}
$$

where $I_{z}$ is defined by (A17); note that $I_{z}=A_{z}$ when $z \in \Omega \backslash V$. We can define a sheaf $F$ on $\Omega$ by

$$
F_{z}=A_{z} / I_{z}, \quad z \in \Omega
$$

that is the following sequence is exact

$$
0 \longrightarrow I \longrightarrow A \longrightarrow F \longrightarrow 0
$$

For $z \in \Omega \backslash V I_{z}=A_{z}$, thus $F_{z}=0$. Thus $F$ is only non-trivial in points of V , so we consider the restriction $F^{\prime}$ to V

$$
F^{\prime}=\bigcup_{z \in V} F_{z}
$$

which is a sheaf on $V$. In accordance with the footnote on page 85 we can regard $F^{\prime}$ as the sheaf of germs of analytic functions on $V$. A section $f$ in $\Gamma\left(V, F^{\prime}\right)$ is a holomorphic function in $V$; regarded as a section $f_{1}$ in $\Gamma(\Omega, F)$ we would have $f_{1}(z)=f(z)$ for $z \in V$ and $f_{1}(z)=0$ for $z \in \Omega \backslash V$. So, we may just as well consider the sections in $\Gamma(\Omega, F)$ as the holomorphic functions in $V$. In case $I_{z}^{\prime}$ is a prime ideal for all $z \in \Omega$, (A18) holds and the sheaf $I$ is finitely generated by polynomials. This also holds when we consider a sheaf of ideals on a multiplicity variety. Hence theorem A14 may be applied. A1so, generally for any sheaf $I$ of ideals Cartan's theorem B may be applied, since $I$ is coherent ([6] IV.D.2). However, in the case occurring in this paper $I_{z}^{\prime}$ is prime for all $z \in \Omega$. Hence in the same way as theorem A15 was obtained from Cartan's theorem B we here get

$$
0 \longrightarrow \Gamma(\Omega, I) \longrightarrow \Gamma(\Omega, A) \longrightarrow \Gamma(\Omega, F) \longrightarrow 0,
$$

so that

$$
\begin{equation*}
A(\Omega) /_{\Gamma(\Omega, I)} \cong \Gamma(\Omega, F) \cong \Gamma\left(V, F^{\prime}\right) \tag{A19}
\end{equation*}
$$

Thus any function holomorphic in $V$ is the restriction of a function in $\mathrm{A}(\Omega)$ and in case (A18) holds any function f in $\mathrm{A}(\Omega)$ that vanishes on V can be written as $f(z)=\sum_{k=1}^{q} h_{k}(z) g_{k}(z), z \in \Omega$ for some $g_{k} \in A(\Omega)$. In this paper we will derive a quantitative version of (A19) for a special variety $V$.

When $\Omega=\mathbb{C}^{\mathrm{n}}$, when $I^{\prime}$ is an ideal generated by polynomials and when $W$ is its associated multiplicity variety (thus (A18) holds), the isomorphism (A19) with $V$ replaced by W and with bounds (that is all the occurring functions satisfy moreover certain estimates at infinity) is Ehrenpreis' fundamental principle (theorem IV 4.1 in [3]; a survey of this theorem and its proof can be found in [1] IV). The fundamental principle holds for
modules generated by polynomials too (th. IV 4.2.[3]), however, in that case the definition of global vector multiplicity varieties is a more delicate question (see page 100 [3]).

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