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# ON THE POLYNOMIAL-LIKE BEHAVIOUR OF CERTAIN ALGEBRAIC FUNCTIONS 

by Charles FEFFERMAN(*) \& Raghavan NARASIMHAN

## 0. Introduction.

The purpose of this paper is to prove an extension theorem which, in particular, implies that certain families of algebraic functions have the growth and smoothness properties of polynomials. Our work was motivated by that of Parmeggiani $[\mathrm{P}]$ on pseudodifferential operators. We begin with a few words about $[\mathrm{P}]$ and how algebraic functions enter there.

It is known from the work of Stein and his collaborators [FS], [RS], [NSW] that a subelliptic differential operator is governed by a family of nonEuclidean balls. The purpose of $[\mathrm{P}]$ is to associate non-Euclidean balls in the cotangent bundle to a pseudodifferential operator $P(x, D)$. In dimensions one and two, the results in $[\mathrm{P}]$ provide a geometrical understanding of these non-Euclidean balls. It would be of interest to extend Parmeggiani's work to higher dimensions.

The method used in $[\mathrm{P}]$ requires writing the symbol $P(x, \xi)$ in a normal form; it is likely that extending the results in $[\mathrm{P}]$ to higher dimensions will also require doing this. It is this normal form which brings in algebraic functions. More precisely, fix a point $\left(x^{0}, \xi^{0}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{n}$ where $P\left(x^{0}, \xi^{0}\right)=0$. We assume that $P(x, \xi) \geq 0$. Expanding $P(x, \xi)$ in a Taylor

[^0]series to a high order at $\left(x^{0}, \xi^{0}\right)$ and rescaling, we obtain a polynomial $p(x, \xi)$ bounded a priori on the unit cube $Q=\left\{\left(x_{1}, \ldots, x_{n}, \xi_{1}, \ldots, \xi_{n}\right) \mid\right.$ $\left.\left|x_{j}\right| \leq 1,\left|\xi_{j}\right| \leq 1,1 \leq j \leq n\right\}$. Here $p(0,0)=0$ and $p \geq 0$ on $Q$ modulo a small error, which we ignore. Suppose that the Hessian $\left(\frac{\partial^{2} p}{\partial \xi_{j} \partial \xi_{k}}\right)$ has rank $r$ at the origin. Rotating coordinates, we may assume that $\left(\frac{\partial^{2} p}{\partial \xi_{j} \partial \xi_{k}}\right)_{1 \leq j, k \leq r}$ is positive definite at 0 .

In the neighborhood of the origin, we can then write $p(x, \xi)$ in the normal form

$$
\begin{align*}
& p\left(x, \xi_{1}, \ldots, \xi_{n}\right)=\sum_{k=1}^{r} e_{k}(x, \xi)\left[\xi_{k}-\theta_{k}\left(x, \xi_{r+1}, \ldots, \xi_{n}\right)\right]^{2}  \tag{1}\\
&+q\left(x, \xi_{r+1}, \ldots, \xi_{n}\right)
\end{align*}
$$

where $e_{k}, \theta_{k}, q$ are smooth, and $e_{k}>0$.
Here, $q\left(x, \xi_{r+1}, \ldots, \xi_{n}\right)$ is an algebraic function. In fact, let $V$ be the (real) algebraic variety $V=\left\{(x, \xi) \in \mathbb{R}^{n} \times \mathbb{R}^{n} \left\lvert\, \frac{\partial p}{\partial \xi_{k}}=0\right.,1 \leq k \leq r\right\}$, and let $\pi: V \rightarrow \mathbb{R}^{2 n-r}$ be the projection $\left(x_{1}, \ldots, x_{n}, \xi_{1}, \ldots, \xi_{n}\right) \mapsto$ $\left(x_{1}, \ldots, x_{n}, \xi_{r+1}, \ldots, \xi_{n}\right)$. Near the origin, $V$ is smooth, $\pi$ is a diffeomorphism, and

$$
\begin{equation*}
q=(p \mid V) \circ \pi^{-1} \tag{2}
\end{equation*}
$$

To get useful information from (1), we need to know that $q$ has the growth and smoothness properties of a polynomial. If we are allowed to restrict attention to a tiny neighborhood of the origin, then one can simply Taylor-expand $q$ to obtain a polynomial. However, the method followed in $[\mathrm{P}]$ requires an understanding of $q$ on a neighborhood of fixed size, and one is forced to study functions of the form (2) on such a neighborhood.

We can formulate the basic problem in a simple, general setting. Suppose that we are given a variety $V \subset \mathbb{R}^{n}$, defined as the set of common zeros of polynomials $P_{1}, \ldots, P_{r}, 1 \leq r \leq n-1$. We make the following:

Assumptions.
(I) The polynomials $P_{j}$ have degree at most $D$, and their coefficients have absolute value at most $C$.
(II) We have $P_{1}(0)=\cdots=P_{r}(0)=0$ and $\left|\operatorname{det}\left(\frac{\partial P_{j}}{\partial x_{k}}(0)\right)_{1 \leq j, k \leq r}\right| \geq$ $c>0$.

Let $\pi: V \rightarrow \mathbb{R}^{n-r}$ be the projection $\left(x_{1}, \ldots, x_{n}\right) \mapsto\left(x_{r+1}, \ldots, x_{n}\right)$. In view of (II), $\pi$ has a smooth local inverse $\pi^{-1}: B\left(0, \delta_{0}\right) \rightarrow V$ defined on a small ball. Our goal is then to understand functions of the form $F=p \circ \pi^{-1}$ where $p$ is a polynomial of degree at most $D$ on $\mathbb{R}^{n}$, and to do this uniformly with respect to the polynomials $P_{j}$ satisfying our assumptions.

The following result (a somewhat weaker form of the main theorem 1.1 of this paper) provides very good control over this family of functions.

Theorem 1. - Let $P_{1}, \ldots, P_{r}$ be polynomials satisfying assumptions (I) and (II), and let $V=\left\{x \in \mathbb{R}^{n} \mid P_{1}(x)=\cdots=P_{r}(x)=0\right\}$. Then, there exist constants $\delta_{*}, C_{*}, D_{*}>0$ depending only on the constants $n, c, C, D$ occurring in (I), (II) such that if $p$ is a polynomial of degree at most $D$ on $\mathbb{R}^{n}$, we can find polynomials $f$ and $g$ with the following properties:
(A) $f, g$ have degree at most $D_{*}$;
(B) On $B\left(0,2 \delta_{*}\right)$, we have $\frac{1}{2}<g<2$ and $|f| \leq C_{*} \sup _{V \cap B\left(0, \delta_{*}\right)}|p|$;
(C) $p=f / g$ on $V \cap B\left(0,2 \delta_{*}\right)$.

We note explicitly that the denominator $g$ in this theorem cannot be taken $\equiv 1$; see $[\mathrm{FN}]$. We can, however, take $g$ to depend only on the $P_{j}$, not on $p$.

As an easy consequence of this extension theorem, we prove the following result which shows that the family of algebraic functions mentioned above behaves like a family of polynomials of bounded degree.

Theorem 2. - Let $V$ be as in Theorem 1, and let $\pi: V \rightarrow \mathbb{R}^{n-r}$ be the projection. There exist constants $\delta_{*}, C_{*}>0$ depending only on $n, c, C, D$ (the constants in (I), (II)) such that $\pi$ has a smooth inverse $\pi^{-1}: B\left(0, \delta_{*}\right) \rightarrow V$ and such that if $p$ is a polynomial of degree $\leq D$ and $F=p \circ \pi^{-1}$ on $B\left(0, \delta_{*}\right)$, then the following inequalities hold:
(A) Polynomial Growth. For $0<2 \delta<\delta_{*}$,

$$
\sup _{B(0,2 \delta)}|F| \leq C_{*} \sup _{B(0, \delta)}|F| .
$$

(B) Bernstein's Inequality. If $\nabla F$ denotes the gradient of $F$, then

$$
\sup _{B(0, \delta)}|\nabla F| \leq \frac{C_{*}}{\delta} \sup _{B(0, \delta)}|F| \text { for } 0<\delta<\delta_{*}
$$

(C) Equivalence of Norms.

$$
\sup _{B(0, \delta)}|F| \leq \frac{C_{*}}{\delta^{n-r}} \int_{B(0, \delta)}|F| \text { for } 0<\delta<\delta_{*}
$$

For a fixed $V$ and $\delta$, estimates of this type are obvious consequences of the fact that any two norms on a finite dimensional vector space are equivalent. This simple remark gives no control over the constant $C_{*}$, whereas the extension theorem enables us to control the constants easily. This control is, in turn, crucial in understanding the function $q$ in (1) and (2).

Let us mention also that the extension theorem clearly enables us to estimate the function $F$ of Theorem 2 in a complex ball of radius $\delta_{*}$ around 0 in terms of its values on $B\left(0, \delta_{*}\right) \subset \mathbb{R}^{n-r}$.

The plan of the proof of Theorem 1 given here is as follows. Let $W$ be the space of all $\vec{P}=\left(P_{1}, \ldots, P_{r}\right)$ satisfying assumptions (I) and (II). [We work with a very slightly different space $W$ in the body of the paper.] Fix $\vec{P}_{0} \in W$. We say that Theorem 1 holds locally at $\vec{P}_{0}$ if the conclusions of the theorem hold for all $\vec{P}$ in a sufficiently small neighborhood of $\vec{P}_{0}$, with a constant $C_{*}$ depending on $\vec{P}_{0}$ (but with $\delta_{*}$ and $D_{*}$ independent of $\vec{P}_{0}$ ). The space $W$ is compact. It is therefore sufficient to prove that Theorem 1 holds locally at any $\vec{P}_{0} \in W$. We shall take degeneracies into account by showing, by downward induction on $s, 0 \leq s \leq \operatorname{dim} W$, that Theorem 1 holds locally at $\vec{P}_{0}$ for any $\vec{P}_{0}$ outside a semi-algebraic subset $V_{s} \subset W$ whose dimension is $<s$. Since $V_{0}=\emptyset$ if $s=0$ (since $\operatorname{dim} V_{0}<0$ ), this shows that Theorem 1 holds locally everywhere on $W$, thus completing the proof.

Clearly, the main work lies in the induction step. This is modeled on a simpler induction in our earlier paper [FN] in which the theorem is proved for hypersurfaces (i.e. for $r=1$ ). Unlike the argument in [FN] in which we described the sets $V_{s}$ in terms of factorizing the equation defining $V$, we are unable to describe the sets $V_{s}$ explicitly, and resort instead to extensive use of the general structure of semi-algebraic sets.

It is clear from the statement of Theorem 1 that one needs to analyze polynomials vanishing on one connected component of the smooth part of the set of common zeros of polynomials $P_{1}, \ldots, P_{r}$ satisfying (I) and (II). This analysis, essentially obvious in [FN], has to be done differently here,
and one of our main steps in carrying out the induction outlined above is to establish the following result.

Theorem 3. - Suppose that $\left(P_{1}, \ldots, P_{r}\right)$ satisfy (I) and (II). Let $V_{0}\left(P_{1}, \ldots, P_{r}\right)$ be the connected component, containing 0 , of the smooth part of the zero set of $\left(P_{1}, \ldots, P_{r}\right)$.

Then, there exists an auxiliary polynomial $q$ on $\mathbb{R}^{n}$ with the following properties:
(A) $q(0)=1$, and the degree and coefficients of $q$ are bounded a priori in terms of the constants in (I) and (II).
(B) If $f$ is any polynomial vanishing on $V_{0}\left(P_{1}, \ldots, P_{r}\right)$, then we can write $q f$ in the form

$$
q f=G_{1} P_{1}+\cdots+G_{r} P_{r}
$$

with polynomials $G_{\nu}$ whose degrees are bounded a priori in terms of $\operatorname{deg} f$ and the constants in (I) and (II).

In the case of hypersurfaces, i.e. when $r=1$, this is a simple lemma (see [FN]). In the general case, we prove Theorem 3 by using Hörmander's $L^{2}$-estimates for the $\bar{\partial}$-operator on $\mathbb{C}^{n}$.

For a given $\left(P_{1}, \ldots, P_{r}\right)$, the polynomials $g$ in Theorem 2 will be constructed as a product of polynomials $q$ from Theorem 3 associated to finitely many $\left(P_{1}^{(j)}, \ldots, P_{r}^{(j)}\right)$ that lie near $\left(P_{1}, \ldots, P_{r}\right)$.

Note that Theorem 3 asserts no a priori bounds on the coefficients of the polynomials $\left(G_{1}, \ldots, G_{r}\right)$ in (B) above. While we need bounds on their degree, we need no bounds on the coefficients. It would, however, be interesting to decide what the optimal estimates on these coefficients might be; in the case of a weaker variant of Theorem 3, one can obtain good a priori estimates (see Theorem 5.2 below).

It would of course be of interest to estimate the constant $C_{*}$ in Theorems 1 and 2 as a function of $n, c, C, D$. This seems difficult to do. However, for the application to pseudodifferential operators (in particular for the results of Parmeggiani), the crucial point is to have a constant $C_{*}$ independent of the polynomials $P_{j}$ (as long as assumptions (I) and (II) are satisfied). For this reason, the work in the earlier sections, leading to a proof of Theorem 3 stated above (Theorem 5.5 in the text) has to be uniform in the $P_{j}$.

We are grateful to the referee who has pointed out that the methods
we have used in sections 2-5 have been applied before to related questions. We have added some of the references he gave us in the relevant places.

## 1. Notation and statement of the Main Theorem.

Let $n \geq 2$ be an integer. If $z=\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n}$, we set, as usual, $|z|^{2}=\left|z_{1}\right|^{2}+\cdots+\left|z_{n}\right|^{2}$; if $\alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right) \in \mathbb{N}^{n}$ is an $n$-tuple of nonnegative integers, we set $z^{\alpha}=z_{1}^{\alpha_{1}} \cdots z_{n}^{\alpha_{n}}$.

If $f \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right], f=\sum_{\alpha} f_{\alpha} z^{\alpha}$ is a polynomial with complex coefficients in $n$ variables, we define $\|f\|$ by

$$
\|f\|^{2}=\sum_{\alpha}\left|f_{\alpha}\right|^{2}
$$

If $f=\left(f_{1}, \ldots, f_{r}\right)$ is an $r$-tuple of polynomials, we set $\|f\|^{2}=$ $\sum\left\|f_{j}\right\|^{2}$.

We shall have occasion to use this same notation for polynomials in more than $n$ variables; thus if $m \geq 1, f \in \mathbb{C}\left[z_{1}, \ldots, z_{m}\right], f=\sum f_{\alpha} z^{\alpha}$, then $\|f\|^{2}=\sum\left|f_{\alpha}\right|^{2}$.

If $A$ is an $n \times n$ matrix over $\mathbb{C}$, we denote by $\|A\|$ the operator norm of the linear map $A: \mathbb{C}^{n} \rightarrow \mathbb{C}^{n}$.

Let $d \geq 1$ be any integer. We define

$$
\begin{aligned}
H^{d} & =\left\{f \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right] \mid \text { degree }(f) \leq d\right\} \\
\mathcal{H}^{d} & =\left\{f \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right] \mid \text { degree }(f) \leq d\right\}
\end{aligned}
$$

For $\rho>0$, let

$$
Q_{\rho}=\left\{x=\left(x_{1}, \ldots, x_{n}\right) \in \mathbb{R}^{n}| | x_{j} \mid \leq \rho\right\}, Q=Q_{1}
$$

If $a \in \mathbb{C}^{n}$ and $\rho>0$, we set

$$
\begin{aligned}
& Q(a, \rho)=\left\{z=\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n} \mid\right. \\
&\left.\left|\operatorname{Re}\left(z_{j}-a_{j}\right)\right| \leq \rho,\left|\operatorname{Im}\left(z_{j}-a_{j}\right)\right| \leq \rho, j=1, \ldots, n\right\}
\end{aligned}
$$

If $a \in \mathbb{C}^{n}, R>0$, we set

$$
B(a, R)=\left\{z \in \mathbb{C}^{n}| | z-a \mid<R\right\}
$$

it is the open ball of radius $R$ centered at $a ; \bar{B}(a, R)$ will denote its closure. We use the same notation also in $\mathbb{R}^{n}$. However, since we shall use balls in
$\mathbb{R}^{k}$ or $\mathbb{C}^{k}$ with varying $k$, we write $B_{k}(a, R)=\left\{\left(x_{1}, \ldots, x_{k}\right) \in \mathbb{R}^{k}\left(\mathbb{C}^{k}\right) \mid\right.$ $\left.\sum_{j=1}^{k}\left|x_{j}-a_{j}\right|^{2}<R^{2}\right\}$ when this is relevant.

Fix integers $n \geq 2,1 \leq r \leq n-1, D \geq 1$ and a constant $C_{1}>0$. These will remain fixed throughout the paper.

If $p_{1}, \ldots, p_{r} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]\left(\right.$ resp. $\left.\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]\right)$, we set

$$
\begin{aligned}
J_{P}(x) & =\operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}(x)\right)_{1 \leq j, k \leq r}, x \in \mathbb{R}^{n} \\
\text { (resp. } J_{P}(z) & \left.=\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}(z)\right)_{1 \leq j, k \leq r}, z \in \mathbb{C}^{n}\right)
\end{aligned}
$$

Our basic space is:
(1) $W=\left\{P=\left(p_{1}, \ldots, p_{r}\right) \in H^{D} \times \cdots \times H^{D} \mid P(0)=0,\|P\| \leq C_{1}\right.$, $\left.J_{P}(0)=1\right\}$.

If $P=\left(p_{1}, \ldots, p_{r}\right), p_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$, we set

$$
\begin{equation*}
Z(P)=\left\{x \in \mathbb{R}^{n} \mid p_{1}(x)=\ldots \cdot p_{r}(x)=0\right\} \tag{2}
\end{equation*}
$$

If $P(0)=0$ and $J_{P}(0) \neq 0$, let $V^{0}(P)$ be the connected component of $Z(P) \cap\left\{x \in \mathbb{R}^{n} \mid J_{P}(x) \neq 0\right\}$ containing 0 . We let $V(P)$ be the closure in $\mathbb{R}^{n}$ of $V^{0}(P)$. Note that $V(P)$ is not necessarily an algebraic set.

We use similar notation over $\mathbb{C}$. If $P=\left(p_{1}, \ldots, p_{r}\right), p_{j} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$, we set
(3) $\mathcal{Z}(P)=\left\{z \in \mathbb{C}^{n} \mid P(z)=0\right\}$.

If $P(0)=0, J_{P}(0) \neq 0, \mathcal{V}^{0}(P)$ will stand for the connected component of $\mathcal{Z}(P) \cap\left\{z \in \mathbb{C}^{n} \mid J_{P}(z) \neq 0\right\}$ containing 0 , and $\mathcal{V}(P)$ will denote its closure in $\mathbb{C}^{n}$. In this case, $\mathcal{V}(P)$ is the irreducible component of $\mathcal{Z}(P)$ containing 0 (unique since $\mathcal{Z}(P)$ is smooth at 0 because $J_{P}(0) \neq 0$ ). In particular, $\mathcal{V}(P)$ is an algebraic variety.
(4) There is a constant $\rho_{1}, 0<\rho_{1} \leq 1$ depending only on $C_{1}, n, D$ such that for any $P \in W$, we have $Q_{\rho_{1}} \cap Z(P) \subset V^{0}(P)$ (see the remark following Lemma 2.1). For $0<\rho \leq \rho_{1}$, we set

$$
V_{\rho}(P)=Q_{\rho} \cap Z(P)=Q_{\rho} \cap V(P)=\left\{x \in Q_{\rho} \mid P_{1}(x)=\ldots \cdot P_{r}(x)=0\right\}
$$

(5) For $P \in W$, we denote by $\mathcal{G}(P)$ the space of germs at 0 of functions defined on some neighborhood of 0 on $Z(P)$ (or $V(P)$ ).

We can now state our main theorem.
Theorem 1.1 (The Extension Theorem).
Part 1. There exist constants $D^{\prime} \geq 1, C^{\prime}>0, m>0$, depending only on $C_{1}, n, D$ such that the following holds.

Given $P \in W, \rho>0\left(0<\rho \leq \rho_{1}\right)$, there exists $q \in H^{D^{\prime}}$ with $q(0)=1$, $\|q\| \leq C^{\prime}$ such that, if $f \in H^{D}$, we can find $F \in H^{D^{\prime}}$ for which
(i) $F=q f$ on $V_{\rho_{1}}(P)$,
and
(ii) $\|F\| \leq C^{\prime} \rho^{-m} \sup _{V_{\rho}(P)}|f|$.

Part 2. There exists $\rho_{0}>0$ depending only on $C_{1}, n, D$ such that if $P \in W, f \in H^{D}$ and $0<\rho \leq \rho_{0}$, then, there exist $F, q \in H^{D^{\prime}}$ with the following properties:
(a) $\frac{1}{2}<q<2$ on $Q_{2 \rho}$
(b) $f=F / q$ on $V_{\rho}(P)$
(c) $\sup _{Q_{2 \rho}}|F| \leq C^{\prime} \sup _{V_{\rho}(P)}|f| ;$
here, as in Part 1, $C^{\prime}, D^{\prime}$ depend only on $C_{1}, n, D$.

## 2. Preliminaries.

Let $n \geq 2,1 \leq r \leq n-1, D \geq 1, C_{1}>0$ be given. Consider the space $\mathcal{W}_{0}$ of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right)$ with $p_{j} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right], \operatorname{deg} p_{j} \leq D$ (i.e. $p_{j} \in \mathcal{H}^{D}$ ), with

$$
P(0)=0,\|P\| \leq C_{1}, J_{P}(0)=1\left(J_{P}(z)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}(z)\right)_{1 \leq j, k \leq r}\right)
$$

(1) We introduce a new variable $z_{0}$ and write $\tilde{z}=\left(z_{0}, z\right) \in \mathbb{C}^{n+1}$ with $z=\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n}$. If $P \in \mathcal{W}_{0}$, define $p_{0}(\tilde{z})=p_{0}\left(z_{0}, \ldots, z_{n}\right)=$ $\left(z_{0}+1\right) J_{P}(z)-1$. We have $p_{0}(0)=0$.

Consider the system

$$
\tilde{P}=\left(p_{0}, p_{1}, \ldots, p_{r}\right)
$$

The set $\mathcal{Z}(\tilde{P})$ of common zeroes of $p_{0}, \ldots, p_{r}$ in $\mathbb{C}^{n+1}$ is a smooth affine variety isomorphic to $\mathcal{Z}(P) \cap\left\{z \in \mathbb{C}^{n} \mid J_{P}(z) \neq 0\right\} \subset \mathbb{C}^{n}$ (in the category of algebraic varieties). Moreover, if $\tilde{J}(\tilde{z})=\operatorname{det}\left(\frac{\partial p_{i}}{\partial z_{j}}\right)_{0 \leq i, j \leq r}$, we have $\tilde{J}(\tilde{z})=\left(J_{P}(z)\right)^{2}$, and $\tilde{J}$ satisfies the equation

$$
\left(z_{0}+1\right)^{2} \tilde{J}(\tilde{z})-p_{0}\left(z_{0}, \ldots, z_{n}\right)\left(1+\left(z_{0}+1\right) J_{P}(z)\right)=1
$$

Hence

$$
\begin{aligned}
1 & \leq\left(\left|z_{0}+1\right|^{4}+\left|1+\left(z_{0}+1\right) J_{P}(z)\right|^{2}\right)\left(|\tilde{J}(\tilde{z})|^{2}+\left|p_{0}(\tilde{z})\right|^{2}\right) \\
& \leq C\left(1+|\tilde{z}|^{2}\right)^{N}\left(|\tilde{J}(\tilde{z})|^{2}+\sum_{j=0}^{r}\left|p_{j}(\tilde{z})\right|^{2}\right)
\end{aligned}
$$

where $C, N$ depend only on $C_{1}, n, D$.
This leads us to consider the following space:
(1) Let $n \geq 2,1 \leq r \leq n-1, D \geq 1, C_{1}>0, c_{0}>0$ and $N \geq 1$ be given.

Define $\mathcal{W}$ to be the following space:
$\mathcal{W}$ is the space of $P=\left(p_{1}, \ldots, p_{r}\right), p_{j} \in \mathcal{H}^{D}$, such that
(a) $P(0)=0,\|P\| \leq C_{1}, \operatorname{deg} p_{j} \leq D$ (i.e. $p_{j} \in \mathcal{H}^{D}$ ).
(b) $J_{P}(0)=1$.
(c) $\left|J_{P}(z)\right|^{2}+\left|p_{1}(z)\right|^{2}+\cdots+\left|p_{r}(z)\right|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N}, z \in \mathbb{C}^{n}$.
(Recall that $J_{P}(z)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}(z)\right)_{1 \leq j, k \leq r}$. .)
We shall need certain estimates for these polynomials which can be obtained by examining the inverse function theorem. We do this for holomorphic mappings, but remark, for later use, that they remain valid for smooth maps (see remark at the end of Lemma 2.1). In particular, the statement (4) in $\S 1$ is an easy consequence of the smooth version of Lemma 2.1.

The results in this section are closely related to those in $[\mathrm{BT}]$ and [BY]. See, in particular, $[\mathrm{BT}]$ for a version of Lemma 2.6.

Lemma 2.1. - Let $R_{1}, R_{2}>0$, and let $r, s$ be integers $\geq 1$. Let $f: B_{r}\left(0, R_{1}\right) \times B_{s}\left(0, R_{2}\right) \rightarrow \mathbb{C}^{r}$ be a holomorphic map $\left(B_{k}(a, \rho)\right.$ is the ball of radius $\rho$ in $\mathbb{C}^{k}$ centered at $\left.a \in \mathbb{C}^{k}\right)$.

Suppose that $f(0,0)=0$, and let $D(z, \zeta)$ be the matrix $\left(\frac{\partial f_{j}}{\partial z_{k}}(z, \zeta)\right)_{1 \leq j, k \leq r}, z \in B_{r}\left(0, R_{1}\right), \zeta \in B_{s}\left(0, R_{2}\right)$.

Assume that we are given constants $M_{1}, M_{2}, M_{3} \geq 1$ and $0<\mu \leq 1$ such that the following inequalities hold for $|z|<R_{1},|\zeta|<R_{2}$ :
(a) $\left|\frac{\partial f_{j}}{\partial z_{k}}\right| \leq M_{1}, 1 \leq j, k \leq r,|\operatorname{det} D(z, \zeta)| \geq \mu$.
(b) $\left|\frac{\partial^{2} f_{i}}{\partial z_{j} \partial z_{k}}\right| \leq M_{2}, 1 \leq i, j, k \leq r$.
(c) $\left|\frac{\partial f_{j}}{\partial \zeta_{\ell}}\right| \leq M_{3}, 1 \leq j \leq r, 1 \leq \ell \leq s$.

Then, there exists a constant $\gamma_{r}$, depending only on $r$, such that the following statements hold.

Let $0<\theta_{1}, \theta_{2}, \theta<\min \left(R_{1}, R_{2}\right)$ be numbers such that

$$
\begin{equation*}
\theta_{1} \leq \gamma_{r} \frac{\mu}{M_{1}^{r-1} M_{2}}, \theta \leq \frac{1}{2} \gamma_{r} \frac{\mu}{M_{1}^{r-1}} \theta_{1}, \theta_{2}<\frac{1}{2 s M_{3}} \theta \tag{2}
\end{equation*}
$$

Then we have:
(i) For fixed $\zeta \in B_{s}\left(0, R_{2}\right)$, the map $z \mapsto f(z, \zeta)$ is injective on the ball $B_{r}\left(0, \theta_{1}\right) \subset \mathbb{C}^{r}$.
(ii) For fixed $\zeta \in B_{s}\left(0, R_{2}\right)$, the image under $f$ of the ball $B_{r}\left(0, \theta_{1}\right)$ contains the ball of radius $\theta$ centered at $f(0, \zeta)$ in $\mathbb{C}^{r}: f\left(B_{r}\left(0, \theta_{1}\right), \zeta\right) \supset$ $B_{r}(f(0, \zeta), \theta)$.
(iii) If $|\zeta|<\theta_{2}$, there is a unique point $z=z(\zeta) \in B_{r}\left(0, \theta_{1}\right)$ for which $f(z, \zeta)=0$ and the map $\zeta \mapsto z(\zeta)$ is holomorphic. In particular, if $X=\left\{(z, \zeta) \in B_{r}\left(0, R_{1}\right) \times B_{s}\left(0, R_{2}\right) \mid f(z, \zeta)=0\right\}$ then $X \cap B_{r}\left(0, \theta_{1}\right) \times$ $B_{s}\left(0, \theta_{2}\right)=\left\{(z(\zeta), \zeta)| | \zeta \mid<\theta_{2}\right\}$ and this intersection is connected.

Proof. - During the course of this proof, we shall denote by $K$ a constant, not necessarily the same at each occurrence, which depends only on $r$.

Let $F(z, \zeta)=D(0, \zeta)^{-1} \cdot(f(z, \zeta)-f(0, \zeta))$, and set $g(z, \zeta)=F(z, \zeta)-z$ $\left(z \in B_{r}\left(0, R_{1}\right), \zeta \in B_{s}\left(0, R_{2}\right)\right)$.

We have $g(0, \zeta)=0, \frac{\partial g_{i}}{\partial z_{j}}(0, \zeta)=0,1 \leq i, j \leq r,|\zeta|<R_{2}$. The entries of the matrix $D(0, \zeta)^{-1}$ are given by $\pm \frac{d_{i j}(\zeta)}{\operatorname{det} D(0, \zeta)}$, where $d_{i j}(\zeta)$ is the determinant of the $(r-1) \times(r-1)$ matrix obtained from $D(0, \zeta)$ by deleting the $i$-th row and the $j$-th column. Hence, the norm of the linear
map $D(0, \zeta)^{-1}: \mathbb{C}^{r} \rightarrow \mathbb{C}^{r}$ satisfies, because of assumption (a),

$$
\begin{equation*}
\left\|D(0, \zeta)^{-1}\right\| \leq K \cdot \frac{M_{1}^{r-1}}{\mu} \tag{3}
\end{equation*}
$$

Using assumption (b), the definition of $g$ and (3), we obtain

$$
\begin{equation*}
\left|\frac{\partial^{2} g_{i}}{\partial z_{j} \partial z_{k}}\right| \leq K \cdot \frac{M_{1}^{r-1}}{\mu} \cdot M_{2}, 1 \leq i, j, k \leq r,|z|<R_{1},|\zeta|<R_{2} \tag{4}
\end{equation*}
$$

Since $\frac{\partial g_{i}}{\partial z_{j}}(0, \zeta)=0$, this gives

$$
\begin{equation*}
\left|\frac{\partial g_{i}}{\partial z_{j}}(z, \zeta)\right| \leq r \cdot|z| \cdot \frac{K}{\mu} M_{1}^{r-1} M_{2} \tag{5}
\end{equation*}
$$

so that, if $|z| \leq \theta_{1}$ and $\theta_{1}$ satisfies (2), we have

$$
\left|\frac{\partial g_{i}}{\partial z_{j}}\right| \leq \frac{1}{2 r},|z| \leq \theta_{1}
$$

Consequently

$$
\begin{equation*}
|g(z, \zeta)-g(w, \zeta)| \leq \frac{1}{2}|z-w| \text { for }|z|,|w|<\theta_{1} \tag{6}
\end{equation*}
$$

This implies that, for $|z|,|w| \leq \theta_{1}$, we have

$$
|F(z, \zeta)-F(w, \zeta)| \geq|z-w|-|g(z, \zeta)-g(w, \zeta)| \geq \frac{1}{2}|z-w|
$$

so that, for fixed $\zeta, z \mapsto F(z, \zeta)$ is injective on $B_{r}\left(0, \theta_{1}\right)$, and hence so is $z \mapsto f(z, \zeta)$. This proves (i).

To prove (ii), we construct the inverse of $F$ by the standard iteration scheme. Set $\varphi_{0}(w, \zeta)=0 \in \mathbb{C}^{r},|w|<\frac{1}{2} \theta_{1}$, and define $\varphi_{\nu}(w, \zeta)$ for $\nu \geq 1$, $|w|<\frac{1}{2} \theta_{1}$, inductively by

$$
\varphi_{\nu}(w, \zeta)=w-g\left(\varphi_{\nu-1}(w, \zeta), \zeta\right)
$$

Now, (6) with $w=0$ implies that $|g(z, \zeta)| \leq \frac{1}{2}|z|$ for $|z|<\theta_{1}$, so that, if $\left|\varphi_{\nu-1}(w, \zeta)\right|<\theta_{1}$, we have $\left|\varphi_{\nu}(w, \zeta)\right| \leq|w|+\frac{1}{2} \theta_{1}<\theta_{1}$ for $|w|<\frac{1}{2} \theta_{1}$. Thus $\varphi_{\nu}$ is well-defined and maps $B_{r}\left(0, \frac{1}{2} \theta_{1}\right)$ into $B_{r}\left(0, \theta_{1}\right)$. Further

$$
\begin{aligned}
\left|\varphi_{\nu+1}(w, \zeta)-\varphi_{\nu}(w, \zeta)\right| & =\left|g\left(\varphi_{\nu}(w, \zeta), \zeta\right)-g\left(\varphi_{\nu-1}(w, \zeta), \zeta\right)\right| \\
& \leq \frac{1}{2}\left|\varphi_{\nu}(w, \zeta)-\varphi_{\nu-1}(w, \zeta)\right| \text { by }(6)
\end{aligned}
$$

so that, since $\left|\varphi_{1}(w, \zeta)-\varphi_{0}(w, \zeta)\right|=|w|<\frac{1}{2} \theta_{1}$, we have $\mid \varphi_{\nu+1}(w, \zeta)-$ $\varphi_{\nu}(w, \zeta) \mid \leq 2^{-\nu-1} \theta_{1}$ for $\nu \geq 0$. It follows that $\lim _{\nu \rightarrow \infty} \varphi_{\nu}(w, \zeta)=\varphi(w, \zeta)$ exists uniformly on $B_{r}\left(0, \frac{1}{2} \theta_{1}\right) \times B_{s}\left(0, R_{2}\right)$ so that $\varphi$ is holomorphic and satisfies

$$
\varphi(w, \zeta)=w-g(\varphi(w, \zeta), \zeta),|\varphi(w, \zeta)| \leq \theta_{1} \text { for }|w|<\frac{1}{2} \theta_{1}
$$

However, this implies, because of (6), that $|\varphi(w, \zeta)|<\theta_{1}$ for $|w|<\frac{1}{2} \theta_{1}$. Moreover, the equation $\varphi(w, \zeta)=w-g(\varphi(w, \zeta), \zeta)$ can be written

$$
\begin{equation*}
F(\varphi(w, \zeta), \zeta)=w, f(\varphi(w, \zeta), \zeta)=f(0, \zeta)+D(0, \zeta) w,|w|<\frac{1}{2} \theta_{1} \tag{7}
\end{equation*}
$$

In particular, for fixed $\zeta$ with $|\zeta|<R_{2}, F\left(B_{r}\left(0, \theta_{1}\right), \zeta\right) \supset B_{r}\left(0, \frac{1}{2} \theta_{1}\right)$ and $f\left(B_{r}\left(0, \theta_{1}\right), \zeta\right)-f(0, \zeta)$ contains the image of $B_{r}\left(0, \frac{1}{2} \theta_{1}\right)$ under the linear $\operatorname{map} D(0, \zeta)$. Since $\left\|D(0, \zeta)^{-1}\right\| \leq \frac{K}{\mu} M_{1}^{r-1}$ (by (3)), we have

$$
D(0, \zeta) B_{r}\left(0, \frac{1}{2} \theta_{1}\right) \supset B_{r}\left(0, \frac{\mu}{2 K} \frac{\theta_{1}}{M_{1}^{r-1}}\right) \supset B_{r}(0, \theta)
$$

by (2). This proves (ii).
To prove (iii), we remark that $|f(0, \zeta)|=|f(0, \zeta)-f(0,0)| \leq s M_{3} \theta_{2}$ for $|\zeta| \leq \theta_{2}$ (by assumption (c)). If $s M_{3} \theta_{2} \leq \frac{1}{2} \theta$, then 0 lies in $B_{r}(f(0, \zeta), \theta)$, hence in the image of $B_{r}\left(0, \theta_{1}\right)$ under $z \mapsto f(z, \zeta)$. Since by (2) and (3),

$$
\begin{aligned}
\left|D(0, \zeta)^{-1} f(0, \zeta)\right| & \leq \frac{K}{\mu} M_{1}^{r-1} s M_{3} \theta_{2} \\
& <\frac{1}{2} \frac{K}{\mu} M_{1}^{r-1} \theta \leq \frac{1}{2} \theta_{1}
\end{aligned}
$$

the point $-D(0, \zeta)^{-1} f(0, \zeta)$ lies in $B_{r}\left(0, \frac{1}{2} \theta_{1}\right)$, and the point $z(\zeta) \in$ $B_{r}\left(0, \frac{1}{2} \theta_{1}\right)$ with $f(z, \zeta)=0$, unique by (i), is given by

$$
z(\zeta)=\varphi\left(-D(0, \zeta)^{-1} f(0, \zeta), \zeta\right)
$$

This proves (iii).
Remark. - Lemma 2.1 remains valid if we replace $\mathbb{C}^{r}, \mathbb{C}^{s}$ by $\mathbb{R}^{r}, \mathbb{R}^{s}$ and holomorphic maps by smooth (or real-analytic) maps. The proof above implies that the function $\varphi$ is continuous. However, the standard inverse
function theorem implies that it is smooth (or real-analytic). In particular, the connectedness statement in Part (iii) implies that there is $\rho_{1}>0$ (depending only on $n, C_{1}, D$ ) such that for $P \in W, Q_{\rho} \cap Z(P) \subset V^{0}(P)$ for $0<\rho \leq \rho_{1}$.

Recall that $\mathcal{W}$ is the space of all $P=\left(p_{1}, \ldots, p_{r}\right)$ with $p_{j} \in \mathcal{H}^{D}$, $P(0)=0,\|P\| \leq C_{1}, J_{P}(0)=1$ and

$$
\left|J_{P}(z)\right|^{2}+\left|p_{1}(z)\right|^{2}+\cdots+\left|p_{r}(z)\right|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N}, z \in \mathbb{C}^{n}
$$

In particular, $J_{P} \neq 0$ on $\mathcal{Z}(P)$, so that $\mathcal{Z}(P)$ is smooth and $\mathcal{V}(P)=\mathcal{V}^{0}(P)$ is the connected component of $\mathcal{Z}(P)$ through the origin.

Lemma 2.2. - There exist constants $c_{1}>0, N_{1} \geq 1$ depending only on $C_{1}, D, n, c_{0}, N$ such that if $P \in \mathcal{W}$ and $X$ is any connected component of $\mathcal{Z}(P)$, then, for any $z \in X$, we have

$$
\mathcal{Z}(P) \cap B\left(z, c_{1}\left(1+|z|^{2}\right)^{-N_{1}}\right) \subset X
$$

Proof. - In this argument, we denote by $(\gamma, m),\left(\gamma_{j}, m_{j}\right),\left(\gamma^{\prime}, m^{\prime}\right)$, etc. constants depending only on $C_{1}, D, n, c_{0}, N$.

If $z \in \mathcal{Z}(P)$, we have $P(z)=0$ so that $\left|J_{P}(z)\right|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N}$. Since $p_{j} \in \mathcal{H}^{D},\|P\| \leq C_{1}$, we have $\left|J_{P}(w)\right|^{2} \geq \frac{1}{2} c_{0}\left(1+|z|^{2}\right)^{-N}$ if $w \in \mathbb{C}^{n}$, $|w-z| \leq \gamma\left(1+|z|^{2}\right)^{-m}$ (for suitable $\gamma>0, m \geq 1$ ). We apply Lemma 2.1 with $R_{1}=R_{2}=\gamma\left(1+|z|^{2}\right)^{-m}, s=n-r, \zeta_{j}=z_{r+j}(1 \leq j \leq s)$ and $f_{j}=P_{j}$ (replacing the origin by $z$ ). We have $\left|\frac{\partial f_{j}}{\partial z_{k}}\right|,\left|\frac{\partial^{2} f_{i}}{\partial z_{j} \partial z_{k}}\right|,\left|\frac{\partial f_{j}}{\partial \zeta_{\ell}}\right| \leq \gamma^{\prime}\left(1+|z|^{2}\right)^{m^{\prime}}$ on $B(z, 1)(1 \leq i, j, k \leq r, 1 \leq \ell \leq s)$, and we can apply Lemma 2.1 with $\theta_{1}=\gamma_{1}\left(1+|z|^{2}\right)^{-m_{1}}, \theta_{2}=\gamma_{2}\left(1+|z|^{2}\right)^{-m_{2}}, \theta=\gamma_{3}\left(1+|z|^{2}\right)^{-m_{3}}$. Part (iii) of the lemma implies that

$$
\left\{w \in \mathcal{Z}(P)\left|\sum_{1}^{r}\right| w_{j}-\left.z_{j}\right|^{2}<\theta_{1}^{2}, \sum_{r+1}^{n}\left|w_{k}-z_{k}\right|^{2}<\theta_{2}^{2}\right\}
$$

is connected, hence contained in $X$ since $z \in X$. The lemma follows.
Definition 2.3.- Let $E \subset \mathbb{C}^{n}$ and let $\gamma>0, m>0$. We set

$$
T_{\gamma, m}(E)=\bigcup_{z \in E} B\left(z, \gamma\left(1+|z|^{2}\right)^{-m}\right)
$$

Lemma 2.4. - There exist constants $\gamma, m>0$ depending only on the constants defining $\mathcal{W}$ (viz. $C_{1}, n, D, c_{0}, N$ ) such that

If $P \in \mathcal{W}$ and $\mathcal{V}_{1}, \ldots, \mathcal{V}_{t}$ are the connected components of $\mathcal{Z}(P)$, then $T_{\gamma, m}\left(\mathcal{V}_{i}\right) \cap T_{\gamma, m}\left(\mathcal{V}_{j}\right)=\emptyset$ if $i \neq j$.

This follows easily from Lemma 2.2.
Lemma 2.5. - Let $\gamma_{0}, m_{0}>0$ be given. There exist constants $\gamma, m>0$ depending only on $\gamma_{0}, m_{0}$ and the constants defining $\mathcal{W}$ such that
(8) If $P \in \mathcal{W}$ and $z \notin T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$, then,

$$
|P(z)|^{2}=\left|p_{1}(z)\right|^{2}+\cdots+\left|p_{r}(z)\right|^{2} \geq \gamma\left(1+|z|^{2}\right)^{-m}
$$

Proof. - For any $z \in \mathbb{C}^{n}$, we have $\left|J_{P}(z)\right|^{2}+|P(z)|^{2} \geq$ $c_{0}\left(1+|z|^{2}\right)^{-N}$. Let $w \in \mathbb{C}^{n} \backslash T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$ and suppose that

$$
\begin{equation*}
|P(w)|^{2} \leq \frac{1}{2} c_{0}\left(1+|w|^{2}\right)^{-N} \tag{9}
\end{equation*}
$$

Then $\left|J_{P}(w)\right|^{2} \geq \frac{1}{2} c_{0}\left(1+|w|^{2}\right)^{-N}$. Now, since $P \in \mathcal{W}$ and $w \notin T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$, there exist constants $\gamma^{\prime}, m^{\prime}>0$ depending only on $\gamma_{0}, m_{0}$ and the constants in $\mathcal{W}$ such that
(a) $\left|J_{P}(z)\right|^{2} \geq \frac{1}{4} c_{0}\left(1+|w|^{2}\right)^{-N}$ if $z \in B\left(w, \gamma^{\prime}\left(1+|w|^{2}\right)^{-m^{\prime}}\right)$
(b) $B\left(w, \gamma^{\prime}\left(1+|w|^{2}\right)^{-m^{\prime}}\right) \cap \mathcal{Z}(P)=\emptyset$.

Lemma 2.1 implies that there are constants $\gamma^{\prime \prime}, m^{\prime \prime}>0$ (depending only on $\gamma^{\prime}, m^{\prime}$ and the constants in $\left.\mathcal{W}\right)$ such that the image of $B\left(w, \gamma^{\prime}(1+\right.$ $\left.|w|^{2}\right)^{-m^{\prime}}$ ) under the map $P: \mathbb{C}^{n} \rightarrow \mathbb{C}^{r}$ contains the ball $B(P(w), \theta)$ where $\theta=\gamma^{\prime \prime}\left(1+|w|^{2}\right)^{-m^{\prime \prime}}$. By (b) above, this image does not contain $0 \in \mathbb{C}^{r}$, so that

$$
\begin{equation*}
|P(w)| \geq \gamma^{\prime \prime}\left(1+|w|^{2}\right)^{-m^{\prime \prime}} \tag{10}
\end{equation*}
$$

Thus if (9) holds, so does (10), which proves the lemma.
Lemma 2.6. - There exist constants $\gamma_{0}, m_{0}>0$ depending only on the constants in $\mathcal{W}$ such that the following holds.

Let $P \in \mathcal{W}$ and $\Omega=T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$. Then, there exists a holomorphic map $\pi: \Omega \rightarrow \mathbb{C}^{n}$ with the following properties:
(i) $\pi(\Omega) \subset \mathcal{Z}(P), \pi(z)=z$ if $z \in \mathcal{Z}(P)$;
(ii) If $z=\left(z_{1}, \ldots, z_{n}\right)$ and $\pi(z)=\left(\pi_{1}(z), \ldots, \pi_{n}(z)\right)$, we have $\pi_{k}(z)=$ $z_{k}$ for $r<k \leq n$;
(iii) $|\pi(z)-z|<1$ for any $z \in \Omega$.

Proof. - Since $\left|J_{P}(w)\right|^{2} \geq c_{0}\left(1+|w|^{2}\right)^{-N}$ if $w \in \mathcal{Z}(P)$, we can choose $\gamma_{1}, m_{1}>0$, depending only on the constants in $\mathcal{W}$, such that

$$
\left|J_{P}(z)\right|^{2} \geq \frac{1}{2} c_{0}\left(1+|z|^{2}\right)^{-N} \text { for } z \in T_{\gamma_{1}, m_{1}}(\mathcal{Z}(P))
$$

By Lemma 2.1, there exist $\left(\gamma_{2}, m_{2}\right),\left(\gamma_{3}, m_{3}\right)$ (depending only on the constants in $\mathcal{W}$ ) such that if $w \in \mathcal{Z}(P)$ and we set

$$
\Delta(w)=\left\{\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n}\left|\sum_{1}^{r}\right| z_{j}-\left.w_{j}\right|^{2}<\theta_{1}^{2}, \sum_{r+1}^{n}\left|z_{k}-w_{k}\right|^{2}<\theta_{2}^{2}\right\}
$$

with $\theta_{1}=\gamma_{2}\left(1+|w|^{2}\right)^{-m_{2}}, \theta_{2}=\gamma_{3}\left(1+|w|^{2}\right)^{-m_{3}}$, then

$$
\begin{array}{r}
\Delta(w) \cap \mathcal{Z}(P)=\left\{\left(\varphi_{1}^{(w)}\left(z_{r+1}, \ldots, z_{n}\right), \ldots, \varphi_{r}^{(w)}\left(z_{r+1}, \ldots, z_{n}\right), z_{r+1}, \ldots, z_{n}\right) \mid\right. \\
\left.\sum_{r+1}^{n}\left|z_{k}-w_{k}\right|^{2}<\theta_{2}^{2}\right\}
\end{array}
$$

where $\varphi^{(w)}=\left(\varphi_{1}^{(w)}, \ldots, \varphi_{r}^{(w)}\right)$ is a holomorphic map of $B_{n-r}\left(w^{\prime \prime}, \theta_{2}\right) \subset$ $\mathbb{C}^{n-r}$ into $B_{r}\left(w^{\prime}, \theta_{1}\right) \subset \mathbb{C}^{r} ;$ here we have written $w^{\prime}=\left(w_{1}, \ldots, w_{r}\right), w^{\prime \prime}=$ $\left(w_{r+1}, \ldots, w_{n}\right)$. Moreover, if $\tilde{w} \in \mathcal{Z}(P)$ and $\Delta(w) \cap \Delta(\tilde{w}) \neq \emptyset$, Part (i) of Lemma 2.1 implies that $\left(\varphi^{(w)}, z_{r+1}, \ldots, z_{n}\right)=\left(\varphi^{(\tilde{w})}, z_{r+1}, \ldots, z_{n}\right)$ on $\Delta(w) \cap \Delta(\tilde{w})$.

Thus, we may define a holomorphic map

$$
\pi: \bigcup_{w \in \mathcal{Z}(P)} \Delta(w) \rightarrow \mathcal{Z}(P)
$$

by setting $\pi\left(z_{1}, \ldots, z_{n}\right)=\left(\varphi_{1}^{(w)}\left(z_{r+1}, \ldots, z_{n}\right), \ldots, \varphi_{r}^{(w)}\left(z_{r+1}, \ldots, z_{n}\right), z_{r+1}\right.$, $\ldots, z_{n}$ ) for $z \in \Delta(w)$. If we choose $\gamma_{0}, m_{0}$ such that $\gamma_{0}<1$ and $T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P)) \subset \bigcup_{w \in \mathcal{Z}(P)} \Delta(w)$, then $\pi \mid T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$ has the properties stated.

Lemma 2.7. - Let $0<\gamma \leq 1$ and $m$ be an integer $\geq 1$. Let $U=\left\{z \in \mathbb{C}^{n}| | \operatorname{Re} z_{j}\left|<1,\left|\operatorname{Im} z_{j}\right|<1, j=1, \ldots, n\right\}\right.$.

There exists a constant $C_{*}$ depending only on $n$ and a sequence $\left\{U_{\nu}\right\}_{\nu \geq 0}$ of open sets of the form

$$
U_{\nu}=a_{\nu}+\rho_{\nu} U, \quad a_{\nu} \in \mathbb{C}^{n}, \quad \rho_{\nu}>0
$$

with the following properties:
(a) $\bigcup_{\nu} U_{\nu}=\mathbb{C}^{n}$.
(b) $C_{*}^{-m} \gamma\left(1+\left|a_{\nu}\right|^{2}\right)^{-m} \leq \rho_{\nu} \leq C_{*}^{m} \gamma\left(1+\left|a_{\nu}\right|^{2}\right)^{-m}, \nu \geq 0$.
(c) If, for $t \geq 1$, we set $U_{\nu}(t)=a_{\nu}+t \rho_{\nu} U$ when $U_{\nu}=a_{\nu}+\rho_{\nu} U$, then, for any $t \geq 1$ and $\nu \geq 0$, the number of $\mu$ such that $U_{\nu}(t) \cap U_{\mu}(t) \neq \emptyset$ is $\leq C_{*}^{t m}$.

Proof. - Let $N_{0}=[2 / \gamma]$ be the integral part of $2 / \gamma$. We have $\frac{1}{\gamma} \leq \frac{2}{\gamma}-1<N_{0} \leq \frac{2}{\gamma}$. For $k \geq 0$, let $N_{k}=2^{2 k m} N_{0}$. Let $Q=$ $\left\{z \in \mathbb{C}^{n}| | \operatorname{Re} z_{j}\left|\leq 1,\left|\operatorname{Im} z_{j}\right| \leq 1, j=1, \ldots, n\right\}\right.$ be the closed unit cube in $\mathbb{C}^{n}$ and, for $k \geq 0$, let $Q^{(k)}$ be the closed cube of side $2^{k+1}$ in $\mathbb{C}^{n}: Q^{(k)}=2^{k} Q=\left\{z \in \mathbb{C}^{n}| | \operatorname{Re} z_{j}\left|\leq 2^{k},\left|\operatorname{Im} z_{j}\right| \leq 2^{k}\right\} ;\right.$ we set $Q^{(-1)}=\emptyset$.

Divide the interval $\left\{x \in \mathbb{R} \mid-2^{k} \leq x \leq 2^{k}\right\}$ into equal intervals $I_{\ell}^{(k)}$, $\ell=1, \ldots, L$, of length $N_{k}^{-1}$ (so that $L=2^{k+1} N_{k}$ ). Consider the cubes of side $N_{k}^{-1}$ in $\mathbb{C}^{n}$ obtained by taking products of the $I_{\ell}^{(k)}: \operatorname{Re} z_{j} \in I_{\ell_{j}}^{(k)}$, $\operatorname{Im} z_{j} \in I_{\ell_{j}^{\prime}}^{(k)}, 1 \leq \ell_{j}, \ell_{j}^{\prime} \leq L$.

Let $Q_{\alpha}^{(k)}, \alpha=1, \ldots, \alpha_{k}$ be a list of those cubes of side $N_{k}^{-1}$ described above which are not contained in $Q^{(k-1)}$. Then $\bigcup_{k} Q_{\alpha}^{(k)}=Q^{(k)} \backslash$ int $\left(Q^{(k-1)}\right)$.

Now, if $z \in Q=Q^{(0)}, 1 \leq 1+|z|^{2} \leq 2 n+1$, while $1+2^{2(k-1)} \leq$ $1+|z|^{2} \leq(2 n+1) 2^{2 k}$ for $z \in Q^{(k)} \backslash$ int $\left(Q^{(k-1)}\right), k \geq 1$. In particular, if $Q_{\alpha}^{(k)}=a+\rho Q$, then

$$
2^{2 k-2} \leq 1+|a|^{2} \leq(2 n+1) 2^{2 k}, 2 \rho=N_{k}^{-1}, \frac{1}{\gamma} 2^{2 k m} \leq N_{k} \leq \frac{1}{\gamma} 2^{2 k m+1}
$$

It follows that if $Q_{\alpha}^{(k)}=a+\rho Q$, then

$$
\begin{equation*}
2^{-2 m-2} \gamma\left(1+|a|^{2}\right)^{-m} \leq \rho \leq(2 n+1)^{m} \gamma\left(1+|a|^{2}\right)^{-m} \tag{11}
\end{equation*}
$$

Let $0<\epsilon<1$, and set $U_{\alpha}^{(k)}=U_{\alpha}^{(k)}(1+\epsilon)=a+(1+\epsilon) \rho U$ if $Q_{\alpha}^{(k)}=a+\rho Q$; then $Q_{\alpha}^{(k)} \subset U_{\alpha}^{(k)}$; if $\left\{U_{\nu}\right\}_{\nu \geq 0}$ is an enumeration of $\left\{U_{\alpha}^{(k)}\right.$, $\left.k \geq 0, \alpha=1, \ldots, \alpha_{k}\right\}$, then, since $\bigcup_{k \geq 0} \bigcup_{1 \leq \alpha \leq \alpha_{k}} Q_{\alpha}^{(k)}=\mathbb{C}^{n}$, it follows that $\left\{U_{\nu}\right\}_{\nu \geq 0}$ covers $\mathbb{C}^{n}$ which is (a). Moreover, (11) implies that $\left\{U_{\nu}\right\}_{\nu \geq 0}$ satisfies (b) with $C_{*}=4(n+2)$ (since $0<\epsilon<1$ ).

It remains to prove that $\left\{U_{\nu}\right\}_{\nu \geq 0}$ satisfies (c).
We first remark that if $t \geq 1$, and $Q_{\alpha}^{(k)}(t) \cap Q_{\beta}^{(\ell)}(t) \neq \emptyset$ (where $Q_{\alpha}^{(k)}(t)=a+t \rho Q$ if $\left.Q_{\alpha}^{(k)}=a+\rho Q\right)$, then we must have $|k-\ell| \leq t+2$; in fact, if $\ell \geq k+t+2$, then $Q_{\alpha}^{(k)}(t)$ is contained in the cube $\left|\operatorname{Re} z_{j}\right|$, $\left|\operatorname{Im} z_{j}\right| \leq 2^{k}+t-1$ (since any cube $Q_{\alpha}^{(k)}$ has side $\leq 1$ ) while $Q_{\beta}^{(\ell)}(t)$ lies
outside $\left|\operatorname{Re} z_{j}\right|,\left|\operatorname{Im} z_{j}\right|<2^{\ell-1}-t+1$. Hence, if $2^{\ell-1}-t+1>2^{k}+t-1$, in particular if $\ell \geq k+t+2$, we must have $Q_{\alpha}^{(k)}(t) \cap Q_{\beta}^{(\ell)}(t)=\emptyset$.

Fix $k, \alpha$. From the remark above, it follows that if $Q_{\alpha}^{(k)}(t) \cap Q_{\beta}^{(\ell)}(t) \neq \emptyset$, then $Q_{\beta}^{(\ell)} \subset Q_{\alpha}^{(k)}\left(A_{1}^{t m}\right)$ where $A_{1}$ is an absolute constant. Moreover, the side of $Q_{\beta}^{(\ell)}$ is $\geq A_{2}^{-t m} \rho$, where $\rho$ is the side of $Q_{\alpha}^{(k)}$, and $A_{2}$ is also an absolute constant. The number of cubes of side $\geq A_{2}^{-t m} \rho$ contained in a cube of side $A_{1}^{t m} \rho$ and no two of which have interior points in common is $\leq C_{*}^{t m}$. Replacing $t$ by $(1+\epsilon) t \leq 2 t$ we obtain (c).

This proves the lemma.
Lemma 2.8. - Let $0<\gamma_{0} \leq 1$ and let $m_{0}$ be an integer $\geq 1$. There exists a sequence $C_{k}=C_{k}\left(n, \gamma_{0}, m_{0}\right)$ of constants, $k=0,1, \ldots\left(C_{k}\right.$ depends only on $k, n, \gamma_{0}, m_{0}$ ) such that the following holds.

Let $E \subset \mathbb{C}^{n}$ be any subset. For a function $\varphi$, write $\operatorname{supp}(\varphi)$ for its support.

We can find a $C^{\infty}$ function $\varphi: \mathbb{C}^{n} \rightarrow \mathbb{R}$ with the following properties:
(a) $0 \leq \varphi \leq 1, \varphi \mid E \equiv 1, \operatorname{supp}(\varphi) \subset T_{\gamma_{0}, m_{0}}(E)$.
(b) If $D^{(k)}=\frac{\partial^{k}}{\partial z_{1}^{\alpha_{1}} \cdots \partial z_{n}^{\alpha_{n}} \partial \bar{z}_{1}^{\beta_{1}} \cdots \partial \bar{z}_{n}^{\beta_{n}}}, \alpha_{1}+\cdots+\alpha_{n}+\beta_{1}+\cdots+\beta_{n}=k$, denotes any differentiation of order $k \geq 0$, we have

$$
\left|D^{(k)} \varphi(z)\right| \leq C_{k}\left(n, \gamma_{0}, m_{0}\right)\left(1+|z|^{2}\right)^{k m_{0}}, \quad \forall z \in \mathbb{C}^{n}
$$

Proof. - Let $C_{*}$ be the constant in parts (b), (c) of Lemma 2.7. Set $m=m_{0}$, and choose a constant $\gamma>0$ (depending only on $C_{*}, n, \gamma_{0}, m$ ) with the following property:
(12) Let $Q(a, \rho)=a+\rho Q\left(Q=\right.$ closed unit cube in $\left.\mathbb{C}^{n}\right)$. If $b \in Q(a, \rho)$ and $\rho \leq 2 C_{*}^{m} \gamma\left(1+|a|^{2}\right)^{-m}$, then $Q(a, \rho) \subset B\left(b, \rho_{1}\right)$ with $\rho_{1}=\gamma_{0}\left(1+|b|^{2}\right)^{-m}$.

Let $U$ be the open unit cube in $\mathbb{C}^{n}$ and $\left\{U_{\nu}\right\}_{\nu \geq 0}, U_{\nu}=a_{\nu}+\rho_{\nu} U$, be an open covering of $\mathbb{C}^{n}$ with the properties given in Lemma 2.7. We set $V_{\nu}=a_{\nu}+2 \rho_{\nu} U, \nu \geq 0$.

Let $\psi \in C^{\infty}\left(\mathbb{C}^{n}\right), 0 \leq \psi \leq 1$ be such that $\psi=1$ on $Q$ and $\operatorname{supp}(\psi) \subset 2 U$.

Define $\psi_{\nu}(z)=\psi\left(\frac{z-a_{\nu}}{\rho_{\nu}}\right)$. Then $\psi_{\nu}=1$ on $Q\left(a_{\nu}, \rho_{\nu}\right)$ and
$\operatorname{supp}\left(\psi_{\nu}\right) \subset V_{\nu}$. Moreover, if $D^{(k)}$ is a differentiation of order $k$, we have

$$
\left|D^{(k)} \psi_{\nu}(z)\right| \leq \rho_{\nu}^{-k} \sup _{\mathbb{C}^{n}}\left|D^{(k)} \psi\right|, z \in \mathbb{C}^{n}
$$

Since $\operatorname{supp}\left(\psi_{\nu}\right) \subset V_{\nu}$ and $C_{*}^{-m} \gamma\left(1+\left|a_{\nu}\right|^{2}\right)^{-m} \leq \rho_{\nu} \leq C_{*}^{m} \gamma\left(1+\left|a_{\nu}\right|^{2}\right)^{-m}$ (property (b) in Lemma 2.7), this implies that

$$
\left|D^{(k)} \psi_{\nu}(z)\right| \leq \bar{C}_{k}(\gamma, m, n)\left(1+|z|^{2}\right)^{k m}, z \in \mathbb{C}^{n}
$$

Let $S$ be the set of $\nu \geq 0$ such that $V_{\nu} \cap E \neq \emptyset$, and define $\Psi(z)=\sum_{\nu=0}^{\infty} \psi_{\nu}(z), \varphi(z)=\frac{1}{\Psi(z)} \sum_{\nu \in S} \psi_{\nu}(z)$.
$\Psi$ is $C^{\infty}$ since $\left\{V_{\nu}\right\}$ is locally finite, and $\Psi \geq 1$ on $\mathbb{C}^{n}$ since $\psi_{\nu}=1$ on $U_{\nu}$ and $\cup U_{\nu}=\mathbb{C}^{n}$.

If $z \in E$ and $\nu \notin S$, then $z \notin V_{\nu}$, so that $z \notin \operatorname{supp}\left(\psi_{\nu}\right)$. Hence, for $z \in E, \sum_{\nu \in S} \psi_{\nu}(z)=\Psi(z)$, so that $\psi(z)=1$.

Further, if $z \in \operatorname{supp}(\varphi)$, then $z \in \operatorname{supp}\left(\psi_{\nu}\right)$ for some $\nu \in S$ (since $\left\{V_{\mu}\right\}$ is locally finite) so that $z \in V_{\nu}$ for some $\nu$ with $V_{\nu} \cap E \neq \emptyset$. The choice (12) above of $\gamma$ shows that we then have $V_{\nu} \subset B\left(b, \rho_{1}\right), b \in E$, $\rho_{1}=\gamma_{0}\left(1+|b|^{2}\right)^{-m}$, so that $V_{\nu} \subset T_{\gamma_{0}, m_{0}}(E)$. Thus $\operatorname{supp}(\varphi) \subset T_{\gamma_{0}, m_{0}}(E)$.

We have $\left|D^{(k)} \psi_{\nu}(z)\right| \leq \bar{C}_{k}(\gamma, m, n)\left(1+|z|^{2}\right)^{k m}$; by property (c) in Lemma 2.7, we therefore have

$$
\left|D^{(k)} \Psi(z)\right| \leq C_{*}^{2 m} \bar{C}_{k}(\gamma, m, n)\left(1+|z|^{2}\right)^{k m}
$$

and the same argument, applied to the numerator, gives

$$
\left|D^{(k)} \sum_{\nu \in S} \psi_{\nu}(z)\right| \leq C_{*}^{2 m} \bar{C}_{k}(\gamma, m, n)\left(1+|z|^{2}\right)^{k m}
$$

These inequalities, and the fact that $\Psi \geq 1$ on $\mathbb{C}^{n}$ imply that

$$
\left|D^{(k)} \varphi(z)\right| \leq C_{k}^{\prime}(\gamma, m, n)\left(1+|z|^{2}\right)^{k m}
$$

thus proving the lemma.

## 3. Separating components of a smooth algebraic variety.

Let $n \geq 2,1 \leq r \leq n-1, D \geq 1, C_{1}>0, c_{0}>0$ and $N \geq 1$ be given. They define the space $\mathcal{W}$ introduced in §2:
$\mathcal{W}$ is the space of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right), \quad p_{j} \in \mathcal{H}^{D}, \cdot P(0)=0$, $\|P\| \leq C_{1}, \quad J_{P}(0)=1$ such that

$$
\left|J_{P}(z)\right|^{2}+|P(z)|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N}
$$

(where $J_{P}(z)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}\right)_{1 \leq j, k \leq r}$ ).
The aim of this section is the proof of the following theorem.
The formalism of the Koszul complex is also developed in [KT] and applied to a closely related problem.

Theorem 3.1. - There exist constants $C^{\prime}, D^{\prime}>0$ depending only on $n, D, C_{1}, c_{0}, N$ such that the following holds.

Let $P \in \mathcal{W}$ and let $X$ be any connected component of $\mathcal{Z}(P)=\{z \in$ $\left.\mathbb{C}^{n} \mid P(z)=0\right\}$.

We can find a polynomial $F\left(z_{1}, \ldots, z_{n}\right)$ with the following properties:
(i) $\operatorname{deg} F \leq D^{\prime},\|F\| \leq C^{\prime}$.
(ii) $F|X=1, F| \mathcal{Z}(P) \backslash X=0$.

We begin with some algebraic preliminaries.
Let $r \geq 1$ be an integer and let $E_{p}=\stackrel{p}{\wedge} \mathbb{C}^{r}$ be the $p$-th exterior power of $\mathbb{C}^{r}\left(E_{0}=\mathbb{C}, E_{1}=\mathbb{C}^{r}\right.$ and $E_{p}=0$ if $\left.p>r\right)$.

Let $e_{1}, \ldots, e_{r}$ be the standard basis of $\mathbb{C}^{r}\left[e_{i}=(0, \ldots, 1, \ldots, 0)\right.$ with 1 in the $i$-th place].

If $I$ is an increasing $p$-tuple $I=\left(i_{1}, \ldots, i_{p}\right)$ with $1 \leq i_{1}<\cdots<i_{p} \leq r$, set $e_{I}=e_{i_{1}} \wedge \cdots \wedge e_{i_{p}} \in E_{p}$. (If $I$ is not in increasing order, we use the convention that $e_{I}$ is alternating in the indices $i_{1}, \ldots, i_{p}$.) The $\left\{e_{I}\right\}, I$ increasing, form a basis of $E_{p}$.
(1) Let $x_{1}, \ldots, x_{r} \in \mathbb{C}$. We define a $\operatorname{map} \kappa=\kappa_{p}=\kappa_{p}(x)$,

$$
\kappa: E_{p} \rightarrow E_{p-1}
$$

as follows: $\kappa$ is the $\mathbb{C}$-linear map such that

$$
\begin{equation*}
\kappa\left(e_{I}\right)=\sum_{k=1}^{p}(-1)^{k-1} x_{i_{k}} e_{I \backslash\left\{i_{k}\right\}}, e_{I \backslash\left\{i_{k}\right\}}=e_{i_{1}} \wedge \cdots \wedge \hat{e}_{i_{k}} \wedge \cdots \wedge e_{i_{p}} \tag{2}
\end{equation*}
$$

where the hat over $e_{i_{k}}$ indicates that it is to be deleted. On $E_{1}, \kappa: E_{1} \rightarrow \mathbb{C}$ is the map $\sum u_{i} e_{i} \mapsto \sum x_{i} u_{i}$. One checks easily that the map $\kappa \circ \kappa: E_{p+1} \rightarrow$ $E_{p-1}(p \geq 1)$ is $0:$

$$
\kappa^{2}=0
$$

so that one has an algebraic complex, usually called the Koszul complex of $\left(x_{1}, \ldots, x_{r}\right)$ :

$$
\begin{equation*}
0 \rightarrow E_{r} \xrightarrow{\kappa} E_{r-1} \rightarrow \cdots \rightarrow E_{1} \xrightarrow{\kappa} E_{0} \rightarrow 0 . \tag{3}
\end{equation*}
$$

Let $y_{1}, \ldots, y_{r}$ be any $r$ complex numbers. We define a map $h=h_{p}=$ $h_{p}(y)$ :

$$
h: E_{p} \rightarrow E_{p+1}(p \geq 0)
$$

by

$$
\begin{equation*}
h(v)=\left(\sum_{j=1}^{r} y_{j} e_{j}\right) \wedge v, v \in E_{p} \tag{4}
\end{equation*}
$$

Lemma 3.2. - Let $\xi=\sum_{j=1}^{r} x_{j} \cdot y_{j}$. We have:

$$
\begin{align*}
\kappa h(v) & =\xi \cdot v \text { if } v \in E_{0}  \tag{5}\\
(\kappa h+h \kappa)(v) & =\xi \cdot v \text { if } v \in E_{p}, p>0 \tag{6}
\end{align*}
$$

Proof. - This is a standard fact about the Koszul complex. The verification runs as follows.

If $v \in E_{0}=\mathbb{C}, h(v)=v y_{1} e_{1}+\cdots+v y_{r} e_{r}$, and $\kappa(h(v))=\sum x_{j} y_{j} \cdot v$.
Let $p>0$. It is enough to check (6) when $v=e_{I}$ for some $I=\left(i_{1}<\cdots<i_{p}\right)$. We have

$$
\begin{aligned}
h \kappa\left(e_{I}\right) & =\left(\sum_{j=1}^{r} y_{j} e_{j}\right) \wedge \sum_{k=1}^{p}(-1)^{k-1} x_{i_{k}} e_{I \backslash\left\{i_{k}\right\}} \\
& =\sum_{k=1}^{p}(-1)^{k-1} x_{i_{k}} y_{i_{k}} e_{i_{k}} \wedge e_{I \backslash\left\{i_{k}\right\}}+\sum_{j \notin I} \sum_{k=1}^{p}(-1)^{k-1} x_{i_{k}} y_{j} e_{j} \wedge e_{I \backslash\left\{i_{k}\right\}} \\
& =\left(\sum_{k=1}^{p} x_{i_{k}} y_{i_{k}}\right) e_{I}+\sum_{j \notin I} \sum_{k=1}^{p}(-1)^{k-1} x_{i_{k}} y_{j} e_{j} \wedge e_{I \backslash\left\{i_{k}\right\}}
\end{aligned}
$$

while

$$
\begin{aligned}
\kappa h\left(e_{I}\right) & =\kappa\left(\sum_{j=1}^{r} y_{j} e_{j} \wedge e_{I}\right)=\sum_{j \notin I} y_{j} \kappa\left(e_{j} \wedge e_{I}\right) \\
& =\sum_{j \notin I} x_{j} y_{j} e_{I}+\sum_{j \notin I} y_{j} \sum_{k=1}^{p}(-1)^{k} x_{i_{k}} e_{j} \wedge e_{I \backslash\left\{i_{k}\right\}}
\end{aligned}
$$

adding these two equations, we get (6) for $v=e_{I}$, thus proving the lemma.
Let $0 \leq q \leq n, 0 \leq p \leq r$. We denote by $\mathcal{A}^{q}\left(E_{p}\right)$ the space of $C^{\infty}$ forms of type $(0, q)$ on $\mathbb{C}^{r}$ with values in $E_{p}$. If $J=\left(j_{1}<\cdots<j_{q}\right)$ is an
increasing $q$-tuple of integers $j_{k}, 1 \leq j_{k} \leq n$, and we set $d \bar{z}^{J}=d \bar{z}_{j_{1}} \wedge \cdots$ $\wedge d \bar{z}_{j_{q}}$, an element $\omega \in \mathcal{A}^{q}\left(E_{p}\right)$ can be written uniquely in the form

$$
\omega=\sum_{J} v_{J} \otimes d \bar{z}^{J}, v_{J}=\sum_{I} v_{J}^{I} e_{I}, I=\left(i_{1}<\cdots<i_{p}\right)
$$

where the $v_{J}^{I}$ are $C^{\infty}$ functions on $\mathbb{C}^{n}$. We set

$$
|\omega(z)|^{2}=\sum_{J}\left|v_{J}(z)\right|^{2}=\sum_{J} \sum_{I}\left|v_{J}^{I}(z)\right|^{2}
$$

The $\bar{\partial}$-operator extends to a map (denoted again by $\bar{\partial}$ )

$$
\bar{\partial}: \mathcal{A}^{q}\left(E_{p}\right) \rightarrow \mathcal{A}^{q+1}\left(E_{p}\right), p, q \geq 0
$$

If $\omega=\sum_{J, I} v_{J}^{I} e_{I} \otimes d \bar{z}^{J}$, we have

$$
\bar{\partial} \omega=\sum_{I, J} \sum_{\nu=1}^{n} \frac{\partial v_{J}^{I}}{\partial \bar{z}_{\nu}} e_{I} \otimes d \bar{z}_{\nu} \wedge d \bar{z}^{J}
$$

If $f_{1}, \ldots, f_{r}$ are $C^{\infty}$ functions on $\mathbb{C}^{n}$, we define a map $\kappa=\kappa_{p}^{q}(f)$ : $\mathcal{A}^{q}\left(E_{p}\right) \rightarrow \mathcal{A}^{q}\left(E_{p-1}\right)$ by

$$
\begin{equation*}
\kappa\left(\sum_{I, J} v_{J}^{I} e_{I} \otimes d \bar{z}^{J}\right)(z)=\sum_{I, J} v_{J}^{I}(z) \sum_{k=1}^{p}(-1)^{k-1} f_{i_{k}}(z) e_{I \backslash\left\{i_{k}\right\}} \otimes d \bar{z}^{J} \tag{7}
\end{equation*}
$$

Note that if $f_{1}, \ldots, f_{r}$ are holomorphic on $\mathbb{C}^{n}$, then

$$
\kappa \bar{\partial}(\omega)=\bar{\partial} \kappa(\omega), \omega \in \mathcal{A}^{q}\left(E_{p}\right)
$$

In fact, in terms of bases, $\kappa(\omega)$ is given by multiplication by a matrix of holomorphic functions, and this operation commutes with $\bar{\partial}$.

If $\Omega_{0} \subset \mathbb{C}^{n}$ is open, and $g_{1}, \ldots, g_{r} \in C^{\infty}\left(\Omega_{0}\right)$, and we denote by $\mathcal{A}_{0}^{q}\left(E_{p}\right)$ the subspace of $\mathcal{A}^{q}\left(E_{p}\right)$ consisting of forms $\omega$ with $\operatorname{supp}(\omega) \subset \Omega_{0}$, then, we can define a map $h=h_{p}^{q}(g): \mathcal{A}_{0}^{q}\left(E_{p}\right) \rightarrow \mathcal{A}_{0}^{q}\left(E_{p+1}\right)$ by:

$$
\begin{equation*}
h\left(\sum_{I, J} v_{J}^{I} e_{I} \otimes d \bar{z}^{J}\right)=\sum_{J} \sum_{I} \sum_{j=1}^{r} g_{j} v_{J}^{I}\left(e_{j} \wedge e_{I}\right) \otimes d \bar{z}^{J} \tag{8}
\end{equation*}
$$

(i.e. we just operate on the coefficients of the $d \bar{z}^{J}$ ). Lemma 3.2 implies that

$$
\begin{equation*}
\kappa h=\text { multiplication by } \sum_{1}^{r} f_{j} g_{j} \text { on } \mathcal{A}_{0}^{q}\left(E_{0}\right), q \geq 0 \tag{9}
\end{equation*}
$$

and
(10) $\quad \kappa h+h \kappa=$ multiplication by $\sum_{1}^{r} f_{j} g_{j}$ on $\mathcal{A}_{0}^{q}\left(E_{p}\right), p>0, q \geq 0$.

We shall use the following $L^{2}$-existence theorem for the $\bar{\partial}$-operator due to Hörmander [ H ].

Theorem 3.3. - Let $q \geq 1, m \geq 1$, and let $\omega$ be a $C^{\infty}$ form of type $(0, q)$ on $\mathbb{C}^{n}$ with $\bar{\partial} \omega=0$ and such that

$$
\int_{\mathbb{C}^{n}}|\omega|^{2}\left(1+|z|^{2}\right)^{-m} d \lambda<\infty \quad(d \lambda=\text { Lebesgue measure })
$$

Then, there exists a $C^{\infty}$ form $\alpha$ of type $(0, q-1)$ on $\mathbb{C}^{n}$ with

$$
\bar{\partial} \alpha=\omega, \int_{\mathbb{C}^{n}}|\alpha|^{2}\left(1+|z|^{2}\right)^{-m-2} d \lambda \leq \int_{\mathbb{C}^{n}}|\omega|^{2}\left(1+|z|^{2}\right)^{-m} d \lambda
$$

Note. - The existence of a form $\alpha$ with $L_{\text {loc }}^{2}$ coefficients with these properties is proved in Hörmander's book [H], §4.4 (the second edition). If $H$ is the Hilbert space of forms $\beta$ of type ( $0, q-1$ ) with $\int|\beta|^{2}\left(1+|z|^{2}\right)^{-m-2} d \lambda<\infty$, we can choose $\alpha$ such that it is orthogonal, in $H$, to all forms $\beta \in H$ with $\bar{\partial} \beta=0$ [since if $\alpha^{\prime}$ is the projection of $\alpha$ onto the orthogonal complement of $\{\beta \in H, \bar{\partial} \beta=0\}$, then $\bar{\partial} \alpha=\bar{\partial} \alpha^{\prime}$ and $\left.\int\left|\alpha^{\prime}\right|^{2}\left(1+|z|^{2}\right)^{-m-2} d \lambda \leq \int|\alpha|^{2}\left(1+|z|^{2}\right)^{-m-2} d \lambda\right]$.

Now, if $q=1$, any solution of $\bar{\partial} \alpha=\omega$ is smooth. If $q>1, \alpha$ is, in particular, orthogonal to any form $\bar{\partial} \varphi$, where $\varphi$ is of type $(0, q-2)$ and is $C^{\infty}$ with compact support: $\int(\alpha, \bar{\partial} \varphi)\left(1+|z|^{2}\right)^{-m-2} d \lambda=0$. This implies that $\bar{\partial}^{*}\left(\left(1+|z|^{2}\right)^{-m-2} \alpha\right)=0$, where $\bar{\partial}^{*}$ is the formal adjoint of $\bar{\partial}$ in the Euclidean metric, viz., if $f=\sum_{|J|=q-1} f_{J} d \bar{z}^{J}$, then

$$
\bar{\partial}^{*}(f)=\sum_{|K|=q-2}\left(\sum_{j=1}^{n} \frac{\partial f_{j K}}{\partial z_{j}}\right) d \bar{z}^{K}
$$

The equations

$$
\bar{\partial} \alpha=\omega, \omega \operatorname{smooth}, \bar{\partial}^{*}\left(\left(1+|z|^{2}\right)^{-m-2} \alpha\right)=0
$$

are sufficient to guarantee the smoothness of $\alpha$ (see the proof of Theorem 4.2.5 in Hörmander's book $[\mathrm{H}])$. In fact, they can be written $\bar{\partial} \alpha=\omega$, $\bar{\partial}^{*} \alpha=L_{0}(\alpha), L_{0}$ an operator of order 0 with smooth coefficients. In particular

$$
\left(\bar{\partial} \bar{\partial}^{*}+\bar{\partial}^{*} \bar{\partial}\right) \alpha=\left(\bar{\partial} L_{0}\right)(\alpha)+\bar{\partial}^{*} \omega ;
$$

since $\bar{\partial} \bar{\partial}^{*}+\bar{\partial}^{*} \bar{\partial}$ is elliptic of order 2 and $\bar{\partial} L_{0}$ is of order 1 with smooth coefficients, once can simply apply the standard regularity theorem for elliptic operators.

Note that we are only using Hörmander's theorem for $\mathbb{C}^{n}$, not for arbitrary pseudoconvex domains. His method is very simple in this case.

Corollary 3.4. - Let $q \geq 1$ and $\omega \in \mathcal{A}^{q}\left(E_{p}\right), p \geq 0$. Assume that $\bar{\partial} \omega=0$ and that $\int_{\mathbb{C}^{n}}|\omega|^{2}\left(1+|z|^{2}\right)^{-m} d \lambda<\infty$. Then, there exists $\alpha \in \mathcal{A}^{q-1}\left(E_{p}\right)$ with

$$
\bar{\partial} \alpha=\omega, \int_{\mathbb{C}^{n}}|\alpha|^{2}\left(1+|z|^{2}\right)^{-m-2} d \lambda \leq \int_{\mathbb{C}^{n}}|\omega|^{2}\left(1+|z|^{2}\right)^{-m} d \lambda
$$

If $\omega \in \mathcal{A}^{q}\left(E_{p}\right)$ and $\int_{\mathbb{C}^{n}}|\omega|^{2}\left(1+|z|^{2}\right)^{-m} d \lambda<\infty$, we write

$$
\|\omega\|_{m}^{2}=\int_{\mathbb{C}^{n}}|\omega|^{2}\left(1+|z|^{2}\right)^{-m} d \lambda
$$

Proof of Theorem 3.1. - In the course of this proof, all constants $\gamma, \gamma^{\prime}, \gamma_{j}, m, m^{\prime}, m_{j}, C^{(j)}, N_{j}, \ldots$ will, unless otherwise stated, depend only on the constants $n, D, C_{1}, r, c_{0}, N$ defining $\mathcal{W}$; dependence on other parameters will be explicitly indicated.

Let $P \in \mathcal{W}$ and let $X=X_{1}, X_{2}, \ldots, X_{t}$ be the connected components of $\mathcal{Z}(P)=\left\{z \in \mathbb{C}^{n} \mid P(z)=0\right\}$. [The number $t$ is $\leq D^{r}$ as can be proved using Bezout's Theorem; see the remark following Lemma 5.7. We shall not need this fact.]

By Lemma 2.4, we can find $\gamma>0, m \geq 1$ so that

$$
T_{\gamma, m}\left(X_{i}\right) \cap T_{\gamma, m}\left(X_{j}\right)=\emptyset \text { if } i \neq j
$$

Now, if $S \subset \mathbb{C}^{n}$ is connected, so also is $T_{\gamma, m}(S)$ (for any $\gamma>0, m>0$ ). Thus

$$
T_{\gamma, m}(\mathcal{Z}(P))=\bigcup_{i} T_{\gamma, m}\left(X_{i}\right)
$$

is the decomposition of $T_{\gamma, m}(\mathcal{Z}(P))$ into connected components.
Choose $\gamma^{\prime}>0, \gamma^{\prime}<\gamma$ such that if $E=T_{\gamma^{\prime}, m}(\mathcal{Z}(P))$, then $T_{\gamma^{\prime}, m}(E) \subset T_{\gamma, m}(\mathcal{Z}(P))$. Set $\Omega=T_{\gamma^{\prime}, m}(X)=T_{\gamma^{\prime}, m}\left(X_{1}\right)$. Then $T_{\gamma^{\prime}, m}(\Omega) \subset$ $T_{\gamma^{\prime}, m}(E) \subset \bigcup_{i} T_{\gamma, m}\left(X_{i}\right)$ and $T_{\gamma^{\prime}, m}(\Omega)$ is connected; since the $T_{\gamma, m}\left(X_{i}\right)$ are pointwise disjoint and connected, and $T_{\gamma^{\prime}, m}(\Omega) \cap T_{\gamma, m}\left(X_{1}\right) \neq \emptyset$, it follows that

$$
T_{\gamma^{\prime}, m}(\Omega) \subset T_{\gamma, m}(X)
$$

Hence, using Lemma 2.8, we obtain the following:
(11) There exists $\varphi \in C^{\infty}\left(\mathbb{C}^{n}\right), 0 \leq \varphi \leq 1$, such that

$$
\varphi=1 \text { on } \Omega, \operatorname{supp}(\varphi) \subset T_{\gamma, m}(X)
$$

and such that, for $k \geq 0$, there are constants $C_{k}$ depending only on $k, \gamma, \gamma^{\prime}, m, n$ with

$$
\left|D^{(k)} \varphi(z)\right| \leq C_{k}\left(1+|z|^{2}\right)^{k m}, z \in \mathbb{C}^{n}, k \geq 0
$$

for any differentiation $D^{(k)}$ of order $k$.
Clearly $\varphi \mid X_{j}=0$ if $j \geq 2$.
Lemma 2.5 gives us:
(12) There exist constants $\gamma_{1}>0, N_{1} \geq 1$ such that

$$
\begin{aligned}
|P(z)|^{2} & =\sum_{j=1}^{r}\left|p_{j}(z)\right|^{2} \\
& \geq \gamma_{1}\left(1+|z|^{2}\right)^{-N_{1}} \text { if } z \notin T_{\gamma^{\prime}, m}(\mathcal{Z}(P))=\bigcup_{i=1}^{t} T_{\gamma^{\prime}, m}\left(X_{i}\right) .
\end{aligned}
$$

Consider the Koszul complex defined by the functions $\left(p_{1}, \ldots, p_{r}\right)$ :

$$
\begin{equation*}
0 \rightarrow \mathcal{A}^{q}\left(E_{r}\right) \xrightarrow{\kappa} \cdots \rightarrow \mathcal{A}^{q}\left(E_{1}\right) \xrightarrow{\kappa} \mathcal{A}^{q}\left(E_{0}\right) \rightarrow 0, \tag{13}
\end{equation*}
$$

$\kappa$ being defined by:

$$
\begin{aligned}
& \kappa\left(e_{I} \otimes d \bar{z}^{J}\right)=\sum_{k=1}^{p}(-1)^{k-1} p_{i_{k}}(z) e_{I \backslash\left\{i_{k}\right\}} \otimes d \bar{z}^{J} \\
& \quad I=\left(i_{1}<\cdots<i_{p}\right), J=\left(j_{1}<\cdots<j_{q}\right)
\end{aligned}
$$

Let $\mathcal{A}_{0}^{q}\left(E_{p}\right)$ be the space of smooth $(0, q)$ forms $\omega$ on $\mathbb{C}^{n}$ with values in $E_{p}$ such that

$$
\operatorname{supp}(\omega) \cap \mathcal{Z}(P)=\emptyset
$$

We define a map $h: \mathcal{A}_{0}^{q}\left(E_{p}\right) \rightarrow \mathcal{A}_{0}^{q}\left(E_{p+1}\right)$ as in (8) above with $\Omega_{0}=\mathbb{C}^{n} \backslash \mathcal{Z}(P), g_{j}=\bar{p}_{j} /|P|^{2}:$

$$
h\left(\sum_{I, J} v_{J}^{I} e_{I} \otimes d \bar{z}^{J}\right)=\sum_{I, J} v_{J}^{I}\left(\sum_{j=1}^{r} \frac{\bar{p}_{j}}{|P|^{2}} e_{j}\right) \wedge e_{I} \otimes d \bar{z}^{J}
$$

Since $p_{1} \frac{\bar{p}_{1}}{|P|^{2}}+\cdots+p_{r} \frac{\bar{p}_{r}}{|P|^{2}}=1$ on $\Omega_{0}=\mathbb{C}^{n} \backslash \mathcal{Z}(P)$, we obtain, from (9) and (10) above, the following:

$$
\begin{align*}
\kappa h & =\text { identity on } \mathcal{A}_{0}^{q}\left(E_{0}\right), q \geq 0, \text { and }  \tag{14}\\
\kappa h+h \kappa & =\text { identity on } \mathcal{A}_{0}^{q}\left(E_{p}\right) \text { for } p>0, q \geq 0 .
\end{align*}
$$

Lemma 3.5. - There is a sequence $\left\{\bar{C}_{k}\right\}_{k \geq 0}$ of constants depending only on $k$ and the constants in $\mathcal{W}$, and $N_{2}>0$ (depending only on the constants in $\mathcal{W}$ ) such that

If $\omega \in \mathcal{A}_{0}^{q}\left(E_{p}\right)$ is such that $\operatorname{supp}(\omega) \cap T_{\gamma^{\prime}, m}(\mathcal{Z}(P))=\emptyset$, and if we have estimates

$$
\left|D^{(k)} \omega(z)\right| \leq A_{k}\left(1+|z|^{2}\right)^{M_{k}}, z \in \mathbb{C}^{n}, k \geq 0, A_{k} \leq A_{k+1}, M_{k} \leq M_{k+1}
$$

then

$$
\left|D^{(k)} h(\omega)(z)\right| \leq \bar{C}_{k} A_{k}\left(1+|z|^{2}\right)^{M_{k}+(k+1) N_{2}}, k \geq 0, z \in \mathbb{C}^{n}
$$

here $D^{(k)}$ runs over all differentiations of order $k$.
Proof. - The coefficients of the form $h(\omega)$ are linear combinations of functions of the form

$$
\frac{\bar{p}_{j}}{|P|^{2}} u
$$

$u$ being the coefficient of $e_{I} \otimes d \bar{z}^{J}$ in $\omega$ for some $I, J$. Since $\left\|p_{j}\right\| \leq C_{1}$, $\operatorname{deg} p_{j} \leq D$ and $|P(z)|^{2} \geq \gamma_{1}\left(1+|z|^{2}\right)^{-N_{1}}$ on the support of $u$ (by (12)), we have, for $\ell \geq 0$,

$$
\left|D^{(\ell)}\left(\frac{\bar{p}_{j}}{|P|^{2}}\right)\right| \leq C_{\ell}^{\prime}\left(1+|z|^{2}\right)^{3 \ell D+(\ell+1) N_{1}}
$$

( $C_{\ell}^{\prime}$ depending only on $\ell$ and the constants in $\mathcal{W}$ ), while $\left|D^{(k-\ell)} u\right| \leq$ $A_{k}\left(1+|z|^{2}\right)^{M_{k}}(0 \leq \ell \leq k)$. The lemma follows.

Returning to the proof of Theorem 3.1, consider the function $\varphi$ described in (11) above. Since $\varphi=1$ on $\Omega=T_{\gamma^{\prime}, m}(X)$ and $\operatorname{supp}(\varphi) \subset$ $T_{\gamma, m}(X)$, we have

$$
\operatorname{supp}(\bar{\partial} \varphi) \subset T_{\gamma, m}(X) \backslash \Omega
$$

We define, successively,

$$
\varphi_{1}=h(\bar{\partial} \varphi), \varphi_{2}=h\left(\bar{\partial} \varphi_{1}\right), \ldots, \varphi_{p}=h\left(\bar{\partial} \varphi_{p-1}\right), \ldots, \varphi_{r}=h\left(\bar{\partial} \varphi_{r-1}\right)
$$

These are well-defined since $\operatorname{supp}(\bar{\partial} \varphi) \cap \mathcal{Z}(P)=\emptyset$ and neither $\bar{\partial}$ nor $h$ increases supports. Moreover, we have

$$
\varphi_{p} \in \mathcal{A}_{0}^{p}\left(E_{p}\right) \text { for } 1 \leq p \leq r
$$

Now $\bar{\partial} \varphi \in \mathcal{A}_{0}^{1}\left(E_{0}\right)$ and, by (14) $\kappa\left(\varphi_{1}\right)=\kappa h(\bar{\partial} \varphi)=\bar{\partial} \varphi$. Hence $\kappa\left(\bar{\partial} \varphi_{1}\right)=\bar{\partial} \kappa\left(\varphi_{1}\right)=\bar{\partial}(\bar{\partial} \varphi)=0$, and (14) now gives

$$
\bar{\partial} \varphi_{1}=(\kappa h+h \kappa)\left(\bar{\partial} \varphi_{1}\right)=\kappa h\left(\bar{\partial} \varphi_{1}\right)=\kappa\left(\varphi_{2}\right)
$$

By induction, we obtain
(15) For $1 \leq p \leq r, \operatorname{supp}\left(\varphi_{p}\right) \subset T_{\gamma, m}(X) \backslash \Omega$, and

$$
\begin{aligned}
& \kappa\left(\varphi_{p}\right)=\bar{\partial} \varphi_{p-1} \text { for } 2 \leq p \leq r \\
& \kappa\left(\varphi_{1}\right)=\bar{\partial} \varphi
\end{aligned}
$$

Moreover, by Lemma 3.5 and induction on $p$, we obtain [since $\operatorname{supp}(\bar{\partial} \varphi) \cap T_{\gamma^{\prime}, m}(\mathcal{Z}(P))=\emptyset$ and $\left|D^{(k)} \bar{\partial} \varphi(z)\right| \leq C_{k+1}\left(1+|z|^{2}\right)^{(k+1) m}$ by (11)]:
(16) There is a constant $N_{3}>0$ (depending only on the constants in $\mathcal{W}$ ) and a sequence $\left\{C_{k}^{\prime \prime}\right\}_{k \geq 0}$ of constants depending only on $k$ and the constants in $\mathcal{W}$ such that

$$
\left|D^{(k)} \varphi_{p}(z)\right| \leq C_{k}^{\prime \prime}\left(1+|z|^{2}\right)^{(k+1) N_{3}}, z \in \mathbb{C}^{n}, k \geq 0,1 \leq p \leq r
$$

We now make the following remark:
(17) For $q \geq 0$, the $\operatorname{map} \kappa: \mathcal{A}_{0}^{q}\left(E_{r}\right) \rightarrow \mathcal{A}_{0}^{q}\left(E_{r-1}\right)$ is injective.

In fact, if $v \in E_{r}=\stackrel{r}{\wedge} \mathbb{C}^{r}, v \neq 0$, we can write $v=u e_{1} \wedge \cdots \wedge e_{r} u \in \mathbb{C}$, $u \neq 0$ and

$$
\kappa\left(u e_{1} \wedge \cdots \wedge e_{r} \otimes d \bar{z}^{J}\right)=\sum_{k=1}^{r}(-1)^{k-1} p_{k}(z) e_{1} \wedge \cdots \wedge \hat{e}_{k} \wedge \cdots \wedge e_{r} \otimes d \bar{z}^{J}
$$

which is $\neq 0$ unless all the $p_{k}(z)=0$, hence is $\neq 0$ if $z \notin \mathcal{Z}(P)$.
Now, $|P(z)|^{2} \leq C^{(0)}\left(1+|z|^{2}\right)^{D}$ since $p_{j} \in \mathcal{H}^{D}$ and $\|P\| \leq C_{1}$. Hence:

$$
\begin{equation*}
|\kappa(\omega)(z)|^{2} \leq C^{(1)}\left(1+|z|^{2}\right)^{D}|\omega(z)|^{2}, z \in \mathbb{C}^{n}, \omega \in \mathcal{A}^{q}\left(E_{p}\right) \tag{18}
\end{equation*}
$$

By construction, $\bar{\partial} \varphi_{r} \in \mathcal{A}_{0}^{r+1}\left(E_{r}\right)$ and $\kappa\left(\bar{\partial} \varphi_{r}\right)=\bar{\partial}\left(\kappa\left(\varphi_{r}\right)\right)=\bar{\partial}\left(\bar{\partial} \varphi_{r-1}\right)$ (by (15)) $=0$ [ $\kappa$ and $\bar{\partial}$ commute since the $p_{j}$ are holomorphic]. Thus, by (17), we must have $\bar{\partial} \varphi_{r}=0$. Moreover, if we take $p=r, k=0$ in (16), we find that

$$
\int_{\mathbb{C}^{n}}\left|\varphi_{r}\right|^{2}\left(1+|z|^{2}\right)^{-N_{3}-n-1} d \lambda \leq C^{(2)}
$$

By Corollary 3.4, there exists $\psi_{r-1} \in \mathcal{A}^{r-1}\left(E_{r}\right)$ with

$$
\begin{equation*}
\bar{\partial} \psi_{r-1}=\varphi_{r},\left\|\psi_{r-1}\right\|_{m_{1}}^{2}=\int_{\mathbb{C}^{n}}\left|\psi_{r-1}\right|^{2}\left(1+|z|^{2}\right)^{-m_{1}} d \lambda \leq C^{(2)} \tag{19}
\end{equation*}
$$

where $m_{1}=N_{3}+n+3$.
Consider now $\varphi_{r-1}-\kappa\left(\psi_{r-1}\right) \in \mathcal{A}^{r-1}\left(E_{r-1}\right)$. By (15), (19), we have $\bar{\partial}\left(\varphi_{r-1}-\kappa\left(\psi_{r-1}\right)\right)=\bar{\partial} \varphi_{r-1}-\kappa\left(\bar{\partial} \psi_{r-1}\right)=\bar{\partial} \varphi_{r-1}-\kappa\left(\varphi_{r}\right)=0$. Further, by (16), (18), (19), there exist $m_{2}^{\prime}>0$ and $C^{(3)}>0$ so that

$$
\begin{equation*}
\left\|\varphi_{r-1}-\kappa\left(\psi_{r-1}\right)\right\|_{m_{2}^{\prime}}^{2} \leq C^{(3)} \tag{20}
\end{equation*}
$$

and we can apply Corollary 3.4 again to find $\psi_{r-2} \in \mathcal{A}^{r-2}\left(E_{r-1}\right)$ such that, with $m_{2}=m_{2}^{\prime}+2$,

$$
\bar{\partial} \psi_{r-2}=\varphi_{r-1}-\kappa\left(\psi_{r-1}\right),\left\|\psi_{r-2}\right\|_{m_{2}}^{2} \leq C^{(3)}
$$

We can iterate this procedure to find $\psi_{r}=0, \psi_{r-1}, \ldots, \psi_{0}$, and constants $m_{2}>0, m_{3}>0, \ldots, m_{r}>0, C^{(2)}, C^{(3)}, \ldots, C^{(r+1)}>0$ such that

$$
\begin{align*}
\psi_{r-q} & \in \mathcal{A}^{r-q}\left(E_{r-q+1}\right) \\
\bar{\partial} \psi_{r-q} & =\varphi_{r-q+1}-\kappa\left(\psi_{r-q+1}\right)  \tag{21}\\
\left\|\psi_{r-q}\right\|_{m_{q}}^{2} & \leq C^{(q+1)} \quad \text { for } \quad 1 \leq q \leq r
\end{align*}
$$

[Note that if the results in (21) have been proved for $1 \leq q \leq q_{0}-1$, then

$$
\begin{aligned}
\bar{\partial}\left(\varphi_{r-q_{0}+1}-\kappa\left(\psi_{r-q_{0}+1}\right)\right) & =\bar{\partial} \varphi_{r-q_{0}+1}-\kappa\left(\bar{\partial} \psi_{r-q_{0}+1}\right) \\
& =\bar{\partial} \varphi_{r-q_{0}+1}-\kappa\left(\varphi_{r-q_{0}+2}-\kappa\left(\psi_{r-q_{0}+2}\right)\right) \\
& =\bar{\partial} \varphi_{r-q_{0}+1}-\kappa\left(\psi_{r-q_{0}+2}\right)=0
\end{aligned}
$$

by (15), and we can apply Hörmander's theorem, Cor. 3.4, to obtain (21) for $q=q_{0}$.]

We have $\psi_{0} \in \mathcal{A}^{0}\left(E_{1}\right)$. Consider $\varphi-\kappa\left(\psi_{0}\right) \in \mathcal{A}^{0}\left(E_{0}\right)=C^{\infty}\left(\mathbb{C}^{n}\right)$. Now, $\bar{\partial}\left(\varphi-\kappa\left(\psi_{0}\right)\right)=\bar{\partial} \varphi-\kappa\left(\bar{\partial} \psi_{0}\right)=\bar{\partial} \varphi-\kappa\left(\varphi_{1}-\kappa\left(\psi_{1}\right)\right)=\bar{\partial} \varphi-\kappa\left(\varphi_{1}\right)=0$ by (15). Thus:
(22) The function $F=\varphi-\kappa\left(\psi_{0}\right)$ is holomorphic on $\mathbb{C}^{n}$.

Moreover, by (18) and (21), we have

$$
\int_{\mathbb{C}^{n}}\left|\kappa\left(\psi_{0}\right)\right|^{2}\left(1+|z|^{2}\right)^{-m_{r}-D} d \lambda \leq C^{(1)} C^{(r+1)}
$$

while $\int_{\mathbb{C}^{n}}|\varphi|^{2}\left(1+|z|^{2}\right)^{-n-1} d \lambda \leq \int_{\mathbb{C}^{n}}\left(1+|z|^{2}\right)^{-n-1} d \lambda(<\infty)$. Consequently, there exist constants $\bar{m}(>n+1)$ and $\bar{C}>0$ depending only on the constants in $\mathcal{W}$ such that

$$
\begin{equation*}
\int_{\mathbb{C}^{n}}|F(z)|^{2}\left(1+|z|^{2}\right)^{-\bar{m}} d \lambda \leq \bar{C} \tag{23}
\end{equation*}
$$

Now, $\psi_{0} \in \mathcal{A}^{0}\left(E_{1}\right)$, and we can write $\psi_{0}=\sum_{j=1}^{r} \psi_{0 j} e_{j}$, where $\psi_{0 j} \in$ $C^{\infty}\left(\mathbb{C}^{n}\right)$. Hence $\kappa\left(\psi_{0}\right)=\sum_{j=1}^{r} \psi_{0 j} p_{j}$. In particular, $\kappa\left(\psi_{0}\right) \mid \mathcal{Z}(P) \equiv 0$. Thus we have

$$
\begin{equation*}
F=\varphi \text { on } \mathcal{Z}(P) ; \text { i.e. } F|X \equiv 1, F| X_{j} \equiv 0 \text { for } j>1 \tag{24}
\end{equation*}
$$

Thus, we have only to show that $F$ is a polynomial satisfying the required estimates to complete the proof of Theorem 3.1.

Since $F$ is holomorphic, the function $|F|^{2}$ is subharmonic. Hence, if $w \in \mathbb{C}^{n}$ and $R>0$, we have

$$
|F(w)|^{2} \leq \frac{1}{\operatorname{vol}(B(w, R))} \int_{B(w, R)}|F(z)|^{2} d \lambda
$$

If $|w| \geq 1$, we take $R=\frac{1}{2}|w|$ and remark that $1+|z|^{2} \leq 3\left(1+|w|^{2}\right)$ for $|z-w|<\frac{1}{2}|w|$, and $\operatorname{vol}(B(w, R))=\operatorname{vol}\left(B\left(0, \frac{1}{2}\right)\right) \cdot|w|^{2 n}$. Hence

$$
\begin{aligned}
|F(w)|^{2} & \leq c_{1}(\bar{m}, n) \cdot\left(1+|w|^{2}\right)^{\bar{m}}|w|^{-2 n} \int_{\mathbb{C}^{n}} \frac{|F(z)|^{2}}{\left(1+|z|^{2}\right)^{\bar{m}}} d \lambda \\
& \leq c_{2}(\bar{m}, n) \cdot \bar{C}\left(1+|w|^{2}\right)^{\bar{m}-n} \quad(\text { by }(23))
\end{aligned}
$$

where $c_{1}(\bar{m}, n)$ and $c_{2}(\bar{m}, n)$ depend only on $\bar{m}$ and $n$. It follows that $F$ is a polynomial of degree $\leq \bar{m}-n$.

Also, for $|w| \leq n$, we have

$$
\begin{aligned}
|F(w)|^{2} & \leq c_{n} \int_{|z| \leq 2 n}|F(z)|^{2} d \lambda \leq c_{3}(\bar{m}, n) \int_{\mathbb{C}^{n}} \frac{|F(z)|^{2}}{\left(1+|z|^{2}\right)^{\bar{m}}} d \lambda \\
& \leq C_{\#}, C_{\#} \text { depending only on the constants in } \mathcal{W}
\end{aligned}
$$

By Cauchy's inequalities:

$$
\left|\frac{\partial^{|\alpha|} F}{\partial z_{1}^{\alpha_{1}} \cdots \partial z_{n}^{\alpha_{n}}}(0)\right| \leq \alpha!C_{\#}, \alpha=\left(\alpha_{1}, \ldots, \alpha_{n}\right)
$$

Since $F$ is a polynomial of degree $\leq D^{\prime}=\bar{m}-n$, this proves that $\|F\| \leq C^{\prime}$ where $C^{\prime}$ depends only on the constants in $\mathcal{W}$. This and (24) complete the proof of Theorem 3.1.

A similar argument, combined with Lemma 2.6, can be used to prove the following extension theorem.

Theorem 3.6. - There exists a constant $D^{\prime}>0$ depending only on the constants in $\mathcal{W}$, and, for $M \geq 0$, a constant $C(M)$ depending only on $M$ and the constants in $\mathcal{W}$ such that the following holds.

Let $P \in \mathcal{W}$ and let $f$ be a holomorphic function on $\mathcal{Z}(P)$. Assume that there are constants $A>0, M \geq 0$ such that

$$
|f(z)|^{2} \leq A\left(1+|z|^{2}\right)^{M} \text { for } z \in \mathcal{Z}(P)
$$

Then, there exists a polynomial $F$ of degree $\leq M+D^{\prime}$ such that $F \mid \mathcal{Z}(P)=f$ and $\|F\|^{2} \leq C(M) A$.

Sketch of Proof. - By Lemma 2.6, we can find $\gamma_{0}, m_{0}$ and a holomorphic retraction $\pi: T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P)) \rightarrow \mathcal{Z}(P),|\pi(z)-z|<1$. Choose $\gamma>0\left(\gamma<\gamma_{0}\right)$ such that if $\Omega=T_{\gamma, m_{0}}(\mathcal{Z}(P))$ then $T_{\gamma, m_{0}}(\Omega) \subset$ $T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$, and $\varphi \in C^{\infty}\left(\mathbb{C}^{n}\right)$ with $\varphi=1$ on $\Omega, \operatorname{supp}(\varphi) \subset T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$, $\left|D^{(k)} \varphi(z)\right| \leq C_{k}\left(1+|z|^{2}\right)^{k m_{0}}(k \geq 0)$. Consider the $C^{\infty}$ function on
$\mathbb{C}^{n}: \varphi_{0}=\varphi \cdot(f \circ \pi)$ on $T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P)), \varphi_{0}=0$ outside $T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$. The form $\omega \in \mathcal{A}^{1}\left(E_{0}\right)$ defined by $\omega=\bar{\partial} \varphi_{0}(=(f \circ \pi) \cdot \bar{\partial} \varphi)$ has support outside $\Omega$. We can repeat the proof of Theorem 3.1 to find a controlled function $\psi_{0} \in \mathcal{A}^{0}\left(E_{1}\right)$ such that $\bar{\partial}\left(\varphi_{0}-\kappa\left(\psi_{0}\right)\right)=0$, so that $F=\varphi_{0}-\kappa\left(\psi_{0}\right)$ is holomorphic and $F\left|\mathcal{Z}(P)=\varphi_{0}\right| \mathcal{Z}(P)=f$. The argument in Theorem 3.1 gives the estimate

$$
\int_{\mathbb{C}^{n}}|F|^{2}\left(1+|z|^{2}\right)^{-M-D^{\prime}} d \lambda \leq C^{\prime} A
$$

( $C^{\prime}, D^{\prime}$ depending only on the constants in $\mathcal{W}$ ), and the proof is completed as above, using the subharmonicity of $|F|^{2}$.

## 4. Polynomials vanishing on a smooth variety.

In this section, we prove the following theorem.
The referee has informed us that a stronger version of Theorem 4.1 below is contained in recent work of F. Amoroso [A].

Theorem 4.1. - Let $n \geq 2, D \geq 1,1 \leq r \leq n-1, C_{1}>0$, $c_{0}>0, N \geq 1$ be given, and let $\mathcal{W}$ be, as before, the space of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right)$ with $p_{j} \in \mathcal{H}^{D}, P(0)=0,\|P\| \leq C_{1}, J_{P}(0)=1$ and $\left|J_{P}(z)\right|^{2}+|P(z)|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N}$, where $J_{P}(z)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}\right)_{1 \leq j, k \leq r}$.

There exists a constant $D^{\prime} \geq 1$ depending only on the constants defining $\mathcal{W}$, and, for $d \geq 1$ a constant $C(d)$ depending only on $d$ and the constants in $\mathcal{W}$ such that the following holds.

Let $d \geq 1$ be an integer, and let $F \in \mathcal{H}^{d}$ be such that $F(z)=0$, $\forall z \in \mathcal{Z}(P)\left(=\left\{z \in \mathbb{C}^{n} \mid P(z)=0\right\}\right)$.

Then, there exist polynomials $F_{1}, \ldots, F_{r}$ such that

$$
\begin{array}{ll}
\text { (1.a) } & F=\sum_{j=1}^{r} F_{j} p_{j} \\
\text { (1.b) } & \operatorname{deg} F_{j} \leq d+D^{\prime}  \tag{1.b}\\
\text { (1.c) } & \left\|F_{j}\right\| \leq C(d)\|F\| .
\end{array}
$$

Proof. - During the course of this proof, constants $\gamma, \gamma^{\prime}, \gamma_{j}, m$, $m^{\prime}, m_{j}, C, C^{\prime}, C^{(j)}$ etc. will, unless otherwise stated, depend only on the constants defining $\mathcal{W}$. If $k, d, \ldots$ are other parameters, we write $C(k, \mathcal{W})$,
$C(d, \mathcal{W})$ etc. to indicate constants depending only on the parameters in question and on the constants defining $\mathcal{W}$.

If $z=\left(z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n}$, we write $z=(x, \zeta)$ with $x=\left(x_{1}, \ldots, x_{r}\right)=$ $\left(z_{1}, \ldots, z_{r}\right)$ and $\zeta=\left(\zeta_{1}, \ldots, \zeta_{n-r)}\right)=\left(z_{r+1}, \ldots, z_{n}\right)$.

Given $G \in \mathcal{H}^{d}$, we can write

$$
\begin{equation*}
G(x, \zeta)-G(y, \zeta)=\sum_{j=1}^{r}\left(x_{j}-y_{j}\right) A_{G}^{(j)}(x, y, \zeta) \tag{2}
\end{equation*}
$$

where

$$
\begin{equation*}
A_{G}^{(j)}(x, y, \zeta)=\int_{0}^{1} \frac{\partial G}{\partial x_{j}}(t x+(1-t) y, \zeta) d t \tag{3}
\end{equation*}
$$

Clearly
(4) $A_{G}^{(j)}$ is a polynomial of degree $\leq d-1$ in $x, y, \zeta$ and $\left\|A_{G}^{(j)}\right\| \leq$ $c_{1}(d, n)\|G\|$, where $\left\|A_{G}^{(j)}\right\|$ denotes the norm in the space of polynomials in $n+r$ variables and $c_{1}(d, n)$ depends only on $d$ and $n$.

By Lemma 2.6, we have the following:
(5) There exist $\gamma_{0}>0, m_{0} \geq 1$ such that, if $P \in \mathcal{W}$ and $\Omega=T_{\gamma_{0}, m_{0}}(\mathcal{Z}(P))$, we can find a holomorphic map

$$
\pi(x, \zeta)=\left(\pi_{0}(x, \zeta), \zeta\right), \quad((x, \zeta) \in \Omega)
$$

of $\Omega$ into $\mathbb{C}^{n}$ with the properties: $\pi(\Omega) \subset \mathcal{Z}(P), \pi \mid \mathcal{Z}(P)=$ identity, and $\left|\pi_{0}(x, \zeta)-x\right|<1$. (By definition, $\pi_{0}$ is a holomorphic map of $\Omega$ into $\mathbb{C}^{r}$.)

We write $\pi_{0}(x, \zeta)=\left(\pi_{1}(x, \zeta), \ldots, \pi_{r}(x, \zeta)\right)$. By Cauchy's inequalities, if $k \geq 0$ and $(x, \zeta)=z \in T_{\gamma_{0} / 2, m_{0}}(\mathcal{Z}(P))$, then

$$
\begin{equation*}
\left|D^{(k)} \pi_{j}(x, \zeta)\right| \leq C_{1}(k, \mathcal{W})\left(1+|z|^{2}\right)^{k m_{0}+1} \tag{6}
\end{equation*}
$$

We use formulae (2), (3) for $p_{i}$ to write

$$
p_{i}(x, \zeta)-p_{i}(y, \zeta)=\sum_{j=1}^{r} p_{i j}(x, y, \zeta)\left(x_{j}-y_{j}\right)
$$

with $p_{i j}(x, y, \zeta)=\int_{0}^{1} \frac{\partial p_{i}}{\partial x_{j}}(t x+(1-t) y, \zeta) d t$. If $z=(x, \zeta) \in \Omega$, we have

$$
p_{i}(x, \zeta)=\sum_{j=1}^{r} P_{i j}(x, \zeta)\left(x_{j}-\pi_{j}(x, \zeta)\right)
$$

where $P_{i j}(x, \zeta)=P_{i j}\left(x, \pi_{0}(x, \zeta), \zeta\right)$. Using (4) with $G=p_{i}$ and the fact that $\left|\pi_{0}(x, \zeta)-x\right|<1$, we obtain:
(7) if $z=(x, \zeta) \in \Omega,\left|P_{i j}(x, \zeta)\right| \leq C^{(1)}(1+|z|)^{D-1} \leq C^{(2)}\left(1+|z|^{2}\right)^{D / 2}$.

Now, if $z=(x, \zeta) \in \mathcal{Z}(P)$, we have $\pi(z)=z, \pi_{0}(x, \zeta)=x$, so that
(8) if $(x, \zeta) \in \mathcal{Z}(P)$, then $P_{i j}(x, \zeta)=\int_{0}^{1} \frac{\partial p_{i}}{\partial z_{j}}(t x+(1-t) x, \zeta) d t=\frac{\partial p_{i}}{\partial z_{j}}(x, \zeta)$.

Since $\left|\operatorname{det}\left(\frac{\partial p_{i}}{\partial z_{j}}\right)\right|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N}$ if $z \in \mathcal{Z}(P)$, we have

$$
\begin{equation*}
\left|\operatorname{det}\left(P_{i j}(x, \zeta)\right)\right|^{2} \geq c_{0}\left(1+|z|^{2}\right)^{-N} \text { if } z \in \mathcal{Z}(P) \tag{9}
\end{equation*}
$$

From (7) and Cauchy's inequalities we obtain:
(10) If $z=(x, \zeta) \in T_{\gamma_{0} / 2, m_{0}}(\mathcal{Z}(P))$, we have

$$
\left|D^{(k)} P_{i j}(x, \zeta)\right| \leq C_{2}(k, \mathcal{W})\left(1+|z|^{2}\right)^{D / 2+k m_{0}}, k \geq 0
$$

From (9) and (10), we deduce the following:
(11) There are constants $\gamma_{1}, m_{1}$ and $\gamma_{2}, m_{2}, \gamma_{\nu}>0, m_{\nu} \geq 1$, such that $T_{\gamma_{1}, m_{1}}\left((\mathcal{Z}(P)) \subset T_{\gamma_{0} / 2, m_{0}}(\mathcal{Z}(P))\right.$ and $\left|\operatorname{det}\left(P_{i j}(x, \zeta)\right)\right|^{2} \geq \gamma_{2}\left(1+|z|^{2}\right)^{-m_{2}}$ for $(x, \zeta)=z \in T_{\gamma_{1}, m_{1}}(\mathcal{Z}(P))$.

From (7) and (11) we obtain:
(12) There exist constants $C^{(3)}, m_{3} \geq 1$ such that the matrix $\left(P_{i j}(x, \zeta)\right)$ is invertible if $(x, \zeta) \in T_{\gamma_{1}, m_{1}}(\mathcal{Z}(P))$ and the inverse matrix $\left(Q_{i j}(x, \zeta)\right)$ satisfies

$$
\left|Q_{i j}(x, \zeta)\right| \leq C^{(3)}\left(1+|z|^{2}\right)^{m_{3}},(x, \zeta)=z \in T_{\gamma_{1}, m_{1}}(\mathcal{Z}(P))
$$

(13) For $1 \leq i \leq r$

$$
\begin{aligned}
\sum_{k=1}^{r} Q_{i k}(x, \zeta) p_{k}(x, \zeta) & =\sum_{j=1}^{r} \sum_{k=1}^{r} Q_{i k}(x, \zeta) P_{k j}(x, \zeta)\left(x_{j}-\pi_{j}(x)\right) \\
& =x_{i}-\pi_{i}(x, \zeta)
\end{aligned}
$$

Let now $F$ be a polynomial of degree $\leq d$ such that $F \mid \mathcal{Z}(P)=0$. If $z=(x, \zeta) \in T_{\gamma_{1}, m_{1}}(\mathcal{Z}(P))$, we have (using (2), (3) with $G=F$ )

$$
\begin{align*}
F(z) & =F(x, \zeta)-F\left(\pi_{0}(x, \zeta), \zeta\right)  \tag{14}\\
& =\sum_{j=1}^{r} A_{F}^{(j)}\left(x, \pi_{0}(x, \zeta), \zeta\right)\left(x_{j}-\pi_{j}(x, \zeta)\right) \\
& =\sum_{j, k=1}^{r} A_{F}^{(j)}\left(x, \pi_{0}(x, \zeta), \zeta\right) Q_{j k}(x, \zeta) p_{k}(x, \zeta)
\end{align*}
$$

Using (4) for $G=F$ and the fact that $\left|\pi_{0}(x, \zeta)-x\right|<1$, we see that $\left|A_{F}^{(j)}\left(x, \pi_{0}(x, \zeta), \zeta\right)\right|^{2} \leq c_{3}(d, \mathcal{W})\|F\|^{2}\left(1+|z|^{2}\right)^{d}$. This, (13) and (14) give us:
(15) There exist constants $\gamma>0, m, \bar{m} \geq 1$ depending only on the constants in $\mathcal{W}$, and a constant $\bar{C}(d, \mathcal{W})$ (depending only on $d \geq 1$ and the constants in $\mathcal{W}$ ) such that the following holds.

If $P \in \mathcal{W}, d \geq 1, F \in H^{d}$ and $F \mid \mathcal{Z}(P)=0$, then we can find holomorphic functions $g_{1}, \ldots, g_{r}$ on $T_{\gamma, m}(\mathcal{Z}(P))$ such that, for $z \in T_{\gamma, m}(\mathcal{Z}(P))$,

$$
F(z)=g_{1}(z) p_{1}(z)+\cdots+g_{r}(z) p_{r}(z)
$$

and

$$
\left|g_{1}(z)\right|^{2}+\cdots+\left|g_{r}(z)\right|^{2} \leq \bar{C}(d, \mathcal{W})\|F\|^{2}\left(1+|z|^{2}\right)^{d+\bar{m}}
$$

We choose $\left(\gamma^{\prime}, m^{\prime}\right),\left(\gamma^{\prime \prime}, m^{\prime \prime}\right)$ such that, if $P \in \mathcal{W}, B\left(w, \gamma^{\prime}(1+\right.$ $\left.\left.|w|^{2}\right)^{-m^{\prime}}\right) \subset T_{\gamma^{\prime \prime}, m^{\prime \prime}}(\mathcal{Z}(P))$ for any $w \in T_{\gamma^{\prime}, m^{\prime}}(\mathcal{Z}(P))$ and $B\left(w, \gamma^{\prime \prime}(1+\right.$ $\left.\left.|w|^{2}\right)^{-m^{\prime \prime}}\right) \subset T_{\gamma, m}(\mathcal{Z}(P))$ for any $w \in T_{\gamma^{\prime \prime}, m^{\prime \prime}}(\mathcal{Z}(P))$. We may assume that $\gamma^{\prime \prime}<\gamma_{1}, m^{\prime \prime}>m_{1}\left(\gamma_{1}, m_{1}\right.$ as in (11)).

By Lemma 2.8, we can find $\varphi \in C^{\infty}\left(\mathbb{C}^{n}\right)$ such that $0 \leq \varphi \leq 1, \varphi=1$ on $T_{\gamma^{\prime}, m^{\prime}}(\mathcal{Z}(P)), \operatorname{supp}(\varphi) \subset T_{\gamma^{\prime \prime}, m^{\prime \prime}}(\mathcal{Z}(P))$ and $\left|D^{(k)} \varphi(z)\right| \leq C_{k}(\mathcal{W})(1+$ $\left.|z|^{2}\right)^{k m^{\prime}}, k \geq 0, C_{k}(\mathcal{W})$ being a constant depending only on $k$ and the constants in $\mathcal{W}$.

Let $P \in \mathcal{W}$, let $F \in \mathcal{H}^{d}, F \mid \mathcal{Z}(P)=0$. Let $g_{j}(1 \leq j \leq r)$ be the holomorphic functions on $T_{\gamma, m}(\mathcal{Z}(P))$ constructed in (15).

For $1 \leq j \leq r$, we define

$$
\Phi_{j}=\varphi g_{j}+(1-\varphi) \frac{\bar{p}_{j} F}{|P|^{2}}
$$

Let $\gamma_{3}>0, m_{3} \geq 1$ be such that

$$
\begin{equation*}
|P(z)|^{2} \geq \gamma_{3}\left(1+|z|^{2}\right)^{-m_{3}} \quad \text { if } z \in \mathbb{C}^{n} \backslash T_{\gamma^{\prime}, m^{\prime}}(\mathcal{Z}(P)) \tag{16}
\end{equation*}
$$

these exist by Lemma 2.5.
Lemma 4.2. - The functions $\Phi_{j}$ defined above have the following properties:
(i) $\Phi_{j} \in C^{\infty}\left(\mathbb{C}^{n}\right), \Phi_{j} \mid T_{\gamma^{\prime}, m^{\prime}}(\mathcal{Z}(P))$ is holomorphic.
(ii) $F(z)=\sum_{1}^{r} \Phi_{j}(z) p_{j}(z)$ for $z \in \mathbb{C}^{n}$.
(iii) If $k \geq 0$ and $D^{(k)}$ is any differentiation of order $k$, we have

$$
\left|D^{(k)} \Phi_{j}(z)\right|^{2} \leq C_{k}(d, \mathcal{W})\|F\|^{2}\left(1+|z|^{2}\right)^{d+(k+1) N_{2}}
$$

Here $N_{2} \geq 1$ is a constant depending only on the constants defining $\mathcal{W}$, while $C_{k}(d, \mathcal{W})$ depends only on $k, d$ and the constants defining $\mathcal{W}$.

Proof of Lemma 4.2. - $\quad$ Since $\varphi=1$ on $T_{\gamma^{\prime}, m^{\prime}}(\mathcal{Z}(P))$ and $g_{j}$ is holomorphic on this set, (i) is obvious.

For (ii), if $z \in \mathcal{Z}(P)$, then $\Phi_{j}(z)=g_{j}(z)$ and (ii) follows from (15). If $z \in T_{\gamma, m}(\mathcal{Z}(P)) \backslash \mathcal{Z}(P)$, we have

$$
\begin{aligned}
\sum \Phi_{j}(z) p_{j}(z) & =\varphi(z) \sum g_{j}(z) p_{j}(z)+(1-\varphi(z)) \sum \frac{\overline{p_{j}(z)} F(z)}{|P(z)|^{2}} p_{j}(z) \\
& =F(z)
\end{aligned}
$$

Finally, if $z \notin T_{\gamma, m}(\mathcal{Z}(P))$, we have

$$
\sum \Phi_{j}(z) p_{j}(z)=\sum \frac{\overline{p_{j}(z)} F(z)}{|P(z)|^{2}} p_{j}(z)=F(z)
$$

this proves (ii).
To prove (iii), we make the following remarks. We have $\operatorname{supp}(\varphi) \subset$ $T_{\gamma^{\prime \prime}, m^{\prime \prime}}(\mathcal{Z}(P))$, and $B\left(w, \gamma^{\prime \prime}\left(1+|z|^{2}\right)^{-m^{\prime \prime}}\right) \subset T_{\gamma, m}(\mathcal{Z}(P))$ if $w \in$ $T_{\gamma^{\prime \prime}, m^{\prime \prime}}(\mathcal{Z}(P))$. From the estimates for $g_{j}$ on $T_{\gamma, m}(\mathcal{Z}(P))$ given in (15) and Cauchy's inequalities, we deduce that

$$
\left|D^{(k)} g_{j}(z)\right|^{2} \leq C_{k}^{\prime}(d, \mathcal{W})\|F\|^{2}\left(1+|z|^{2}\right)^{d+\bar{m}+2 k m^{\prime \prime}}
$$

for $z \in T_{\gamma^{\prime \prime}, m^{\prime \prime}}(\mathcal{Z}(P))$. Now, if $z \notin T_{\gamma^{\prime}, m^{\prime}}(\mathcal{Z}(P))$ (in particular, if $z \in$ $\operatorname{supp}(1-\varphi)$ ), we have $\left|D^{(k)} \frac{\bar{p}_{j}}{|P|^{2}}(z)\right|^{2} \leq C_{k}(\mathcal{W})\left(1+|z|^{2}\right)^{(k+1) m_{4}}$ because of (16) ( $m_{4}$ depending only on the constants in $\mathcal{W}$ ). These estimates, and the estimates for $D^{(k)} \varphi$ in the defining properties of $\varphi$ imply (iii).

We can now use the method of proof of Theorem 3.1 to complete that of Theorem 4.1.

We consider again the space $\mathbb{C}^{r}$ with standard basis $e_{1}, \ldots, e_{r}$. As in $\S 3$, let $E_{p}=\wedge{ }_{\wedge}^{r} \mathbb{C}^{r} \mathcal{A}^{q}\left(E_{p}\right)$ is the space of smooth $(0, q)$ forms on $\mathbb{C}^{n}$ with values in $E_{p}$, and, for $P \in \mathcal{W}, \mathcal{A}_{0}^{q}\left(E_{p}\right)$ is the subspace of forms with support $\subset \mathbb{C}^{n} \backslash \mathcal{Z}(P)$. We consider again the Koszul complex

$$
\begin{equation*}
0 \rightarrow \mathcal{A}^{q}\left(E_{r}\right) \xrightarrow{\kappa} \mathcal{A}^{q}\left(E_{r-1}\right) \rightarrow \cdots \rightarrow \mathcal{A}^{q}\left(E_{1}\right) \xrightarrow{\kappa} \mathcal{A}^{q}\left(E_{0}\right) \rightarrow 0 \tag{17}
\end{equation*}
$$

with $\kappa$ defined by

$$
\begin{align*}
& \kappa\left(e_{I} \otimes d \bar{z}^{J}\right)=\sum_{k=1}^{p}(-1)^{k-1} p_{i_{k}}(z) e_{I \backslash\left\{i_{k}\right\}} \otimes d \bar{z}^{J}  \tag{18}\\
& I=\left(i<\cdots<i_{p}\right), J=\left(j_{1}<\cdots<j_{q}\right) .
\end{align*}
$$

We also have the homotopy operator

$$
h: \mathcal{A}_{0}^{q}\left(E_{p}\right) \rightarrow \mathcal{A}_{0}^{q}\left(E_{p+1}\right): h(\omega)=\left(\sum_{j=1}^{r} \frac{\bar{p}_{j}}{|P|^{2}} e_{j}\right) \wedge \omega
$$

( $h=0$ on $E_{r}$ ) with the property that

$$
\left\{\begin{array}{l}
\kappa h=\text { identity on } \mathcal{A}_{0}^{q}\left(E_{0}\right)  \tag{19}\\
\kappa h+h \kappa=\text { identity on } \mathcal{A}_{0}^{q}\left(E_{p}\right), p>0 .
\end{array}\right.
$$

(See §3).
Let $\left(\Phi_{1}, \ldots, \Phi_{r}\right)$ be as in Lemma 4.2. Define

$$
\begin{equation*}
\Phi^{(1)}=\sum_{j} \Phi_{j} e_{j} \in \mathcal{A}^{0}\left(E_{1}\right) \tag{20}
\end{equation*}
$$

We have $\kappa\left(\Phi^{(1)}\right)=\sum \Phi_{j} p_{j}=F$. Consider $\bar{\partial} \Phi^{(1)} \in \mathcal{A}^{1}\left(E_{1}\right)$. By property (i) of the $\Phi_{j}, \bar{\partial} \Phi^{(1)} \in \mathcal{A}_{0}^{1}\left(E_{1}\right)$ (i.e. $\left.\operatorname{supp}\left(\bar{\partial} \Phi^{(1)}\right) \cap \mathcal{Z}(P)=\emptyset\right)$. Using the fact that $\bar{\partial}$ and $h$ do not increase supports, we define, successively,

$$
\begin{aligned}
& \Phi^{(2)}=h\left(\bar{\partial} \Phi^{(1)}\right) \in \mathcal{A}_{0}^{1}\left(E_{2}\right) \\
& \Phi^{(3)}=h\left(\bar{\partial} \Phi^{(2)}\right) \in \mathcal{A}_{0}^{2}\left(E_{3}\right), \\
& \vdots \\
& \Phi^{(r)}=h\left(\bar{\partial} \Phi^{(r-1)}\right) \in \mathcal{A}_{0}^{r-1}\left(E_{r}\right) .
\end{aligned}
$$

We have $\kappa\left(\bar{\partial} \Phi^{(1)}\right)=\bar{\partial} \kappa\left(\Phi^{(1)}\right)=\bar{\partial} F=0$ so that $\kappa\left(\Phi^{(2)}\right)=\kappa h\left(\bar{\partial} \Phi^{(1)}\right)=$ $(\kappa h+h \kappa)\left(\bar{\partial} \Phi^{(1)}\right)=\bar{\partial} \Phi^{(1)}$. We obtain, by induction:

$$
\left\{\begin{array}{l}
\kappa\left(\bar{\partial} \Phi^{(p)}\right)=0 \text { for } 1 \leq p \leq r  \tag{21}\\
\kappa\left(\Phi^{(p+1)}\right)=\bar{\partial} \Phi^{(p)} \text { for } 1 \leq p \leq r-1
\end{array}\right.
$$

By Lemma 3.5 and Lemma 4.2, Part (iii), we have the following:
(22) For $k \geq 0, d \geq 1$, there exists a constant $C_{k}(d, \mathcal{W})$ depending only on $k, d$ and the constants in $\mathcal{W}$, and a constant $N_{3} \geq 1$ depending only on the constants in $\mathcal{W}$, such that, for $1 \leq p \leq r$,

$$
\left|D^{(k)} \Phi^{(p)}(z)\right|^{2} \leq C_{k}(d, \mathcal{W})\|F\|^{2}\left(1+|z|^{2}\right)^{d+(k+1) N_{3}}, k \geq 0, z \in \mathbb{C}^{n}
$$

$D^{(k)}$ being, as usual, any differentiation of order $k$.

In the rest of the proof, we denote by $\Gamma_{*}$ a constant depending only on those in $\mathcal{W}$, and by $C_{*}(d)$, a constant depending only on $d$ and the constants in $\mathcal{W}$; they are not necessarily the same at each occurrence.

Let $N_{3}$ be as in (22), and define $\mu_{r-1}=d+N_{3}+n+1, \mu_{q}=\mu_{q+1}+D+2$ for $r-2 \geq q \geq 1$.

We have $\Phi^{(r)} \in \mathcal{A}_{0}^{r-1}\left(E_{r}\right), \kappa\left(\bar{\partial} \Phi^{(r)}\right)=0$ (by (21)). Hence, since $\kappa$ is injective on $\mathcal{A}_{0}^{q}\left(E_{r}\right)$ (see $\S 3,(17)$ ), we have $\bar{\partial} \Phi^{(r)}=0$. Moreover, by (22),

$$
\int_{\mathbb{C}^{n}}\left|\Phi^{(r)}\right|^{2}\left(1+|z|^{2}\right)^{-\mu_{r-1}} d \lambda \leq C_{*}(d)\|F\|^{2}
$$

Hence, by Hörmander's theorem (Cor. 3.4), there is $\Psi^{(r-1)} \in \mathcal{A}^{r-2}\left(E_{r}\right)$ such that

$$
\bar{\partial} \Psi^{(r-1)}=\Phi^{(r)}, \int_{\mathbb{C}^{n}}\left|\Psi^{(r-1)}\right|^{2}\left(1+|z|^{2}\right)^{-\mu_{r-1}-2} d \lambda \leq C_{*}(d)\|F\|^{2}
$$

Since, for any $\Psi \in \mathcal{A}^{q}\left(E_{p}\right)$, we have, by $\S 3$, (18),

$$
|\kappa(\Psi)(z)|^{2} \leq \Gamma_{*}\left(1+|z|^{2}\right)^{D}|\Psi(z)|^{2}
$$

we have also

$$
\int_{\mathbb{C}^{n}}\left|\Phi^{(r-1)}-\kappa\left(\Psi^{(r-1)}\right)\right|^{2}\left(1+|z|^{2}\right)^{-\mu_{r-1}-2-D} d \lambda \leq C_{*}(d)\|F\|^{2}
$$

Moreover, $\Phi^{(r-1)}-\kappa\left(\Psi^{(r-1)}\right) \in \mathcal{A}^{r-2}\left(E_{r-1}\right)$ and $\bar{\partial}\left(\Phi^{(r-1)}-\kappa\left(\Psi^{(r-1)}\right)\right)=$ $\bar{\partial} \Phi^{(r-1)}-\kappa \bar{\partial} \Psi^{(r-1)}=\bar{\partial} \Phi^{(r-1)}-\kappa\left(\Phi^{(r)}\right)=0$ by (21); hence we can find $\Psi^{(r-2)} \in \mathcal{A}^{r-3}\left(E_{r-1}\right)$ so that $\left.\bar{\partial} \Psi^{(r-2)}\right)=\Phi^{(r-1)}-\kappa\left(\Psi^{(r-1)}\right)$ and satisfying

$$
\int_{\mathbb{C}^{n}}\left|\Psi^{(r-2)}\right|^{2}\left(1+|z|^{2}\right)^{-\mu_{r-2}-2} d \lambda \leq C_{*}(d)\|F\|^{2}
$$

Thus, we solve, successively, the equations

$$
\begin{aligned}
\bar{\partial} \Psi^{(r-1)} & =\Phi^{(r)} \\
\bar{\partial} \Psi^{(r-2)} & =\Phi^{(r-1)}-\kappa\left(\Psi^{(r-1)}\right) \\
& \vdots \\
\bar{\partial} \Psi^{(1)} & =\Phi^{(2)}-\kappa\left(\Psi^{(2)}\right)
\end{aligned}
$$

where $\Psi^{(q)} \in \mathcal{A}^{q-1}\left(E_{q+1}\right)(r-1 \geq q \geq 1)$ and

$$
\begin{equation*}
\int_{\mathbb{C}^{n}}\left|\Psi^{(q)}\right|^{2}\left(1+|z|^{2}\right)^{-\mu_{q}-2} d \lambda \leq C_{*}(d)\|F\|^{2} \tag{23}
\end{equation*}
$$

Consider now $G=\Phi^{(1)}-\kappa\left(\Psi^{(1)}\right) \in \mathcal{A}^{0}\left(E_{1}\right)$. We have

$$
\begin{aligned}
\bar{\partial} G=\bar{\partial} \Phi^{(1)}-\kappa\left(\bar{\partial} \Psi^{(1)}\right) & =\bar{\partial} \Phi^{(1)}-\kappa\left(\Phi^{(2)}-\kappa\left(\Psi^{(2)}\right)\right) \\
& =\bar{\partial} \Phi^{(1)}-\kappa\left(\Phi^{(2)}\right)=0 \text { by }(21)
\end{aligned}
$$

Thus $G$ is holomorphic. Further $\kappa(G)=\kappa\left(\Phi^{(1)}\right)=F$, so that, if $G=$ $\sum G_{j} e_{j}$, we have $F=\sum G_{j} p_{j}$. The estimate (23) for $\Psi^{(1)}$, and (22) and §3, (18), imply that

$$
\int_{\mathbb{C}^{n}}\left|G_{j}\right|^{2}\left(1+|z|^{2}\right)^{-\mu_{1}-2-D} d \lambda \leq C_{*}(d)\|F\|^{2}
$$

Using the fact that $\left|G_{j}\right|^{2}$ is subharmonic as at the end of the proof of Theorem 3.1, we conclude that the $G_{j}$ are polynomials of degree $\leq \mu_{1}+2+D-n \leq d+D^{\prime}$ (with $D^{\prime}$ depending only on $D, n$ and $N_{3}$ in (22)) and that

$$
\left\|G_{j}\right\|^{2} \leq C_{*}(d)\|F\|^{2}
$$

This completes the proof of Theorem 4.1.

## 5. The auxiliary function.

Let $n \geq 2, D \geq 1,1 \leq r \leq n-1, C_{1}>0$ be given. Let $\mathcal{W}_{0}$ be the space of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right)$ with $p_{j} \in \mathcal{H}^{D}$ (i.e. $p_{j} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$, $\left.\operatorname{deg} p_{j} \leq D\right)$ such that: $P(0)=0,\|P\| \leq C_{1}, J_{P}(0)=1$, where, as before, $J_{P}(z)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}\right)_{1 \leq j, k \leq r}$.

Theorem 5.1. - There exist constants $D^{\prime} \geq 1, C^{\prime}>0$ depending only on $n, D, C_{1}$, and, for $d \geq 1$, a constant $C(d)$ depending only on $d, n, D, C_{1}$ such that the following holds.

Let $P \in \mathcal{W}_{0}$. There exists $Q \in \mathcal{H}^{D^{\prime}}$ with $Q(0)=1,\|Q\| \leq C^{\prime}$ such that if $F \in \mathcal{H}^{d}$ and $F$ vanishes on $\mathcal{Z}(P) \cap\left\{z \in \mathbb{C}^{n}| | z \mid<\epsilon\right\}$ for some $\epsilon>0$, then we can find polynomials $F_{1}, \ldots, F_{r}$ having the following properties:
(i) $J_{P}^{d} Q F=F_{1} p_{1}+\cdots+F_{r} p_{r}$.
(ii) $\operatorname{deg} F_{j} \leq n D\left(d+D^{\prime}\right)$.
(iii) $\left\|F_{j}\right\| \leq C(d)\|F\|$.

Proof. - As at the beginning of $\S 2$, we introduce a new variable $z_{0}$ and a polynomial $f_{0}\left(z_{0}, \ldots, z_{n}\right)=\left(z_{0}+1\right) J_{P}(z)-1, z=\left(z_{1}, \ldots, z_{n}\right)$. We write $w=\left(z_{0}, z\right)=\left(z_{0}, z_{1}, \ldots, z_{n}\right) \in \mathbb{C}^{n+1}$ and $f_{j}(w)=p_{j}(z)$ for $j \geq 1, f_{0}(w)=f_{0}\left(z_{0}, \ldots, z_{n}\right)$.

We have

$$
J_{f}(w)=\operatorname{det}\left(\frac{\partial f_{i}}{\partial w_{j}}\right)_{0 \leq i, j \leq r}=\left(J_{p}(z)\right)^{2}
$$

(1) There exist constants $c_{0}>0, N \geq 1, C_{2} \geq 1, D_{1} \geq 1$ depending only on $n, D, C_{1}$, such that if $\mathcal{W}$ is the space of $(r+1)$-tuples $q=\left(q_{0}, \ldots, q_{r}\right)$ of polynomials in $(n+1)$ variables $w_{0}, \ldots, w_{n}$ satisfying: $\operatorname{deg} q_{j} \leq D_{1}$, $q(0)=0, J_{q}(0)=1,\|q\| \leq C_{2}$ and

$$
\left|J_{q}(w)\right|^{2}+\left|q_{0}(w)\right|^{2}+\cdots+\left|q_{r}(w)\right|^{2} \geq c_{0}\left(1+|w|^{2}\right)^{-N}, w \in \mathbb{C}^{n+1}
$$

then, the following holds.
If $P=\left(p_{1}, \ldots, p_{r}\right) \in \mathcal{W}_{0}$, the $(r+1)$-tuple $f=\left(f_{0}, \ldots, f_{r}\right)$ defined by

$$
f_{0}(w)=\left(z_{0}+1\right) J_{P}(z)-1, f_{j}(w)=p_{j}(z), j \geq 1, w=\left(z_{0}, \ldots, z_{n}\right)
$$

belongs to $\mathcal{W}$. [We can take $D_{1}=D+n(D-1)$ since $\operatorname{deg} f_{0} \leq 1+r(D-1)$.]
This was noted in §2, and follows easily from the relation

$$
\left(1+\left(z_{0}+1\right) J_{P}(z)\right) f_{0}(w)-\left(z_{0}+1\right)^{2} J_{f}(w)=-1
$$

Let $X_{1}, \ldots, X_{t}$ be the connected components of $\mathcal{Z}(f) \subset \mathbb{C}^{n+1}, X_{1}$ being the component containing 0 . By Theorem 3.1 (applied with $n, r, D$, $C_{1}, c_{0}, N$ replaced by $\left.n+1, r+1, D_{1}, C_{2}, c_{0}, N\right)$ there exists a polynomial $\Phi \in \mathbb{C}[w]$ such that

$$
\begin{equation*}
\Phi\left|X_{1}=1, \Phi\right| X_{j}=0 \text { for } j>1, \operatorname{deg} \Phi \leq D_{2},\|\Phi\| \leq C_{3} \tag{2}
\end{equation*}
$$

where $D_{2}, C_{3}$ depend only on the constants defining $\mathcal{W}$ above, hence only on $n, D, C_{1}$.

Let $F \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ and suppose that $\operatorname{deg} F \leq d$ and that $F$ vanishes on $\mathcal{Z}(P) \cap\left\{z \in \mathbb{C}^{n}| | z \mid<\epsilon\right\}$ for some $\epsilon>0$. Then $F$, considered as a function of $w=\left(z_{0}, z\right)$ vanishes on a neighborhood of 0 on $X_{1}$, so that, by the principle of analytic continuation, $F \mid X_{1} \equiv 0$. Thus $\Phi \cdot F \equiv 0$ on $\mathcal{Z}(f)$. By Theorem 4.1, we obtain the following:
(3) There exist polynomials $G_{0}(w), \ldots, G_{r}(w)$ with the following properties:
(a) $\Phi F=G_{0} f_{0}+\cdots+G_{r} f_{r}$
(b) $\operatorname{deg} G_{j} \leq d+D_{3}$
(c) $\left\|G_{j}\right\| \leq C_{1}(d)\|\Phi F\|$;
here $D_{3}$ depends only on $n, D, C_{1}$, and $C_{1}(d)$ on $d, n, D, C_{1}$.
Since $\|\Phi\| \leq C_{3}$, (c) above implies
(c') $\left\|G_{j}\right\| \leq C_{2}(d)\|F\|, C_{2}(d)$ depending only on $d, n, D, C_{1}$.

Let $D^{\prime \prime}=\max \left(D_{2}, D_{3}\right)$ with $D_{2}$ as in (2), $D_{3}$ as in (3, b), and set

$$
\left\{\begin{array}{l}
Q\left(z_{1}, \ldots, z_{n}\right)=\left(J_{P}(z)\right)^{D^{\prime \prime}} \Phi\left(-1+\frac{1}{J_{P}(z)}, z_{1}, \ldots, z_{n}\right)  \tag{4}\\
F_{j}\left(z_{1}, \ldots, z_{n}\right)=\left(J_{P}(z)\right)^{d+D^{\prime \prime}} G_{j}\left(-1+\frac{1}{J_{P}(z)}, z_{1}, \ldots, z_{n}\right), j \geq 1
\end{array}\right.
$$

Since $\operatorname{deg} G_{j} \leq d+D^{\prime \prime}, F_{j}$ is a polynomial of degree $\leq d+D^{\prime \prime}+$ $\left(d+D^{\prime \prime}\right) \operatorname{deg} J_{P} \leq\left(d+D^{\prime \prime}\right)+(1+r(D-1)) \leq n D\left(d+D^{\prime \prime}\right)$. Moreover, if $G_{j}(w)=\sum_{|\alpha| \leq d+D^{\prime \prime}} c_{\alpha} w^{\alpha}$, then $c_{\alpha}\left(J_{P}(z)\right)^{d+D^{\prime \prime}}\left(-1+\frac{1}{J_{P}(z)}\right)^{\alpha_{0}} z_{1}^{\alpha_{1}} \cdots z_{n}^{\alpha_{n}}$ clearly has norm at most $\left|c_{\alpha}\right| C_{3}(d), C_{3}(d)$ depending only on $d, n, D, C_{1}$. Thus

$$
\begin{equation*}
\operatorname{deg} F_{j} \leq n D\left(d+D^{\prime \prime}\right),\left\|F_{j}\right\| \leq C_{4}(d)\left\|G_{j}\right\| \leq C(d)\|F\| \tag{5}
\end{equation*}
$$

(by ( $c^{\prime}$ ) above). In the same way

$$
\left\{\begin{array}{l}
\operatorname{deg} Q \leq D_{2}+D^{\prime \prime} \cdot n(D-1)=D^{\prime}(\text { say })  \tag{6}\\
\|Q\| \leq C^{\prime}, C^{\prime} \text { depending only on } n, D, C_{1}
\end{array}\right.
$$

Further $Q(0)=J_{P}(0)^{D^{\prime \prime}} \Phi(0)=1$ (since $\left.-1+\frac{1}{J_{P}(0)}=0\right)$. Finally, $f_{0}\left(-1+\frac{1}{J_{P}(z)}, z_{1}, \ldots, z_{n}\right) \equiv 0$, so that

$$
\begin{align*}
&\left(J_{P}(z)\right)^{d} Q(z) F(z)  \tag{7}\\
&=\left(J_{P}(z)\right)^{d+D^{\prime \prime}} \sum_{j=1}^{r} G_{j}\left(-1+\frac{1}{J_{P}(z)}, z_{1}, \ldots, z_{n}\right) p_{j}\left(z_{1}, \ldots, z_{n}\right) \\
&=\sum_{j=1}^{r} F_{j}(z) p_{j}(z)
\end{align*}
$$

Theorem 5.1 follows from (5), (6), (7).
Theorem 5.1 has a real analogue which is, in fact, what we shall need.
Theorem 5.2. - Let $n \geq 2, D \geq 1,1 \leq r \leq n-1$, and $C_{1}>0$ be given. Let $W$ be the space of r-tuples $P=\left(p_{1}, \ldots, p_{r}\right)$, where $p_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right], \operatorname{deg} p_{j} \leq D$ (i.e. $p_{j} \in H^{D}$ ), satisfying $P(0)=0$, $\|P\| \leq C_{1}, J_{P}(0)=1$, with $J_{P}(x)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}\right)_{1 \leq j, k \leq r}$.

There exist constants $D^{\prime} \geq 1, C^{\prime}>0$ depending only on $n, D, C_{1}$ and, for $d \geq 1$ a function $C(d)$ of $d, n, D, C_{1}$ such that the following holds.

Let $P \in W$ and $Z(P)=\left\{z \in \mathbb{R}^{n} \mid p_{1}(x)=\cdots p_{r}(z)=0\right\}$. There exists $Q \in H^{D^{\prime}}$ with $Q(0)=1$ and $\|Q\| \leq C^{\prime}$ such that: if $F \in H^{d}$ and
$F=0$ on $Z(P) \cap\left\{x \in \mathbb{R}^{n}| | x \mid<\epsilon\right\}$ for some $\epsilon>0$, then we can find $F_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right], j=1, \ldots, r$ having the following properties:
(i) $\left(J_{P}(x)\right)^{d} Q(x) F(x)=\sum_{j=1}^{r} F_{j}(x) p_{j}(x), x \in \mathbb{R}^{n}$.
(ii) $\operatorname{deg} F_{j} \leq n D\left(d+D^{\prime}\right)$.
(iii) $\left\|F_{j}\right\| \leq C(d)\|F\|$.

Proof. - The implicit function theorem implies that there are holomorphic functions $h_{j}(\zeta), j=1, \ldots, r, \zeta=\left(z_{r+1}, \ldots, z_{n}\right)$ defined in a ball $B(0, \rho) \subset \mathbb{C}^{n-r}$, which are real for real values of $\zeta$ such that for arbitrarily small neighborhoods $U$ of 0 in $\mathbb{C}^{n}$, we have

$$
\begin{aligned}
\mathcal{Z}(P) \cap U & =\left\{\left(h_{1}(\zeta), \ldots, h_{r}(\zeta), \zeta\right)\left|\zeta \in \mathbb{C}^{n-r},|\zeta|<\rho\right\} \cap U\right. \\
Z(P) \cap U & =\mathcal{Z}(P) \cap \mathbb{R}^{n} \cap U \\
& =\left\{\left(h_{1}(\zeta), \ldots, h_{r}(\zeta), \zeta\right)\left|\zeta \in \mathbb{R}^{n-r},|\zeta|<\rho\right\} \cap U .\right.
\end{aligned}
$$

Hence, if $F \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ and $F=0$ on $Z(P) \cap U$, then $F\left(h_{1}(\zeta), \ldots, h_{r}(\zeta), \zeta\right)$ vanishes for real $\zeta$ sufficiently close to 0 , hence also for all small complex $\zeta$. Hence $F \mid \mathcal{Z}(P)$ is zero in some neighborhood of 0 . Theorem 5.1 gives us polynomials $Q, F_{j}$ with complex coefficients satisfying (i), (ii), (iii) above. Replacing them by the real polynomials whose coefficients are the real parts of the $Q, F_{j}$, we obtain the desired polynomials.

Remark. - For the application of these results to the proof of the extension theorem, we need the bounds given on $Q$, and the bounds on the degrees of the $F_{j}$ (not those on the norms). These bounds on the degrees of the $F_{j}$ can be obtained by using a purely algebraic theorem, although the construction of $Q$ with the bounds given above seems to necessitate analytic methods.

We are grateful to Burt Totaro who told us that the next theorem was known. The reference to work of Bayer and Stillman that he gave us led us to the paper [ He ] of Grete Hermann in which it is proved. We formulate this theorem, and indicate its proof, although our argument is not very different from that of Hermann. The theorem and its proof remain valid if $\mathbb{R}$ is replaced by $\mathbb{C}$.

THEOREM 5.4. - Let $n, r, s, D, d$ be integers $\geq 1$. There is an integer $D^{\prime} \geq 1$ depending only on $n, r, D, d$ such that the following holds.

Let $p_{i j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right], 1 \leq i \leq r, 1 \leq j \leq s$, be polynomials with $\operatorname{deg} p_{i j} \leq D . \operatorname{Let} F_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right], j=1, \ldots, s$, and suppose that $\operatorname{deg} F_{j} \leq d$.

Assume that the system of equations

$$
\begin{equation*}
F_{j}=\sum_{i=1}^{r} g_{i} p_{i j}, j=1, \ldots, s \tag{8}
\end{equation*}
$$

has a solution $g=\left(g_{1}, \ldots, g_{r}\right)$ with $g_{i} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$.
Then, there is a solution of (8) with $\operatorname{deg} g_{i} \leq D^{\prime}$.

Outline of Proof. - The proof is by induction on $n$; the result is trivial for $n=0$. Assume the theorem proved for $\mathbb{R}\left[x_{1}, \ldots, x_{n-1}\right]$.

We may assume that the vectors $P_{j}=\left(p_{1 j}, \ldots, p_{r j}\right), j=1, \ldots, s$ are linearly independent over the field $\mathbb{R}(x)$ of rational functions. In fact, if $P_{1}, \ldots, P_{s^{\prime}}, s^{\prime} \leq s$, are independent, and $P_{j}, j>s^{\prime}$, are linear combinations of these, then it suffices to solve the equations

$$
F_{j}=\sum_{i=1}^{r} g_{i} p_{i j}, j=1, \ldots, s^{\prime}
$$

the other equations in (8) necessarily follow from these. In particular, we may assume that $s \leq r$.

By a change of notation, we may assume that the matrix

$$
P=\left(p_{i j}\right)_{1 \leq i \leq s, 1 \leq j \leq s}
$$

has a non-zero determinant $\Delta \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$. Its degree is $\leq s D, \Delta \not \equiv 0$. The matrix $\Delta P^{-1}$ is a matrix of polynomials of degree $\leq(s-1) D$ (its entries are, up to sign, the determinants of the matrices obtained from $P$ by deleting one row and one column).

Multiplying the equation

$$
\left(F_{1}, \ldots, F_{s}\right)=\left(g_{1}, \ldots, g_{r}\right)\left(p_{i j}\right)
$$

on the right by $\Delta P^{-1}$, we see that, with our assumptions, (8) is equivalent to a system

$$
G_{1}=g_{1} \Delta+\sum_{s<i \leq r} g_{i} q_{i 1}
$$

$$
\begin{equation*}
G_{s}=g_{s} \Delta+\sum_{s<i \leq r} g_{i} q_{i s} \tag{9}
\end{equation*}
$$

where the $q_{i j}$ are linear combinations of the $p_{i j}$, the coefficients being entries of $\Delta P^{-1}$. Hence $\operatorname{deg} q_{i j} \leq s D$. The $G_{j}$ are linear combinations of the $F_{j}$ with coefficients entries of $\Delta P^{-1}$, so that $\operatorname{deg} G_{j} \leq d+(s-1) D$.

By a linear change of coordinates, we may assume that if $\operatorname{deg} \Delta=N$ $(\leq s D)$, then the coefficient of $x_{n}^{N}$ is a non-zero constant (i.e. the coefficient of the monomial $x_{1}^{0} \cdots x_{n-1}^{0} x_{n}^{N}$ is $\neq 0$ ).

If $\left(g_{1}, \ldots, g_{r}\right)$ is a solution of (9) and $h \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$, the vector

$$
\left(g_{1}+h q_{s+1,1}, g_{2}+h q_{s+1,2}, \ldots, g_{s}+h q_{s+1, s}, g_{s+1}-h \Delta, g_{s+2}, \ldots, g_{r}\right)
$$

is also a solution of (9). By our assumption on the coefficient of $x_{n}^{N}$ in $\Delta$, we can choose $h$ so that $\operatorname{deg}_{x_{n}}\left(g_{s+1}-h \Delta\right)<N$ ( $\operatorname{deg}_{x_{n}}$ being the degree in $\left.x_{n}\right)$. Repeating the process with $g_{s+2}, \ldots, g_{r}$, we obtain:

If (9) has a solution, it has one in which

$$
\operatorname{deg}_{x_{n}}\left(g_{i}\right)<N \text { for } s<i \leq r
$$

Now, if $\left(g_{1}, \ldots, g_{r}\right)$ satisfies this condition, we can use (9) to conclude that for $1 \leq j \leq s$,

$$
\begin{aligned}
\operatorname{deg}_{x_{n}}\left(g_{j} \Delta\right) & \leq \max \left(\operatorname{deg}_{x_{n}} G_{j}, \operatorname{deg}_{x_{n}} q_{i j}+N-1\right) \\
& \leq \max (d+(s-1) D, 2 s D-1)<d+2 s D
\end{aligned}
$$

Thus, if (9) has a solution, it has one, $\left(g_{1}, \ldots, g_{r}\right)$, with $\operatorname{deg}_{x_{n}} g_{i}<d+2 s D$ for $i=1, \ldots, r$.

We now simply write $G_{j}, g_{i}, q_{i j}, \Delta$, in the form

$$
\sum_{0 \leq \nu<d+2 s D} a_{\nu}\left(x_{1}, \ldots, x_{n-1}\right) x_{n}^{\nu}, a_{\nu} \in \mathbb{R}\left[x_{1}, \ldots, x_{n-1}\right]
$$

equating coefficients of powers of $x_{n}$ in the system (9), we obtain a system of equations of the same form as (8), but in the variables $x_{1}, \ldots, x_{n-1}$, and we can proceed by induction.

Theorem 5.2 can be improved in one respect at the cost of losing the bounds in Theorem 5.2 (iii); we do not, however lose the bounds on $Q$ (which are essential for us in what follows). Since this provides what appears to be the right analogue of the factorization into the main factor and the other factor that we used in [FN], we shall now discuss this.

The following result is what we called Theorem 3 in the Introduction.
Theorem 5.5. - Let $n, D, r, C_{1}$ and $W$ be as in Theorem 5.2. There exist constants $D^{\prime} \geq 1, C^{\prime}>0$ depending only on these data, and, for $d \geq 1$, a function $\bar{d}$ of $d, n, D, r$ and $C_{1}$ such that the following holds.

Given $P=\left(p_{1}, \ldots, p_{r}\right) \in W$, there exists $Q \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ (depending only on $P$ ) with
(a) $Q(0)=1, \operatorname{deg} Q \leq D^{\prime},\|Q\| \leq C^{\prime}$
such that
(b) If $F \in H^{d}$ and $F=0$ as a germ in $\mathcal{G}(P)$ (i.e. $F$ vanishes on some neighborhood of 0 in $V(P))$, then there exist $F_{1}, \ldots, F_{r} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ with

$$
\operatorname{deg} F_{j} \leq \bar{d} \text { and } Q \cdot F=\sum_{j=1}^{r} F_{j} \cdot p_{j}
$$

The proof is based on Theorem 5.6 below (which is simply the real analogue of the following theorem).

Theorem 5.6. - Let $n \geq 2, D \geq 1,1 \leq r \leq n-1$ be given integers. Denote by $\mathcal{W}_{1}$ the space of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right), p_{j} \in \mathcal{H}^{D}$, with $P(0)=0, J_{P}(0)=1$.

There exist integers $D_{*}, r_{*} \geq 1$ depending only on $n, D, r$, and for $d \geq 1$, a function $d_{*}$ of $d, n, D, r$ such that:

Given $P \in \mathcal{W}_{1}$, there exist $f_{1}, \ldots, f_{r_{*}} \in \mathcal{H}^{D_{*}}$ with the following properties:
(i) $f_{j}=0$ on some neighborhood of 0 in $\mathcal{Z}(P), 1 \leq j \leq r_{*}$; and
(ii) if $f \in \mathcal{H}^{d}$ and $f=0$ on some neighborhood of 0 in $\mathcal{Z}(P)$, then, we can find $g_{j} \in \mathcal{H}^{d_{*}}, j=1, \ldots, r_{*}$, such that

$$
f=\sum_{1 \leq j \leq r_{*}} g_{j} f_{j}
$$

The real analogue is
Theorem 5.6'. - Let $n \geq 2, D \geq 1,1 \leq r \leq n-1$, and let $W_{1}$ be the space of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right)$ of polynomials $p_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ with $\operatorname{deg} p_{j} \leq D, P(0)=0, J_{P}(0)=1$.

There exist integers $D_{*}, r_{*} \geq 1$ depending only on $n, D, r$ and for $d \geq 1$, a function $d_{*}$ of $d, n, D, r$ such that, given $P \in W_{1}$, we can find polynomials $f_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right], 1 \leq j \leq r_{*}$ with $\operatorname{deg} f_{j} \leq D_{*}$ which vanish on some neighborhood of 0 in $V(P)$ and such that if $f \in H^{d}$ and $f=0$ as a germ in $\mathcal{G}(P)$, then, there exist $g_{j} \in H^{d_{*}}, 1 \leq j \leq r_{*}$ with $f=\sum_{1 \leq j \leq r_{*}} g_{j} f_{j}$.

We first remark that Theorem $5.6^{\prime}$ follows from Theorem 5.6. In fact, if $P=\left(p_{1}, \ldots, p_{r}\right) \in W_{1} \subset \mathcal{W}_{1}$, and $f \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ is such that it vanishes on $\left\{x \in \mathbb{R}^{n} \mid P(x)=0\right\}$ near 0 , then it vanishes on $\left\{z \in \mathbb{C}^{n} \mid P(z)=0\right\}$ near 0 (see proof of Theorem 5.2). Theorem 5.6 gives us polynomials $f_{1}, \ldots, f_{r_{*}} \in \mathcal{H}^{D_{*}}$ vanishing on $\mathcal{Z}(P)$ near 0 and generating all such $f$; we have only to replace $f_{j}$ by the two polynomials $f_{1 j}, f_{2 j}$ whose coefficients are the real and imaginary parts of those of $f_{j}$. [Note that if $x \in \mathbb{R}^{n}$ and $f_{j}(x)=0$, then $f_{1 j}(x)=f_{2 j}(x)=0$; further, if $f=\sum g_{j} f_{j}$ with $g_{j} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ and we write $g_{j}=g_{1 j}+i g_{2 j}$ where $g_{1 j}, g_{2 j}$ are polynomials with real coefficients, and if $f$ also has real coefficients, we have $f=\sum g_{1 j} f_{1 j}-\sum g_{2 j} f_{2 j}$.]

Next, we remark that Theorem $5.6^{\prime}$ and Theorem 5.2 imply Theorem 5.5. To see this, if $P \in W \subset W_{1}$, Theorem 5.2 implies that there is $Q_{0} \in H^{D^{\prime}}, Q_{0}(0)=1,\left\|Q_{0}\right\| \leq C^{\prime}$ such that, with the $f_{j}$ as in Theorem 5.6 ${ }^{\prime}$, we have

$$
J_{P}^{D *} Q_{0} f_{j}=\sum_{k=1}^{r} f_{j k} p_{k}, f_{j k} \in H^{D^{\prime \prime}}, 1 \leq j \leq r_{*}
$$

here $D^{\prime}, C^{\prime}, D^{\prime \prime}, D_{*}$ depend only on $n, D, r, C_{1}$. If $F \in H^{d}$ and vanishes on some neighborhood of 0 in $V(P)$, then

$$
F=\sum_{1 \leq j \leq r_{*}} g_{j} f_{j} \text { with } g_{j} \in H^{d_{*}}
$$

and we obtain

$$
J_{P}^{D *} Q_{0} F=\sum_{k=1}^{r} F_{k} p_{k} \text { with } F_{k}=\sum_{1 \leq j \leq r_{*}} f_{j k} g_{j}
$$

so that $\operatorname{deg} F_{k} \leq d_{*}+D^{\prime \prime}$. We have only to take $Q=J_{P}^{D *} Q_{0}$ to complete Theorem 5.5.

Thus, to prove Theorem 5.5, we have only to prove Theorem 5.6. Before starting on the proof of this theorem, we make some preliminary remarks.

Consider $\mathbb{C}^{n}$ as an open subset of projective space $\mathbb{P}^{n}$ with homogeneous coordinates $\left(z_{0}: \cdots: z_{n}\right), \mathbb{C}^{n}$ being defined by $z_{0} \neq 0$. The hyperplane at $\infty: H=\left\{\left(z_{0}: \cdots: z_{n}\right) \in \mathbb{P}^{n} \mid z_{0}=0\right\}$ can be naturally identified with $\mathbb{P}^{n-1}$.

Let $X$ be an affine algebraic variety, $X \subset \mathbb{C}^{n}$, of pure dimension $k$. Its closure $\bar{X}$ in $\mathbb{P}^{n}$ is a projective variety of pure dimension $k$, and $\operatorname{dim}(\bar{X} \cap H)=k-1$.

Any linear subspace of dimension $n-k$ in $\mathbb{P}^{n}$ meets $\bar{X}$, but there is a linear subspace $L \subset H$ of dimension $n-k-1$ such that $L \cap \bar{X}=\emptyset$.
(These are standard facts; see e.g. Mumford's book [M].)
A linear space $L \subset H$ as above defines a projection $\varphi: \mathbb{P}^{n} \backslash L \rightarrow \mathbb{P}^{k}$ as follows.

Fix a linear subspace of dimension $k$ in $\mathbb{P}^{n}$ disjoint from $L$; we shall simply call it $\mathbb{P}^{k}$. If $w \in \mathbb{P}^{n} \backslash L, w$ and $L$ span an $(n-k)$-dimensional linear subspace which meets $\mathbb{P}^{k}$ in a unique point $\varphi(w)$. The map $\varphi$ sends $w$ to $\varphi(w)$. Now, $\mathbb{C}^{n} \subset \mathbb{P}^{n} \backslash L$, and, with a suitable choice of linear coordinates, $\varphi \mid \mathbb{C}^{n}$ is just the standard projection of $\mathbb{C}^{n}$ onto $\mathbb{C}^{k}$.

We set $H_{0}=\mathbb{P}^{k} \cap H ; H_{0}$ is a hyperplane in $\mathbb{P}^{k}$, and $\mathbb{P}^{k} \backslash H_{0}=\mathbb{C}^{k}$. We have $\varphi(H \backslash L) \subset H_{0}$ and $\varphi^{-1}\left(H_{0}\right) \subset H$ (since, if $w \in H \backslash L$, the span of $w$ and $L$ lies in $H$ ).
(10) Let $X \subset \mathbb{C}^{n}$ be an affine variety of pure dimension $k$, and let $\pi=\varphi \mid X$. Then $\pi(X) \subset \mathbb{C}^{k}=\mathbb{P}^{k} \backslash H_{0} ; \pi: X \rightarrow \mathbb{C}^{k}$ is a proper map and its fibres $\pi^{-1} \pi(x), x \in X$, are finite.

In fact, if $K \subset \mathbb{C}^{k}$ is compact, then $\pi^{-1}(K)=(\varphi \mid \bar{X})^{-1}(K)$ (since $\left.\varphi^{-1}\left(H_{0}\right) \cap \bar{X}=\bar{X} \cap H\right)$ and so is compact since $\bar{X}$ is compact. The fibres $\pi^{-1} \pi(x)$ are compact analytic sets in $X \subset \mathbb{C}^{n}$ and so are finite sets.
(11) There exists an algebraic variety $B \subset \mathbb{C}^{k}$ of dimension $k-1$ (the branch locus) such that, if $\tilde{B}=\pi^{-1}(B)$, then $\tilde{B}$ has dimension $k-1$ and $\pi \mid X \backslash \tilde{B} \rightarrow \mathbb{C}^{k} \backslash B$ is a finite unramified covering; $X \backslash \tilde{B}$ is smooth, and $\pi$ is of maximal rank $k$ at every point of $X \backslash \tilde{B}$.

This again is standard; for a proof, see e.g. [M].
The number of points in $\pi^{-1}(y), y \in \mathbb{C}^{k} \backslash B$, is independent of $y$ (since $\mathbb{C}^{k} \backslash B$ is connected). We call it the degree of $\pi$ and denote it by $\mu$.

Let $P=\left(p_{1}, \ldots, p_{r}\right) \in \mathcal{W}_{1}$ (so that $p_{j} \in \mathcal{H}^{D}, P(0)=0$, $\left.\operatorname{det}\left(\frac{\partial p_{j}}{\partial z_{k}}(0)\right)=1\right)$. Let $\mathcal{V}(P)$ be the irreducible component of $\mathcal{Z}(P)=\{z \in$ $\left.\mathbb{C}^{n} \mid P(z)=0\right\}$ containing the origin (there is only one such component since $\mathcal{Z}(P)$ is smooth at 0 because $\left.J_{P}(0)=1\right)$.

We take $X=\mathcal{V}(P)(\operatorname{dim} X=n-r)$ and project from a linear subspace $L \subset H$ of dimension $n-(n-r)-1=r-1$ not meeting $\bar{X}$. We use the notation introduced above.

Let $\tilde{S}$ be the union of $\tilde{B}$ and the intersection of $\mathcal{V}(P)$ with all other irreducible components of $\mathcal{Z}(P)$, and let $S=\pi(\tilde{S}) ; S$ and $\tilde{S}$ are algebraic
varieties of dimension $\leq n-r-1$.
We shall need the following:
Lemma 5.7. - The degree $\mu$ of $\pi: \mathcal{V}(P) \rightarrow \mathbb{C}^{n-r}$ satisfies $\mu \leq D^{r}$.
The proof uses Bezout's Theorem stated below. For a proof, see van der Waerden, Algebra, vol. 2, Chap. XI, §83.

Bezout's Theorem. - Let $f_{1}, \ldots, f_{n}$ be homogeneous polynomials of degrees $d_{1}, \ldots, d_{n}$ respectively in $n+1$ variables $z_{0}, \ldots, z_{n}$. Assume that the set

$$
\left\{z=\left(z_{0}: \cdots: z_{n}\right) \in \mathbb{P}^{n} \mid f_{1}(z)=\cdots=f_{n}(z)=0\right\}
$$

is finite. Then, the number of points in this set is $\leq d_{1} \cdots d_{n}$.
In fact, counted with the proper multiplicities (which are $\geq 1$ ), the sum of the multiplicities of these points equals $d_{1} \cdots d_{n}$.

Proof of Lemma 5.7. - If $p \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right], p=\sum c_{\alpha} z_{1}^{\alpha_{1}} \cdots z_{n}^{\alpha_{n}}$ has degree $d$, we denote by $\tilde{p}$ the "polynomial made homogeneous", viz. $\tilde{p}\left(z_{0}, \ldots, z_{n}\right)=\sum c_{\alpha} z_{0}^{d-|\alpha|} z_{1}^{\alpha_{1}} \cdots z_{n}^{\alpha_{n}}$. If $P=\left(p_{1}, \ldots, p_{r}\right) \in W$, the fibre $\pi^{-1}(c)$ of $\pi: \mathcal{V}(P) \rightarrow \mathbb{C}^{n-r}$ over $c \in \mathbb{C}^{n-r}, c=\left(c_{1}, \ldots, c_{n-r}\right)$ is the intersection of $\mathcal{V}(P)$ with the affine subspace $\varphi\left(z_{1}, \ldots, z_{n}\right)=c$; if the coordinates are chosen so that $\varphi \mid \mathbb{C}^{n}$ is the standard projection $\left(z_{1}, \ldots, z_{n}\right) \mapsto\left(z_{1}, \ldots, z_{n-r}\right), \pi^{-1}(c)=\left\{z \in \mathcal{V}(P) \mid z_{\nu}-c_{\nu}=0\right.$, $\nu=1, \ldots, n-r\}$. Let $\ell_{\nu}\left(z_{0}, \ldots, z_{n}\right)=z_{\nu}-c_{\nu} z_{0}$. Then $\pi^{-1}(c) \subset\{z \in$ $\left.\mathbb{C}^{n} \mid \tilde{p}_{j}\left(1, z_{1}, \ldots, z_{n}\right)=0=\ell_{\nu}\left(1, z_{1}, \ldots, z_{n}\right)\right\}$.

If $c \in \mathbb{C}^{n-r} \backslash S$, the map $\left(z_{1}, \ldots, z_{n}\right) \mapsto\left(p_{1}(z), \ldots, p_{r}(z), \ell_{1}(1, z)\right.$, $\left.\ldots, \ell_{n-r}(1, z)\right)$ is of maximal rank $n$ at every point $z^{(\alpha)}$ of $\pi^{-1}(c), \alpha=$ $1, \ldots, \mu$. Let $U_{\alpha}$ be a small neighborhood of $z^{(\alpha)}, U_{\alpha} \cap U_{\beta}=\emptyset$ if $\alpha \neq \beta$. There exists $\epsilon>0$ such that if $q_{j}$ is a homogeneous polynomial in $z_{0}, \ldots, z_{n}$ of degree $d_{j}=\operatorname{deg} p_{j}$, if $\lambda_{\nu}$ is a homogeneous linear form, and if $\left\|\tilde{p}_{j}-q_{j}\right\|<\epsilon,\left\|\ell_{\nu}-\lambda_{\nu}\right\|<\epsilon$ (the norms being for polynomials in $n+1$ variables), then each $U_{\alpha}$ contains a point at which $q_{j}\left(1, z_{1}, \ldots, z_{n}\right)=0=$ $\lambda_{\nu}\left(1, z_{1}, \ldots, z_{n}\right)(1 \leq j \leq r, 1 \leq \nu \leq n-r)$. Hence: (12) $\quad \mu \leq \#\left\{z=\left(z_{0}: \cdots: z_{n}\right) \in \mathbb{P}^{n} \mid\right.$

$$
\left.q_{1}(z)=\cdots=q_{r}(z)=0=\lambda_{1}(z)=\cdots=\lambda_{n-r}(z)\right\}
$$

whenever $\operatorname{deg} q_{j}=d_{j},\left\|q_{j}-\tilde{p}_{j}\right\|<\epsilon,\left\|\lambda_{\nu}-\ell_{\nu}\right\|<\epsilon$.
We now remark that for any $\epsilon>0$, we can choose $q_{j}, \lambda_{\nu}$ close to $\tilde{p}_{j}, \ell_{\nu}$ so that the set

$$
\begin{equation*}
\left\{z \in \mathbb{P}^{n} \mid q_{1}(z)=\cdots=q_{r}(z)=0=\lambda_{1}(z)=\cdots=\lambda_{n-r}(z)\right\} \tag{13}
\end{equation*}
$$

is finite. In fact, take $q_{1}=\tilde{p}_{1}$ and choose $q_{2},\left\|q_{2}-\tilde{p}_{2}\right\|<\epsilon$, such that $q_{2}$ does not vanish identically on any irreducible component of the set $\left\{z \in \mathbb{P}^{n} \mid q_{1}(z)=0\right\}$. The variety $\left\{z \in \mathbb{P}^{n} \mid q_{1}(z)=q_{2}(z)=0\right\}$ then has dimension $\leq n-2$. Proceeding in this way, we find $q_{j}, \lambda_{\nu}$ close to $\tilde{p}_{j}, \ell_{\nu}$, so that the set $\left\{q_{j}(z)=0=\lambda_{\nu}(z)\right\}$ has dimension 0 , i.e. is finite.

Bezout's Theorem implies that the number of points in (13) is $\leq d_{1}, \cdots d_{r} 1 \cdots 1 \leq D^{r}$. This and (12) prove Lemma 5.7.

Remark. - If we apply this argument to the whole variety $\mathcal{Z}(P)$ with $P \in \mathcal{W}$, we find that the number $t$ of connected components of $\mathcal{Z}(P)$ is $\leq D^{r}$ as stated in the proof of Theorem 3.1.

Proof of Theorem 5.6. - We consider the affine variety $\mathcal{V}(P) \subset \mathbb{C}^{n}$, $P \in \mathcal{W}$, and project onto $\mathbb{C}^{n-r} \subset \mathbb{P}^{n-r}$ from a linear subspace $L \subset H$ (of dimension $r-1), L \cap \overline{\mathcal{V}(P)}=\emptyset$ as above. We retain the notation above. Suppose linear coordinates $\left(x_{1}, \ldots, x_{n-r}, y_{1}, \ldots, y_{r}\right)$ chosen on $\mathbb{C}^{n}$ so that the map $\varphi: \mathbb{P}^{n} \backslash L \rightarrow \mathbb{P}^{n-r}$, restricted to $\mathbb{C}^{n}$, is simply $\varphi(x, y)=x$.

We first prove
Lemma 5.8. - Let $\ell(z)=\ell(x, y), z=(x, y)$ be a linear function on $\mathbb{C}^{n}$ (i.e. a polynomial of degree $\leq 1$ ). For $1 \leq i \leq \mu(\mu=$ degree of $\left.\pi: \mathcal{V}(P) \rightarrow \mathbb{C}^{n-r}\right)$, there are polynomials $b_{i}(x)$ in $x$ of degree $\leq i$ such that
(i) $(\ell(z))^{\mu}+\sum_{1 \leq i \leq \mu} b_{i}(\pi(z))(\ell(z))^{\mu-i}=0$ for any $z \in \mathcal{V}(P)$.
(ii) If $x \in \mathbb{C}^{n-r} \backslash S$ (notation as above), then $b_{i}(x)$ is the $i$-th elementary symmetric function in $\ell\left(z^{(1)}\right), \ldots, \ell\left(z^{(\mu)}\right)$, where $\pi^{-1}(z)=$ $\left(z^{(1)}, \ldots, z^{(\mu)}\right)$.
[Note: the elementary symmetric functions $b_{i}$ in $\xi_{1}, \ldots, \xi_{\mu}$ are defined by $\prod_{\nu=1}^{\mu}\left(T-\xi_{\nu}\right)=T^{\mu}+b_{1} T^{\mu-1}+\cdots+b_{\mu}, T$ being an indeterminate.]

Proof of Lemma 5.8. - For $x \in \mathbb{C}^{n-r} \backslash S$, define $b_{i}(x)$ by (ii) above. The $b_{i}$ are holomorphic on $\mathbb{C}^{n-r} \backslash S$. If $x_{0} \in S$ and $U$ is a bounded open set in $\mathbb{C}^{n-r}$ with $x_{0} \in U$, then $\pi^{-1}(\bar{U})$ is compact in $\mathcal{V}(P)$, so that $\ell$ is bounded on $\pi^{-1}(U)$. Hence $b_{i} \mid U \backslash S$ is bounded, so that it extends to a holomorphic function on $U$ (by the Riemann extension theorem). Since $x_{0} \in S$, is arbitrary this gives us holomorphic functions $b_{i}$ on $\mathbb{C}^{n-r}$.

The equation (i) holds on $\mathcal{V}(P) \backslash \pi^{-1}(S)$ by definition, hence on all $\mathcal{V}(P)$ since $\pi^{-1}(S)$ has dimension $<n-r$ and so is nowhere dense in $\mathcal{V}(P)$.

To prove that $b_{i}$ is a polynomial of degree $\leq i$, let $a \in H_{0}=H \cap \mathbb{P}^{n-r}$; and let $U_{0}$ be a small neighborhood of $a$ in $\mathbb{P}^{n-r}$. We may suppose that there are homogeneous coordinates $\left(w_{0}: \cdots: w_{n-r}\right)$ on $\mathbb{P}^{n-r}$ so that $w_{1} \neq 0$ on $U_{0}$ and $U_{0} \cap H_{0}=\left\{w \in U_{0} \mid w_{0}=0\right\}$. Let $h$ be the holomorphic function $w_{0} / w_{1}$ on $U_{0}$. Since $\ell$ is linear (so has only simple poles on $H$ ), the function $z \mapsto \ell(z) h(\varphi(z))$ is holomorphic on $\varphi^{-1}\left(U_{0}\right) \subset \mathbb{P}^{n} \backslash L$. In particular, if $K_{0}$ is a compact neighborhood of $a, K_{0} \subset U_{0}$, there is a constant $C>0$ with $|\ell(z) h(\pi(z))|<C$ if $z \in \pi^{-1}\left(K_{0} \backslash H_{0}\right)$. Moreover, $\frac{w_{1}}{w_{0}}, \ldots, \frac{w_{n-r}}{w_{0}}$ form linear coordinates in $\mathbb{C}^{n-r}$. Hence, if $x \in K_{0} \backslash H_{0},|h(x)|=1 /\left|\frac{w_{1}}{w_{0}}(x)\right| \geq \frac{C^{\prime}}{1+|x|}$, $C^{\prime}>0$ being a constant. Thus $|\ell(z)| \leq C^{\prime \prime}(1+|\pi(z)|), z \in \pi^{-1}\left(K_{0} \backslash H_{0}\right)$, $C^{\prime \prime}$ being a constant. Since $H_{0}$ is compact, it follows that there is a constant $\bar{C}>0$ so that

$$
\begin{equation*}
|\ell(z)| \leq \bar{C}(1+|\pi(z)|), z \in \mathcal{V}(P) \tag{14}
\end{equation*}
$$

Hence, if $x \in \mathbb{C}^{n-r} \backslash S$ and $\pi^{-1}(x)=\left(z^{(1)}, \ldots, z^{(\mu)}\right)$, we have $\left|b_{i}(x)\right|=$ $\left|(-1)^{i} \sum_{1 \leq \nu_{1}<\cdots<\nu_{i} \leq \mu} \ell\left(z^{\left(\nu_{1}\right)}\right) \cdots \ell\left(z^{\left(\nu_{i}\right)}\right)\right| \leq$ const. $(1+|z|)^{i}$ since $\pi\left(z^{(\nu)}\right)=x$ for $1 \leq \nu \leq \mu$. Since $b_{i}$ is holomorphic on $\mathbb{C}^{n-r}$, this inequality holds everywhere on $\mathbb{C}^{n-r}$ and shows that $b_{i}$ is a polynomial of degree $\leq i$.

Choose an $r \times r$ invertible complex matrix $A$ such that if $\left(u_{1}, \ldots, u_{r}\right)=$ $\left(y_{1}, \ldots, y_{r}\right) A$, then $u_{1}$ separates the points of $\pi^{-1}\left(x_{0}\right)=\left(z_{0}^{(1)}, \ldots, z_{0}^{(\mu)}\right)$ for some $x_{0} \in \mathbb{C}^{n-r} \backslash S_{1}$ (possible since the $r$-tuple ( $y_{1}, \ldots, y_{r}$ ) takes distinct values at the points of $\pi^{-1}\left(x_{0}\right)$ ).

Let $F_{j}\left(x, u_{j}\right)=u_{j}^{\mu}+\sum_{\nu=1}^{\mu} b_{j \nu}(x) y^{\mu-\nu}(1 \leq j \leq r)$ be the polynomial constructed in Lemma 5.8 for the linear form $\ell=u_{j}$. We shall simply write $b_{\nu}$ for $b_{1 \nu}$. We have $F_{j}\left(x, u_{j}\right) \mid \mathcal{V}(P) \equiv 0$. Moreover, if $G\left(x, u_{1}\right)$ is a polynomial in $x, u_{1}$ alone with $\operatorname{deg}_{u_{1}} G<\mu$, then if $G \mid \mathcal{V}(P) \equiv 0$, $G$ must be the zero polynomial (because $u_{1}\left(z^{(\nu)}\right), 1 \leq \nu \leq \mu$, where $\pi^{-1}(x)=\left(z^{(1)}, \ldots, z^{(\mu)}\right)$ would be $\mu$ distinct roots of $G\left(x, u_{1}\right)$ for $x$ near $x_{0}$ ).

Let $\delta=\delta(x)$ be the discriminant of the polynomial $F_{1}\left(x, u_{1}\right)$. Then $\delta \not \equiv 0$; in fact $\delta\left(x_{0}\right) \neq 0$ since $F_{1}\left(x_{0}, u_{1}\right)$ has $\mu$ distinct roots. If $S_{1}=S \cup\left\{x \in \mathbb{C}^{n-r} \mid \delta(x)=0\right\}$, then $\pi^{-1}\left(S_{1}\right)$ has dimension $<n-r$ and, for $x \in \mathbb{C}^{n-r} \backslash S_{1}$, the polynomial $F_{1}\left(x, u_{1}\right)$ has $\mu$ distinct roots, and $u_{1}$ takes distinct values on $\pi^{-1}(x)$. The discriminant $\delta$ can be defined as the determinant of a $(2 \mu-1) \times(2 \mu-1)$ matrix whose entries are the $b_{\nu}(x)$ and $(\mu-\nu) \cdot b_{\nu}(x)$ (the coefficients of $F_{1}$ and $\left.\partial F_{1} / \partial u_{1}\right)$. Hence the degree
of $\delta$ is $\leq \mu(2 \mu-1)$ since $\operatorname{deg} b_{i} \leq i \leq \mu$.
We now use $\left(x_{1}, \ldots, x_{n-r}, u_{1}, \ldots, u_{r}\right)=(x, u)$ as coordinates on $\mathbb{C}^{n}$.
(15) Let $d \geq 1$. There is a constant $C(d, \mu)$ depending only on $d$ and $\mu$ such that the following holds.

If $f \in \mathbb{C}[x, u], \operatorname{deg} f \leq d$, there exist $g_{0}, \ldots, g_{\mu-1} \in \mathbb{C}[x]$, such that $\operatorname{deg} g_{\nu} \leq C(d, \mu)$ and

$$
\delta(x) f(x, u)-\sum_{\nu=0}^{\mu-1} g_{\nu}(x) u_{1}^{\nu} \equiv 0 \text { on } \mathcal{V}(P)
$$

Proof of (15). - We write $F_{1}\left(x, u_{1}\right)=\sum_{i=0}^{\mu} b_{i}(x) u_{1}^{\mu-i}$ with $b_{0} \equiv 1$. Let $x \in \mathbb{C}^{n-r} \backslash S_{1}, \pi^{-1}(x)=\left(z^{(1)}, \ldots, z^{(\mu)}\right)$. Since $\delta(x) \neq 0$, the values $u_{1}\left(z^{(\nu)}\right)$, $1 \leq \nu \leq \mu$ are distinct. Consider the sum

$$
\begin{aligned}
\sum_{\nu=1}^{\mu} \frac{F_{1}\left(x, u_{1}\right)}{u_{1}-u_{1}\left(z^{(\nu)}\right)} f\left(z^{(\nu)}\right) & =\sum_{\nu=1}^{\mu} \frac{F_{1}\left(x, u_{1}\right)-F_{1}\left(x, u_{1}\left(z^{(\nu)}\right)\right)}{u_{1}-u_{1}\left(z^{(\nu)}\right)} f\left(z^{(\nu)}\right) \\
& =\sum_{\nu=1}^{\mu} f\left(z^{(\nu)}\right) \sum_{i=0}^{\mu} b_{i}(x) \frac{u_{1}^{\mu-i}-\left(u_{1}\left(z^{(\nu)}\right)\right)^{\mu-i}}{u_{1}-u_{1}\left(z^{(\nu)}\right)} \\
& =\sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k}
\end{aligned}
$$

where $g_{k}^{\prime}(x)=\sum_{k<i \leq \mu} b_{\mu-i}(x) \sum_{1 \leq \nu \leq \mu}\left(u_{1}\left(z^{(\nu)}\right)\right)^{i-k-1} f\left(z^{(\nu)}\right)$. Since $b_{\mu-i}(x)$ has degree $\leq \mu-i$ and $\left|z^{(\nu)}\right| \leq$ const. $(1+|x|)$ by (14) we obtain, if we note that $u_{1}$ is linear and $f$ has degree $\leq d$,

$$
\begin{aligned}
\left|g_{k}^{\prime}(x)\right| & \leq \text { const. } \max _{0<i \leq \mu}\left[(1+|x|)^{\mu-i}(1+|x|)^{i-k-1}(1+|x|)^{d}\right] \\
& \leq \text { const. }(1+|x|)^{d+\mu-k-1}, x \in \mathbb{C}^{n-r} \backslash S_{1} .
\end{aligned}
$$

This implies that $g_{k}^{\prime}$ extends holomorphically to $\mathbb{C}^{n-r}$ and is a polynomial of degree $\leq d+\mu-k-1$.

If we substitute $u_{1}=u_{1}\left(z^{(j)}\right), 1 \leq j \leq \mu,(x, u)=z^{(j)}$, in the above equation, we obtain, for $x \notin S_{1}$,

$$
\frac{\partial F_{1}}{\partial u_{1}}\left(x, u_{1}\left(z^{(j)}\right)\right) f\left(x, u\left(z^{(j)}\right)\right)=\sum_{k=0}^{\mu-1} g_{k}^{\prime}(x)\left(u_{1}\left(z^{(j)}\right)\right)^{k}, 1 \leq j \leq \mu
$$

Since $1 \leq j \leq \mu$ is arbitrary, this gives:

$$
\begin{equation*}
\frac{\partial F_{1}}{\partial u_{1}}\left(x, u_{1}\right) f(x, u)-\sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k}=0 \text { on } \mathcal{V}(P) \tag{16}
\end{equation*}
$$

if $x \notin S_{1}$; since $\pi^{-1}\left(S_{1}\right)$ is nowhere dense in $\mathcal{V}(P)$, (16) holds for all $(x, u) \in \mathcal{V}(P)$. Here $g_{k}^{\prime}$ is a polynomial of degree $\leq d+\mu-k-1$.

We now use the following fact about discriminants (see van der Waerden, Algebra, vol. 1).
(17) Let $F_{1}\left(x, u_{1}\right)=u_{1}^{\mu}+b_{1}(x) u_{1}^{\mu-1}+\cdots+b_{\mu}(x)$. There exist polynomials $A\left(x, u_{1}\right)$ and $B\left(x, u_{1}\right)$ with the following properties:
(i) $\operatorname{deg}_{u_{1}} A<\mu-1, \operatorname{deg}_{u_{1}} B<\mu$.
(ii) $\delta(x)=A\left(x, u_{1}\right) F_{1}\left(x, u_{1}\right)+B\left(x, u_{1}\right) \frac{\partial F_{1}}{\partial u_{1}}\left(x, u_{1}\right)$.
(iii) The coefficients (of the powers of $u_{1}$ ) in $A\left(x, u_{1}\right), B\left(x, u_{1}\right)$ are polynomials, with integer coefficients depending only on $\mu$, in $b_{1}, \ldots, b_{\mu}$ of total degree $\leq 2(\mu-1)$.

Note. - The existence of $A, B$ satisfying (1) and (ii) is well known. For (iii), if $\delta$ is the determinant of the $(2 \mu-1) \times(2 \mu-1)$ matrix $\Delta$ with entries $b_{\nu},(\mu-\nu) b_{\nu}$ referred to earlier, the coefficients of $A, B$ are, up to sign, the determinants of the $(2 \mu-2) \times(2 \mu-2)$ matrices obtained from $\Delta$ by deleting one row and one column; see the book of van der Waerden cited above.

Since $\operatorname{deg} b_{i} \leq i \leq \mu, A$ and $B$ have total degree $\leq 2 \mu^{2}-\mu$.
From (16) and (17), we obtain

$$
\begin{aligned}
& \delta(x) f(x, u)- B\left(x, u_{1}\right) \sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k} \\
&= A\left(x, u_{1}\right) F_{1}\left(x, u_{1}\right) f(x, u) \\
& \quad+B\left(x, u_{1}\right)\left(\frac{\partial F_{1}}{\partial u_{1}}\left(x, u_{1}\right) f(x, u)-\sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k}\right) \\
& \equiv 0 \text { on } \mathcal{V}(P)
\end{aligned}
$$

We now make an algebraic division of $B\left(x, u_{1}\right) \sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k}$ by the monic polynomial $F_{1}\left(x, u_{1}\right)$ and write

$$
\begin{equation*}
B\left(x, u_{1}\right) \sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k}=\varphi\left(x, u_{1}\right) F_{1}\left(x, u_{1}\right)+\sum_{\nu=0}^{\mu-1} g_{\nu}(x) u_{1}^{\nu} \tag{19}
\end{equation*}
$$

Let $G\left(x, u_{1}\right)=a_{0}(x) u_{1}^{N}+a_{1}(x) u_{1}^{N-1}+\cdots+a_{N}(x)$ be a polynomial in $x, u_{1}$ of total degree $\leq d$. Then $\operatorname{deg} a_{0}(x) \leq d-N$. We claim that if we
write

$$
G\left(x, u_{1}\right)=q\left(x, u_{1}\right) F_{1}\left(x, u_{1}\right)+h\left(x, u_{1}\right)
$$

with $q, h \in \mathbb{C}\left[x, u_{1}\right]$ and $\operatorname{deg}_{u_{1}} h<N$, then the total degree, $\operatorname{deg} h$, of $h$ is $\leq d$. To see this, note first that $q$ and $h$ are uniquely determined by $G$, so that, if $N<\mu$, we have $h=G$. If $N \geq \mu$, and we set $G_{1}\left(x, u_{1}\right)=G\left(x, u_{1}\right)-a_{0}(x) u_{1}^{N-\mu} F_{1}\left(x, u_{1}\right)$, then $\operatorname{deg}_{u_{1}} G_{1}<N$. Further, $F_{1}$ has total degree $\leq d$ (since $\operatorname{deg} b_{i} \leq i$ ) and $\operatorname{deg} a_{0} \leq d-N$. Hence $G_{1}$ has total degree $\leq d$, and $\operatorname{deg}_{u_{1}} G_{1} \leq N-1$. Our claim now follows by induction on $N$ since the remainder on division of $G_{1}$ by $F_{1}$ is the same as for division of $G$ by $F_{1}$.

Thus, in (19),

$$
\begin{aligned}
\operatorname{deg}\left(\sum_{\nu=0}^{\mu-1} g_{\nu}(x) u_{1}^{\nu}\right) & \leq \operatorname{deg} B+\operatorname{deg} \sum_{k=0}^{\mu-1} g_{k}^{\prime}(x) u_{1}^{k} \\
& \leq C(d, \mu)
\end{aligned}
$$

Since $F_{1}\left(x, u_{1}\right) \equiv 0$ on $\mathcal{V}(P),(18)$ and (19) imply that

$$
\delta(x) f(x, u)-\sum_{\nu=0}^{\mu-1} g_{\nu}(x) u_{1}^{\nu} \equiv 0 \text { on } \mathcal{V}(P)
$$

which proves (15).
If $F_{j}\left(x, u_{j}\right)$ is the polynomial constructed in Lemma 5.8 for the linear form $\ell=u_{j}$ and we apply (15) to the polynomials $f(x, u)=u_{\nu}$ $(\nu=2, \ldots, \mu)$, we obtain:
(20) There exist monic polynomials

$$
F_{j}\left(x, u_{j}\right)=u_{j}^{\mu}+\sum_{i=1}^{\mu} b_{j i}(x) u_{j}^{\mu-i}, 1 \leq j \leq r
$$

where $\operatorname{deg} b_{j i} \leq i, \operatorname{deg} F_{j}=\mu$, and polynomials $G_{\nu}\left(x, u_{1}\right), \nu=2, \ldots, r$ with $\operatorname{deg}_{u_{1}} G_{\nu}<\mu, \operatorname{deg} G_{\nu} \leq C(\mu)$ such that

$$
F_{j}\left(x, u_{j}\right), j=1, \ldots, r, \delta(x) u_{\nu}-G_{\nu}\left(x, u_{1}\right), \nu=2, \ldots, r
$$

all vanish on $\mathcal{V}(P)$.
Lemma 5.9. - Let $N=\max (\mu,(\mu-1)(r-1))$. We have the following:
(21) If $f \in \mathbb{C}[x, u]$, then $f \equiv 0$ on $\mathcal{V}(P)$ if and only if there exist polynomials $A_{j}(x, u), j=1, \ldots, r$, such that

$$
(\delta(x))^{N} f(x, u)=A_{1}(x, u) F_{1}\left(x, u_{1}\right)+\sum_{\nu=2}^{r} A_{\nu}(x, u)\left(\delta(x) u_{\nu}-G_{\nu}\left(x, u_{1}\right)\right)
$$

Proof of Lemma 5.9. - If $f$ satisfies this equation, then $\delta^{N} f \equiv 0$ on $\mathcal{V}(P)$, hence also $f \equiv 0$ on $\mathcal{V}(P)$ since $\{(x, u) \in \mathcal{V}(P) \mid \delta(x) \neq 0\}$ is dense in $\mathcal{V}(P)$.

To prove the converse, let $f(x, u) \in \mathbb{C}[x, u]$. We make an algebraic division of $f$ by the monic polynomial $F_{1}\left(x, u_{1}\right)$ to write

$$
f(x, u)=A_{1}^{\prime}(x, u) F_{1}\left(x, u_{1}\right)+\sum_{\nu=0}^{\mu-1} f_{\nu}\left(x, u_{2}, \ldots, u_{r}\right) u_{1}^{\nu}, A_{1}^{\prime} \in \mathbb{C}[x, u]
$$

Dividing the coefficients $f_{\nu}$ by $F_{2}\left(x, u_{2}\right)$ and repeating this process with $F_{3}\left(x, u_{3}\right), \ldots, F_{r}\left(x, u_{r}\right)$, we can write

$$
f(x, u)=\sum_{j=1}^{r} A_{j}^{\prime}(x, u) F_{j}\left(x, u_{j}\right)+\sum_{0 \leq \nu_{1}<\mu} f_{\nu_{1} \cdots \nu_{r}}(x) u_{1}^{\nu_{1}} \cdots u_{r}^{\nu_{r}}
$$

where $A_{j}^{\prime} \in \mathbb{C}[x, u]$ and $f_{\nu_{1} \cdots \nu_{r}} \in \mathbb{C}[x]$.
We remark that for $j=2, \ldots, r$, if we substitute $G_{j}\left(x, u_{1}\right)$ for $\delta(x) u_{j}$, we find that

$$
(\delta(x))^{\mu} F_{j}\left(x, u_{j}\right) \equiv R_{j}^{\prime}\left(x, u_{1}\right) \bmod \left(\delta(x) u_{j}-G_{j}\left(x, u_{1}\right)\right)
$$

where $R_{j}^{\prime} \in \mathbb{C}\left[x, u_{1}\right]$. Dividing $R_{j}^{\prime}\left(x, u_{1}\right)$ by $F_{1}\left(x, u_{1}\right)$ in $\mathbb{C}\left[x, u_{1}\right]$, we see that

$$
(\delta(x))^{\mu} F_{j}\left(x, u_{j}\right) \equiv R_{j}\left(x, u_{1}\right) \bmod \left(F_{1}\left(x, u_{1}\right), \delta(x) u_{j}-G_{j}\left(x, u_{1}\right)\right)
$$

where $\operatorname{deg}_{u_{1}} R_{j}\left(x, u_{1}\right)<\mu$. Since $F_{j}\left(x, u_{j}\right), F_{1}\left(x, u_{1}\right), \delta(x) u_{j}-G_{j}\left(x, u_{1}\right)$ all vanish on $\mathcal{V}(P)$, so also does $R_{j}\left(x, u_{1}\right)$. But, since $\operatorname{deg}_{u_{1}} R_{j}\left(x, u_{1}\right)<\mu$, this implies that $R_{j} \equiv 0$ as remarked earlier. Hence:
(23) $(\delta(x))^{\mu} F_{j}\left(x, u_{j}\right) \equiv 0 \bmod \left(F_{1}\left(x, u_{1}\right), \delta(x) u_{j}-G_{j}\left(x, u_{1}\right)\right), j=2, \ldots, r$.

Now, for $k=2, \ldots, r, 0 \leq \nu<\mu$,

$$
(\delta(x))^{\mu-1} u_{k}^{\nu} \equiv \delta(x)^{\mu-\nu-1}\left(G_{k}\left(x, u_{1}\right)\right)^{\nu} \bmod \left(\delta(x) u_{k}-G_{k}\left(x, u_{1}\right)\right)
$$

Hence

$$
\begin{aligned}
\delta(x)^{(\mu-1)(r-1)} & \sum_{0 \leq \nu_{i}<\mu} f_{\nu_{1} \cdots \nu_{r}}(x) u_{1}^{\nu_{1}} \cdots u_{r}^{\nu_{2}} \equiv B^{\prime}\left(x, u_{1}\right) \\
& \bmod \left(\delta(x) u_{2}-G_{2}\left(x, u_{1}\right), \ldots, \delta(x) u_{r}-\dot{G}_{r}\left(x, u_{1}\right)\right)
\end{aligned}
$$

where $B^{\prime}\left(x, u_{1}\right)$ is a polynomial in $x, u_{1}$ alone. Dividing $B^{\prime}\left(x, u_{1}\right)$ by $F_{1}\left(x, u_{1}\right)$ in $\mathbb{C}\left[x, u_{1}\right]$, we see that $B^{\prime}\left(x, u_{1}\right) \equiv B\left(x, u_{1}\right) \bmod F_{1}\left(x, u_{1}\right)$, where $\operatorname{deg}_{u_{1}} B\left(x, u_{1}\right)<\mu$. Multiplying (22) by $\delta(x)^{N}, N=\max (\mu,(\mu-1)(r-1))$ and using the above facts together with (23), we have:

Given $f \in \mathbb{C}[x, u]$, there exists a polynomial $B\left(x, u_{1}\right)$ in $x$ and $u_{1}$ alone, with $\operatorname{deg}_{u_{1}} B<\mu$, such that

$$
(\delta(x))^{N} f(x, u) \equiv B\left(x, u_{1}\right) \bmod \left(F_{1}\left(x, u_{1}\right), \delta(x) u_{\nu}-G_{\nu}\left(x, u_{1}\right), 2 \leq \nu \leq r\right)
$$

Now, if $f \equiv 0$ on $\mathcal{V}(P)$, then $B\left(x, u_{1}\right) \equiv 0$ on $\mathcal{V}(P)$. But again, as remarked earlier, a polynomial in $x, u_{1}$ alone, of degree $<\mu$ in $u_{1}$ and vanishing on $\mathcal{V}(P)$ must be the zero polynomial. Thus, $f \equiv 0$ on $\mathcal{V}(P)$ if and only if $B\left(x, u_{1}\right) \equiv 0$, i.e. if and only if $(\delta(x))^{N} f(x, u) \equiv 0$ $\bmod \left(F_{1}\left(x, u_{1}\right), \delta(x) u_{\nu}-G_{\nu}\left(x, u_{1}\right), 2 \leq \nu \leq r\right)$, which is Lemma 5.9.

To complete the proof of Theorem 5.6, we invoke the following theorem, which is also proved in the paper [He] of Grete Hermann cited earlier. It can be proved in exactly the same way as Theorem 5.4.

Theorem. - Let $n, D, r, s$ be integers $\geq 1$. There are integers $D^{\prime} \geq 1$ and $r^{\prime} \geq 1$ depending only on $n, D, r$ such that the following holds.

Let $\left(p_{i j}\right)_{1 \leq i \leq r, 1 \leq j \leq s}$ be a matrix of polynomials $p_{i j} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ with $\operatorname{deg} p_{i j} \leq D$. Consider the system of homogeneous equations

$$
\begin{equation*}
\sum_{i=1}^{r} g_{i} p_{i j}=0, j=1, \ldots, s, g_{i} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right] \tag{24}
\end{equation*}
$$

There exist $r^{\prime}$ solutions $g^{(\alpha)}=\left(g_{\alpha 1}, \ldots, g_{\alpha r}\right), \alpha=1, \ldots, r^{\prime}$ of (24) where $g_{\alpha i} \in \mathbb{C}\left[z_{1}, \ldots, z_{n}\right], \operatorname{deg} g_{\alpha i} \leq D^{\prime}$, such that
(i) $\sum_{i} g_{\alpha i} p_{i j}=0, j=1, \ldots, s, \alpha=1, \ldots, r^{\prime}$.
(ii) If $g=\left(g_{1}, \ldots, g_{r}\right)$ satisfies (24), then we can find $h_{\alpha} \in$ $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right], \alpha=1, \ldots, r^{\prime}$, with $g=\sum_{\alpha} h_{\alpha} g^{(\alpha)}$.

Applying this to the equation in Lemma 5.9, (21), we obtain the following:

There exist constants $\bar{D}, \bar{r} \geq 1$ depending only on $\mu, n, r, D, N$ hence only on $n, r, D$ (since $\operatorname{deg} \delta \leq \mu(2 \mu-1), \mu \leq D^{r}$ and $N \leq \mu r$ ) and a set of $\bar{r}(r+1)$-tuples

$$
\left(f_{\alpha}, A_{1}^{(\alpha)}, \ldots, A_{r}^{(\alpha)}\right), \alpha=1, \ldots, \bar{r}
$$

of polynomials of degree $\leq \bar{D}$ in $(x, u)$ such that
(a) $(\delta(x))^{N} f_{\alpha}(x, u)=A_{1}^{(\alpha)}(x, u) F_{1}\left(x, u_{1}\right)+\sum_{\nu=2}^{r} A_{\nu}^{(\alpha)}(x, u)\left(\delta(x) u_{\nu}-\right.$ $\left.G_{\nu}\left(x, u_{1}\right)\right)$
and
(b) if $\left(f, A_{j}\right)$ satisfies the equation in (21), then, there exist $h_{\alpha} \in$ $\mathbb{C}[x, u]$ with

$$
\left(f, A_{1}, \ldots, A_{r}\right)=\sum_{\alpha} h_{\alpha}\left(f_{\alpha}, A_{1}^{(\alpha)}, \ldots, A_{r}^{(\alpha)}\right)
$$

It follows from Lemma 5.9 that the ideal of $f \in \mathbb{C}[x, u]$ which vanish on $\mathcal{V}(P)$ is generated by the $f_{\alpha}, \alpha=1, \ldots, \bar{r}$. Further, if $f \in \mathbb{C}[x, u]$ vanishes in some neighborhood of 0 in $\mathcal{V}(P)$, then $f \mid \mathcal{V}(P) \equiv 0$ since the set of regular points of $\mathcal{V}(P)$ is connected and dense.

Theorem 5.6 follows from this and Theorem 5.4.
Note. - Madhav Nori has shown us a proof (using generic flatness) that if $I$ is an ideal in $\mathbb{C}\left[z_{1}, \ldots, z_{n}\right]$ generated by $r$ polynomials of degree $\leq D$, then the radical of $I$ is generated by $r^{\prime}$ polynomials of degree $\leq D^{\prime}\left(r^{\prime}, D^{\prime}\right.$ depending only on $\left.n, D, r\right)$. Nori ascribes the method to Grothendieck. We could use this theorem instead of Theorem 5.6 (because of Theorem 3.1).

The proof given above seems to us more elementary, and picks out the radical of one primary component in the decomposition of $I$.

## 6. Semi-algebraic sets and maps: general properties.

In this section, we collect together the general properties of semialgebraic sets and maps that we shall need. Although several of these results are in the literature (see e.g. [BCR]), we have given direct proofs based on the two basic structure theorems of the subject.

## Definition 6.1.

(a) The family of semi-algebraic sets in $\mathbb{R}^{n}$ is the smallest class of sets containing all sets of the form $\left\{x \in \mathbb{R}^{n} \mid P(x)>0\right\}$ where $P \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$, and which is closed under the operations of finite unions, finite intersections and complementation.

The class of semi-algebraic sets is invariant under linear isomorphisms of $\mathbb{R}^{n}$, so that we may speak of semi-algebraic sets in any finite dimensional vector space over $\mathbb{R}$.
(b) If $E \subset \mathbb{R}^{n}$, a map (or function) $f: E \rightarrow \mathbb{R}^{M}$ is called semialgebraic if its graph $\{(x, f(x)) \mid x \in E\}$ is a semi-algebraic subset of $\mathbb{R}^{m+n}$.
(c) A function $f: E \rightarrow \mathbb{R} \cup\{-\infty,+\infty\}$ is called an extended semi-algebraic function if the sets $E_{-}=f^{-1}(-\infty), E_{+}=f^{-1}(+\infty)$, $E_{0}=f^{-1}(\mathbb{R})$ are all semi-algebraic and $f \mid E_{0} \rightarrow \mathbb{R}$ is a semi-algebraic function.
(d) If $E \subset \mathbb{R}^{n}$ is semi-algebraic, a semi-algebraic partition of $E$ is a finite collection $\left\{E_{\nu}\right\}$ of semi-algebraic sets which are pointwise disjoint and whose union is $E$.

A partition $\left\{F_{\mu}\right\}$ of $E$ is a refinement of the partition $\left\{E_{\nu}\right\}$ of $E$ if, for each $\nu, E_{\nu}$ is the union of those $F_{\mu}$ which meet it: $E_{\nu}=\bigcup_{F_{\mu} \cap E_{\mu} \neq \emptyset} F_{\mu}$.

We note, explicitly, that we are not assuming that a semi-algebraic map is continuous. If $E$ is semi-algebraic and $\left\{E_{\nu}\right\}$ is a semi-algebraic partition of $E$, then a function $f: E \rightarrow \mathbb{R}^{m}$ is semi-algebraic if and only if $f \mid E_{\nu}$ is semi-algebraic for all $\nu$.

The basic structure theorem, which enables one to reduce the study of semi-algebraic sets in $\mathbb{R}^{n+1}$ to that of sets in $\mathbb{R}^{n}$ is the following. For a proof see Cohen [C]; see also the book of Bochnak, Coste, Roy [BCR].

The Structure Theorem. - Let $E$ be a semi-algebraic set in $\mathbb{R}^{n+1}=\mathbb{R}^{n} \times \mathbb{R}$. There is a (finite) semi-algebraic partition $\left\{A_{\nu}\right\}$ of $\mathbb{R}^{n}$ and, for each $\nu$, a finite family of functions $t_{\nu j}$ on $A_{\nu}, 0 \leq j \leq r_{\nu}+1$, with values in $\mathbb{R} \cup\{-\infty, \infty\}$ such that

$$
-\infty \equiv t_{\nu, 0}<t_{\nu, 1}<\cdots<t_{\nu, r_{\nu}}<t_{\nu, r_{\nu}+1} \equiv+\infty
$$

$t_{\nu, j}$ is continuous for $1 \leq j \leq r_{\nu}$, and having the following properties:
(a) Each set $\left\{(x, t) \in \mathbb{R}^{n+1} \mid x \in A_{\nu}, t_{\nu, j}(x)<t<t_{\nu, j+1}(x)\right\}$, $0 \leq j \leq r_{\nu}$, is semi-algebraic.
(b) Each set $\left\{(x, t) \in \mathbb{R}^{n+1} \mid x \in A_{\nu}, t=t_{\nu, j}(x)\right\}, 0<j \leq r_{\nu}$, is semi-algebraic.
(c) $E$ is a finite disjoint union of sets of the form (a) or (b).

By considering the graph of a semi-algebraic map $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ and using the structure theorem, one obtains the central theorem about semi-algebraic sets.

Tarski-Seidenberg Theorem. - Let $f: \mathbb{R}^{n} \rightarrow \mathbb{R}^{m}$ be a semialgebraic map and let $E \subset \mathbb{R}^{n}$ be a semi-algebraic set. Then $f(E) \subset \mathbb{R}^{m}$ is again semi-algebraic.

We list some basic properties in the next two results. They are easy consequences of the Tarski-Seidenberg Theorem (see [BCR]).

## Proposition 6.2.

(a) Let $E \subset \mathbb{R}^{n}$ be semi-algebraic, and let $f: E \rightarrow \mathbb{R}^{m}$ be a semialgebraic map. Then, for any semi-algebraic set $F \subset \mathbb{R}^{m}$, the set $f^{-1}(F)$ is again semi-algebraic.
(b) Let $E \subset \mathbb{R}^{n}, F \subset \mathbb{R}^{m}$ be semi-algebraic, let $f: E \rightarrow \mathbb{R}^{m}$ be a semi-algebraic map with $f(E) \subset F$ and let $g: F \rightarrow \mathbb{R}^{k}$ be semi-algebraic. Then $g \circ f: E \rightarrow \mathbb{R}^{k}$ is semi-algebraic.
(c) Let $E \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ be semi-algebraic. Then, the set

$$
\left\{x \in \mathbb{R}^{n} \mid \exists y \in \mathbb{R}^{m} \text { with }(x, y) \in E\right\}
$$

is semi-algebraic.
(d) Let $E \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ be semi-algebraic, let $F \subset \mathbb{R}^{m}$ be semi-algebraic. Then, the set

$$
\left\{x \in \mathbb{R}^{n} \mid(x, y) \in E, \forall u \in F\right\}
$$

is semi-algebraic.
Lemma 6.3. - The closure $\bar{E}$ of a semi-algebraic set $E \subset \mathbb{R}^{n}$ is again semi-algebraic.

We shall also need a second fundamental theorem on semi-algebraic sets which goes back to Whitney's work on real algebraic varieties [W]. The theorem is proved in [BCR].

Before giving the definition needed to state this theorem, we remark that by a smooth (real analytic) submanifold $M$ of $\mathbb{R}^{n}$, we mean a subset $M \subset \mathbb{R}^{n}$ with the following property: for any $a \in M$, we can find an open neighborhood $U$ of $a$ in $\mathbb{R}^{n}$, an integer $d \geq 0$, and smooth (real analytic) coordinates $x_{1}, \ldots, x_{n}$ on $U$ such that $M \cap U=\left\{x \in U \mid x_{d+1}=\cdots=x_{n}=0\right\}$.

Definition 6.4. - Let $E \subset \mathbb{R}^{n}$. A stratification $\left\{S_{i}\right\}$ of $E$ is a finite family of pairwise disjoint subsets $S_{i} \subset E$ with $E=\cup S_{i}$ and having the following properties:
(a) Each $S_{i}$ is a real analytic submanifold of $\mathbb{R}^{n}$.
(b) If $i \neq j$ and $\bar{S}_{i} \cap S_{j} \neq 0$, then $S_{j} \subset \bar{S}_{i}$ and $\operatorname{dim} S_{j}<\operatorname{dim} S_{i}$.

The $\left\{S_{i}\right\}$ are called the strata of the stratification.
If $E$ is semi-algebraic, the stratification $\left\{S_{i}\right\}$ is called semi-algebraic if each stratum $S_{i}$ is semi-algebraic. In particular, $\left\{S_{i}\right\}$ is a semi-algebraic partition of $E$.

The Stratification Theorem. - Let $E \subset \mathbb{R}^{n}$ be semi-algebraic and let $\left\{E_{\nu}\right\}$ be a semi-algebraic partition of $E$.

There exists a semi-algebraic stratification $\left\{S_{i}\right\}$ of $E$ which is also a refinement of $E_{\nu}$, i.e. each $E_{\nu}$ is the union of those $S_{i}$ which meet it.

Definition 6.5. - Let $E \subset \mathbb{R}^{n}$ be semi-algebraic. Let $\left\{S_{i}\right\}$ be a semi-algebraic stratification of $E$. We define the dimension of $E, \operatorname{dim} E$, by:

$$
\operatorname{dim} E=-1 \text { if } E=\emptyset, \operatorname{dim} E=\max _{i} \operatorname{dim} S_{i} \text { if } E \neq \emptyset
$$

Remark. - Given a semi-algebraic set $E \subset \mathbb{R}^{n}$, let $t$ be the largest integer such that $E$ contains the image of the unit ball $B_{t}$ in $\mathbb{R}^{t}$ under a $C^{\infty} \operatorname{map} \varphi: B_{t} \rightarrow \mathbb{R}^{n}$ of rank $t$ everywhere on $B_{t}$. Then $t=\operatorname{dim} E$; in fact if $d=\max \operatorname{dim} S_{i}\left(\left\{S_{i}\right\}\right.$ a semi-algebraic stratification of $\left.E\right)$, then clearly a submanifold of dimension $d$ of $\mathbb{R}^{n}$ contains the image of a ball of dimension $d$ under a smooth map of maximal rank, so that $t \geq d$. If $t>d$ and $\varphi: B_{t} \rightarrow \mathbb{R}^{n}$ has rank $t$, then $\varphi^{-1}\left(S_{i}\right)$ has measure 0 in $B_{t}$ for each $i$ and we cannot have $\varphi\left(B_{t}\right) \subset \cup S_{i}$.

We shall say that a set $E \subset \mathbb{R}^{n}$ contains the diffeomorphic image of a ball of dimension $t$ if there exists a $C^{\infty}$ embedding $\varphi: B_{t} \rightarrow \mathbb{R}^{n}$ (i.e. $\varphi$ is injective of maximal rank) so that $\varphi\left(B_{t}\right) \subset E$. We may formulate our remark above as follows:

If $E \subset \mathbb{R}^{n}$ is semi-algebraic, then $\operatorname{dim} E$ is the largest integer $t$ such that $E$ contains the diffeomorphic image of a ball of dimension $t$.

In particular, $\operatorname{dim} E$ is independent of the stratification $\left\{S_{i}\right\}$ used in Definition 6.5.

Lemma 6.6. - Let $E_{1}, \ldots, E_{N}$ be semi-algebraic sets in $\mathbb{R}^{n}$. Then

$$
\operatorname{dim}\left(E_{1} \cup \cdots \cup E_{N}\right)=\max _{1 \leq \nu \leq N} \operatorname{dim} E_{\nu}
$$

Proof. - If $E=\underset{1 \leq \nu \leq N}{ } E_{\nu}$, then clearly $\operatorname{dim} E \geq \max _{\nu} \operatorname{dim} E_{\nu}$ (since if $E_{\nu}$ contains the diffeomorphic image of a ball of dimension $t$, so does $E$ ).

Let now $t>\max _{\nu} \operatorname{dim} E_{\nu}$ and $\varphi: B_{t} \rightarrow \mathbb{R}^{n}$ be a $C^{\infty}$ embedding. If $\left\{S_{i, \nu}\right\}$ is a stratification of $E_{\nu}$, then $\operatorname{dim} S_{i, \nu}<t$ (for all $i, \nu$ ) so that $\varphi^{-1}\left(S_{i, \nu}\right)$ has measure 0 in $B_{t}$. Hence $\varphi^{-1}(E)$ has measure 0 in $B_{t}$, so that $\varphi\left(B_{t}\right) \not \subset E$. Hence $t>\operatorname{dim} E$, which shows that $\operatorname{dim} E \leq \max _{\nu} \operatorname{dim} E_{\nu}$.

Lemma 6.7. - Let $E \subset \mathbb{R}^{n}$ be semi-algebraic and let $\bar{E}$ be the closure. Then,

$$
\operatorname{dim} \bar{E}=\operatorname{dim} E \text { and } \operatorname{dim}(\bar{E} \backslash E)<\operatorname{dim} E
$$

Proof. - Let $\left\{S_{i}\right\}$ be a semi-algebraic stratification of $\bar{E}$ refining the partition $\{E, \bar{E} \backslash E\}$ of $\bar{E}$. Let $J=\left\{i \mid S_{i} \cap E \neq \emptyset\right\}$. Then $E=\bigcup_{i \in J} S_{i}$, $\bar{E} \backslash E=\bigcup_{j \notin J} S_{j}$. We have $\bar{E}=\bigcup_{i \in J} \bar{S}_{i}$, so that, if $j \notin J$, then $S_{j} \cap \bar{S}_{i} \notin \emptyset$ for some $i \in J$. By condition (b) in Definition 6.4, this implies that $S_{j} \subset \bar{S}_{i}$ and $\operatorname{dim} S_{j}<\operatorname{dim} S_{i} \leq \operatorname{dim} E$. Both statements in the lemma follow from this.

We now make a simple remark which we shall use in Lemma 6.9 below and also later.

Remark 6.8. - Let $X, Y$ be connected $C^{\infty}$ manifolds and $f: X \rightarrow Y$, a smooth map. Assume that for any $y \in Y$, the fibre $f^{-1}(y)$ is discrete (consists of isolated points). Then, the set $\left\{x \in X \mid \operatorname{rank}\left(d f_{x}\right)=\operatorname{dim} X\right\}, d f_{x}$ being the differential of $f$ at $x$, is (open and) dense in $X$.

In fact, let $U \subset X$ be open, $U \neq \emptyset$. Let $p=\max _{x \in U} \operatorname{rank}\left(d f_{x}\right)$, and let $x_{0} \in U$ be such that rank $\left(d f_{x_{0}}\right)=p$. Then, there is a connected open set $U_{0} \subset U, x_{0} \in U_{0}$, so that rank $\left(d f_{x}\right)=p$ for all $x \in U_{0}$. By the rank theorem, the sets $f^{-1} f(x) \cap U_{0}, x \in U_{0}$, are submanifolds of $U_{0}$ of dimension $\operatorname{dim} X-p$; since the fibres are discrete, we must have $p=\operatorname{dim} X$.

Lemma 6.9. - Let $E \subset \mathbb{R}^{n}$ be a semi-algebraic set and let $f: E \rightarrow \mathbb{R}^{m}$ be a semi-algebraic map. Let $\Gamma=\{(x, f(x)) \mid x \in E\}$ be the graph of $f$. Then:
(a) $\operatorname{dim} E=\operatorname{dim} \Gamma$
(b) $\operatorname{dim} f(E) \leq \operatorname{dim} E$.
(c) If, moreover, the fibres of $f$ are finite, then $\operatorname{dim} f(E)=\operatorname{dim} E$.

Proof.
(a) Let $S \subset \Gamma$ be a submanifold of $\mathbb{R}^{n} \times \mathbb{R}^{m}$ of dimension $d$. The projection $\pi: \Gamma \rightarrow E, \pi(x, f(x))=x$ is bijective, so that by 4.8 , there is an open set $U \subset S$ such that $\pi \mid U$ has rank $=\operatorname{dim} S$ at every point of $U$. Hence $E$ contains the image of a ball of dimension $d$, under a map of maximal rank, so that $d \leq \operatorname{dim} E$; thus $\operatorname{dim} \Gamma \leq \operatorname{dim} E$.

On the other hand, if $\varphi: B_{t} \rightarrow \mathbb{R}^{n}$ is a $C^{\infty}$ embedding with $t>\operatorname{dim} \Gamma$, and $\left\{S_{1}, \ldots, S_{N}\right\}$ is a semi-algebraic stratification of $\Gamma$, then $\operatorname{dim} S_{i}<t$ and $\pi\left(S_{i}\right) \subset \mathbb{R}^{n}$ has $t$-dimensional measure 0 . Hence $\varphi\left(B_{t}\right) \not \subset \pi\left(S_{1} \cup \cdots \cup S_{N}\right)=$ $\pi(\Gamma)=E$. Hence $\operatorname{dim} E<t$, i.e. $\operatorname{dim} E \leq \operatorname{dim} \Gamma$.
(b) By (a), it is enough to show that if $\alpha: \Gamma \rightarrow \mathbb{R}^{m}$ is the projection $(x, y) \mapsto y$ restricted to $\Gamma$, then $\operatorname{dim} \alpha(\Gamma) \leq \operatorname{dim} \Gamma$.

If $t=\operatorname{dim} \Gamma=\operatorname{dim} E$ and $\Gamma=\cup S_{i}$ is a semi-algebraic stratification of $\Gamma$, the $(t+1)$-dimensional measure of $\alpha\left(S_{i}\right)$ is 0 for any $i$, so that $\alpha(\Gamma)$ cannot contain the diffeomorphic image of a ball of dimension $\geq t+1$.
(c) If $f: E \rightarrow \mathbb{R}^{m}$ has finite fibres and $\alpha$ is the restriction to $\Gamma$ of the projection $(x, y) \mapsto y$, then $\alpha$ has finite fibres. If $S \subset \Gamma$ is a smooth submanifold with $\operatorname{dim} S=\operatorname{dim} \Gamma$, Remark 6.8 implies that there is an open set $U \subset S, U \neq \emptyset$, so that $\alpha \mid U$ has rank $\operatorname{dim} S$ at every point of $U$. Since $\alpha(U) \subset \alpha(\Gamma)$, we must have $\operatorname{dim} \alpha(\Gamma) \geq \operatorname{dim} S=\operatorname{dim} \Gamma$. This, together with (a) and (b), show that $f(E)=\alpha(\Gamma)$ has the same dimension as $E$.

Remark 6.10. - We shall use the following simple remark. If $A \subset$ $\mathbb{R}^{n}$, denote by $\delta(A)$ the closure in $\mathbb{R}^{n}$ of the set $\bar{A} \backslash A$ ( $\bar{A}$ being the closure of $A$ ). Then $A \backslash \delta(A)$ is closed in $\mathbb{R}^{n} \backslash \delta(A)$.

In fact, if $\left\{x_{\nu}\right\}_{\nu \geq 1}$ is a sequence of points in $A \backslash \delta(A)$ converging to $x_{0} \in \mathbb{R}^{n} \backslash \delta(A)$, then $x_{0} \in \bar{A}$. If $x_{0} \notin A$, we would have $x_{0} \in \bar{A} \bigvee A \subset \delta(A)$. Thus $x_{0} \in A \backslash \delta(A)$, and $A \backslash \delta(A)$ is closed in $\mathbb{R}^{n} \backslash \delta(A)$.

Lemma 6.11. - Let $E \subset \mathbb{R}^{n}$ be a semi-algebraic set of dimension $d$ and let $f: E \rightarrow \mathbb{R}^{m}$ be a semi-algebraic map.

There exists a closed semi-algebraic set $E^{\prime} \subset \mathbb{R}^{n}$ of dimension $<d$ such that $E \backslash E^{\prime}$ is a closed real analytic submanifold of $\mathbb{R}^{n} \backslash E^{\prime}$ and $f \mid E \backslash E^{\prime}$ is real analytic.

Proof. - Let $\left\{S_{i}\right\}, 1 \leq i \leq p$, be a semi-algebraic stratification of $E$, and suppose that $\operatorname{dim} S_{i}=d$ for $1 \leq i \leq q(\leq p), \operatorname{dim} S_{i}<d$ for $q+1 \leq i \leq p$. Let $E_{0}$ be the closure in $\mathbb{R}^{n}$ of $\bigcup_{q<i \leq p} S_{i} \cup \bigcup_{1 \leq j \leq p} \delta\left(S_{j}\right)$ $\left(\delta\left(S_{j}\right)=\right.$ closure of $\bar{S}_{j} \backslash S_{j}$; see remark 6.10). Then $E_{0}$ is a closed semialgebraic set of dimension $<d$ (by Lemmas $6.3,6.6,6.7$ ), and $E \backslash E_{0}$ is a closed real analytic submanifold of $\mathbb{R}^{n} \backslash E_{0}$ of pure dimension $d$ since $\operatorname{dim} E=d=\max \left(\operatorname{dim} E_{0}, \operatorname{dim}\left(E \backslash E_{0}\right)\right)$ by Lemma 6.6.

Let $\Gamma=\left\{(x, f(x)) \in \mathbb{R}^{n} \times \mathbb{R}^{m} \mid x \in E\right\}$ be the graph of $f$; we have $\operatorname{dim} \Gamma=d$ by Lemma 6.9. Moreover, by Lemma 6.9, (c), if $F \subset \mathbb{R}^{n}$ is a semi-algebraic set, then $\operatorname{dim}\left(\left(F \times \mathbb{R}^{m}\right) \cap \Gamma\right)=\operatorname{dim}(E \cap F) \leq \operatorname{dim} F$.

Let $\left\{T_{j}\right\}, 1 \leq j \leq k$, be a semi-algebraic stratification of $\Gamma$; let $\operatorname{dim} T_{j}=d$ for $1 \leq j \leq \ell, \operatorname{dim} T_{j}<d$ for $\ell+1 \leq j \leq k$. If $\pi: \mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$ is the projection, let $E_{1}$ be the closure in $\mathbb{R}^{n}$ of $\pi\left(\bigcup_{\ell+1 \leq j \leq k} T_{j}\right) ; E_{1}$ is a closed semi-algebraic set of dimension $<d$. Let $E_{2}$ be the closure in $\mathbb{R}^{n}$ of $\pi\left(\bigcup_{1 \leq j \leq k} \delta\left(T_{j}\right)\right)\left[\delta\left(T_{j}\right)=\right.$ closure of $\left.\bar{T}_{j} \backslash T_{j}\right]$ and let $E_{3}=E_{0} \cup E_{1} \cup E_{2} ; E_{3}$ is a closed semi-algebraic set of dimension $<d$ (by Lemmas 6.6, 6.7, 6.9).

We have $\tilde{\Gamma}=\Gamma \cap\left(\mathbb{R}^{n} \times \mathbb{R}^{m} \backslash \pi^{-1}\left(E_{3}\right)\right)=\bigcup_{1 \leq j \leq l}\left(T_{j} \backslash \pi^{-1}\left(E_{3}\right)\right), \tilde{\Gamma}$ is a closed real analytic submanifold of $\mathbb{R}^{n} \times \mathbb{R}^{m} \backslash \pi^{-1}\left(E_{3}\right)$ of pure dimension $d, \pi(\tilde{\Gamma})=E \backslash E_{3}$ and $\tilde{\Gamma}$ is the graph of $f \mid E \backslash E_{3}$. Further, $E \backslash E_{3}$ is a closed real analytic submanifold of $\mathbb{R}^{n} \backslash E_{3}$ of pure dimension $d$.

We now make the following remark. Let $M \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ be a semialgebraic set which is a smooth submanifold of $\mathbb{R}^{n} \times \mathbb{R}^{m}$. Let $\varphi: M \rightarrow \mathbb{R}^{n}$ be the restriction to $M$ of the projection $\pi: \mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$. Then, the set $C=\left\{x \in M \mid\right.$ the differential $d \varphi_{x}$ of $\varphi$ at $x$ is not injective $\}$ is semialgebraic.

To see this, let

$$
\begin{aligned}
\mathcal{S}=\left\{(x, y, v) \in \mathbb{R}^{n+m} \times \mathbb{R}^{n+m} \times \mathbb{R}^{n+m}\right. & \mid x, y \in M, x \neq y, \quad \text { and } \\
\exists \lambda & \in \mathbb{R} \text { so that } v=\lambda x+(1-\lambda) y\}
\end{aligned}
$$

(we have written $\mathbb{R}^{n+m}$ for $\mathbb{R}^{n} \times \mathbb{R}^{m}$ ). $\mathcal{S}$ is the family of all secants of $M$ at distinct points and is semi-algebraic. Hence so it its closure $\overline{\mathcal{S}}$. Let

$$
\mathcal{T}=\left\{(x, v) \in \mathbb{R}^{n+m} \times \mathbb{R}^{n+m} \mid(x, x, v) \in \overline{\mathcal{S}}, x \in M\right\}
$$

Then $\mathcal{T}$ is semi-algebraic, and consists of pairs $(x, v)$ where $v$ lies on a tangent line to $M$ at $x \in M$. [Since $M$ is smooth, tangent lines are precisely limits of secant lines at distinct points.] The kernel of $d \varphi_{x}, x \in M$, consists
of those $v \in \mathbb{R}^{n+m}$ such that $(x, x+v) \in \mathcal{T}$ and $v$ projects to 0 in $\mathbb{R}^{n}$, (i.e. has first $n$ coordinates 0 ). Thus

$$
C=\left\{x \in M \mid \exists w \neq x, w \in \mathbb{R}^{n+m} \text { with } \pi(w-x)=0,(x, w) \in \mathcal{T}\right\}
$$

This is clearly semi-algebraic because of the Tarski-Seidenberg theorem.
We return to the proof of Lemma 6.11 and the notation introduced earlier. Let $\varphi=\pi \mid \tilde{\Gamma}: \tilde{\Gamma} \rightarrow E \backslash E_{3} ; \varphi$ is a bijection. Let $C=\{(x, y) \in$ $\tilde{\Gamma} \mid d \varphi_{(x, y)}$ is not injective $\}$. Then $C$ is semi-algebraic and $\tilde{\Gamma} \backslash C$ is open and dense in $\Gamma$ by remark 6.8. We have $\operatorname{dim} C<\operatorname{dim} \tilde{\Gamma}=d$ (because the image in $\tilde{\Gamma}$ of a ball of dimension $d$ under a map of maximal rank contains an open subset of $\tilde{\Gamma}$, hence cannot be contained in $C$ ).

Let $E^{\prime}$ be the union of $E_{3}$ and the closure in $\mathbb{R}^{n}$ of $\pi(C)=\varphi(C)$. $E^{\prime}$ is a closed semi-algebraic set of dimension $<d$ (since $\left.\operatorname{dim} C<d\right)$. By construction $\varphi \mid \tilde{\Gamma} \backslash \pi^{-1}\left(E^{\prime}\right) \rightarrow E \backslash E^{\prime}$ is a real analytic bijection between manifolds of pure dimension $d$ and has maximal rank, so that its inverse is real analytic. Since $f \mid E \backslash E^{\prime}$ is the composite of this inverse with the projection $\mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{m}$, the lemma is proved.

Note. - For the application we have in mind, it would suffice to show the existence of a semi-algebraic $E^{\prime} \subset \mathbb{R}^{n}$ with $\operatorname{dim} E^{\prime}<d$ such that $f \mid E \backslash E^{\prime}$ is continuous. This can be done more simply, using the structure theorem directly. If $\Gamma \subset \mathbb{R}^{n+1}$ is the graph of a semi-algebraic function $f: E \rightarrow \mathbb{R}$, we have only to make sure that $E^{\prime}$ contains the sets $\bar{A}_{\nu} \backslash A_{\nu}$ for the sets of the form (b) in the structure theorem; sets of the form (a) do not occur if $\Gamma$ is a graph.

Lemma 6.12. - Let $E \subset \mathbb{R}^{n}$ and let $f: E \rightarrow \mathbb{R}^{m}$ be a semialgebraic map. Set $f(E)=F$. Then, there exists a semi-algebraic section $\sigma: F \rightarrow E$; in other words, $\sigma$ is a semi-algebraic map $F \rightarrow E$ such that $f(\sigma(y))=y$ for all $y \in F$.

Proof. - Replacing $E$ by the graph of $f$, we may assume that $E \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ and that $f$ is the projection $\pi: E \rightarrow \mathbb{R}^{m}, \pi(x, y)=y$ for $x \in \mathbb{R}^{n}, y \in \mathbb{R}^{m}$. If we write $y=\left(y_{1}, \ldots, y_{m}\right)$ and $\pi_{1}: E \rightarrow \mathbb{R}^{n} \times \mathbb{R}^{m-1}$ is the projection $\left(x,\left(y_{1}, \ldots, y_{m}\right)\right) \rightarrow\left(x,\left(y_{2}, \ldots, y_{m}\right)\right)$, then, by induction on $m$, it suffices to prove the lemma for $\pi_{1}$ (because of Prop. 6.2 (b)). Thus, by a change of notation, it is sufficient to prove the following statement.

Let $E \subset \mathbb{R}^{N+1}$ be semi-algebraic, and let $F=\pi(E), \pi: \mathbb{R}^{N+1} \rightarrow \mathbb{R}^{N}$ being the projection $\left(x_{1}, \ldots, x_{N+1}\right) \mapsto\left(x_{1}, \ldots, x_{N}\right)$. Then, there exists
a semi-algebraic map $\sigma: F \rightarrow \mathbb{R}^{N+1}$ of the form $\sigma\left(x_{1}, \ldots, x_{N}\right)=$ $\left(x_{1}, \ldots, x_{N}, s\left(x_{1}, \ldots, x_{N}\right)\right)$ with the property that $\sigma(x) \in E$ for all $x \in F$.

We use the structure theorem.
There exists a semi-algebraic partition $\left\{A_{\nu}\right\}$ of $F$, and, for each $\nu$, a finite family $\left\{t_{\nu, j}\right\}_{0 \leq j \leq r_{\nu}+1}$ of functions on $A_{\nu}$ such that $-\infty \equiv t_{\nu, 0}<$ $t_{\nu, 1}<\cdots<t_{\nu, r_{\nu}}<t_{\nu, r_{\nu}+1} \equiv+\infty, t_{\nu, j}$ is continuous and semi-algebraic on $A_{\nu}$ for $1 \leq j \leq r_{\nu}$ and such that $E$ is a finite disjoint union of sets of the form
(a) $\left\{(x, t) \in \mathbb{R}^{N+1} \mid x \in A_{\nu}, t_{\nu, j}(x)<t<t_{\nu, j+1}(x)\right\}, 0 \leq j \leq r_{\nu}$
or of the form
(b) $\left\{(x, t) \in \mathbb{R}^{N+1} \mid x \in A_{\nu}, t=t_{\nu, j}(x)\right\}, 1 \leq j \leq r_{\nu}$.

If $E$ contains a set of the form (b) for a given $\nu$, we define $\sigma \mid A_{\nu}$ by $\sigma\left(x_{1}, \ldots, x_{N}\right)=\left(x_{1}, \ldots, x_{N}, t_{\nu, j}(x)\right), x=\left(x_{1}, \ldots, x_{N}\right) \in A_{\nu}$. If, for a given $\nu, E$ contains no set of the form (b), it contains one of the form (a), and we define $\sigma \mid A_{\nu}$ by $\sigma\left(x_{1}, \ldots, x_{N}\right)=\left(x_{1}, \ldots, x_{N}, \tau(x)\right), x \in A_{\nu}$, where $\tau$ is a semi-algebraic function on $A_{\nu}$ with

$$
t_{\nu, j}(x)<\tau(x)<t_{\nu, j+1}(x), \quad x \in A_{\nu}
$$

[If $r_{\nu}=0$, so that we are dealing with the set $\left\{(x, t) \mid x \in A_{\nu}\right.$, $-\infty<t<+\infty\}$, we take $\tau \equiv 0$ on $A_{\nu}$; if $r_{\nu}>0$ and $j=0$, take $\tau(x)=t_{\nu, 1}(x)-1$; if $r_{\nu}>0$ and $j=r_{\nu}$, take $\tau(x)=t_{\nu, r_{\nu}}(x)+1$; if $1 \leq j \leq r_{\nu}-1$, take $\tau(x)=\frac{1}{2}\left(t_{\nu, j}(x)+t_{\nu, j+1}(x)\right)$.]

Since $\sigma \mid A_{\nu}$ is semi-algebraic for each $\nu, \sigma$ is semi-algebraic on $\cup A_{\nu}=$ $F$. Clearly it has the form required.

Lemma 6.13. - Let $E \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ be a semi-algebraic set, and let $\pi: E \rightarrow \mathbb{R}^{n}$ be the restriction to $E$ of the projection $\mathbb{R}^{n} \times \mathbb{R}^{m} \rightarrow \mathbb{R}^{n}$. Let $F=\pi(E)$.

Assume that for each $x \in F$, the set $E_{x}=\left\{y \in \mathbb{R}^{m} \mid(x, y) \in E\right\}$ is a vector subspace of $\mathbb{R}^{m}$.

Then, we have
(a) If $k \geq 0$ is an integer, the set

$$
F_{k} \subset F, F_{k}=\left\{x \in F \mid \operatorname{dim} E_{x}=k\right\}
$$

is semi-algebraic. $\left(F_{k}=\emptyset\right.$ if $\left.k>m\right)$.
(b) For each $k$, there are semi-algebraic maps $v_{1}, \ldots, v_{k}: F_{k} \rightarrow \mathbb{R}^{m}$ such that, if $x \in F_{k}$, the vectors $v_{1}(x), \ldots, v_{k}(x)$ lie in $E_{x}$ and form a basis of $E_{x}$.

Note. - If we define $v_{\alpha}(x)=0$ for $k<\alpha \leq m, x \in F_{k}$, we can formulate (b) in the following equivalent form:
( $\mathrm{b}^{\prime}$ ) There exist semi-algebraic maps $v_{1}, \ldots, v_{m}: F \rightarrow \mathbb{R}^{m}$ such that, for any $x$, the non-zero vectors among $v_{1}(x), \ldots, v_{m}(x)$ form a basis of $E_{x}$.

## Proof.

(a) It is sufficient to show that for any integer $k \geq 0$, the set

$$
F_{k}^{\prime}=\left\{x \in F \mid \operatorname{dim} E_{x} \geq k\right\}
$$

is semi-algebraic. We have $F_{0}^{\prime}=F$. Let $k \geq 1$.
Let $\left(\mathbb{R}^{m}\right)^{k}=\mathbb{R}^{m} \times \cdots \times \mathbb{R}^{m}(k$-times $)$ and let $S_{k} \subset\left(\mathbb{R}^{m}\right)^{k}$ be the set

$$
S_{k}=\left\{\left(v_{1}, \ldots, v_{k}\right) \in\left(\mathbb{R}^{m}\right)^{k} \mid v_{1} \wedge \cdots \wedge v_{k} \neq 0 \text { in } \stackrel{k}{\wedge} \mathbb{R}^{m}\right\}
$$

Clearly, $S_{k}$ is semi-algebraic.
Consider now the set

$$
V_{k} \subset \mathbb{R}^{n} \times\left(\mathbb{R}^{m}\right)^{k}:\left\{\left(x, v_{1}, \ldots, v_{k}\right) \mid\left(x, v_{j}\right) \in E \text { for } 1 \leq j \leq k\right\}
$$

Again, $V_{k}$ is a semi-algebraic set. Now $\operatorname{dim} E_{k} \geq k$ if and only if there exist $v_{1}, \ldots, v_{k} \in E_{x}$ with $v_{1} \wedge \cdots \wedge v_{k} \neq 0$. Hence

$$
\begin{aligned}
& F_{k}^{\prime}=\left\{x \in F \mid \operatorname{dim} E_{x} \geq k\right\}=\left\{x \in F \mid \exists\left(x, v_{1}, \ldots, v_{k}\right) \in V_{k}\right. \\
&\text { such that } \left.\left(v_{1}, \ldots, v_{k}\right) \in S_{k}\right\} .
\end{aligned}
$$

Hence $F_{k}^{\prime}$ is the image under the projection of $\mathbb{R}^{n} \times\left(\mathbb{R}^{m}\right)^{k}$ onto $\mathbb{R}^{n}$, of the semi-algebraic set $V_{k} \cap\left(\mathbb{R}^{n} \times S_{k}\right)$, and so is semi-algebraic (by the Tarski-Seidenberg theorem).
(b) We have $\operatorname{dim} E_{x}=k$ for $x \in F_{k}$. Consider the set $A_{k} \subset \mathbb{R}^{n} \times\left(\mathbb{R}^{m}\right)^{k}$ defined by

$$
\begin{aligned}
& A_{k}=\left\{\left(x, v_{1}, \ldots, v_{k}\right) \mid x \in F_{k},\left(x, v_{j}\right) \in E\right. \\
& \left.\qquad \text { for } 1 \leq j \leq k, \text { and } v_{1} \wedge \cdots \wedge v_{k} \neq 0\right\}
\end{aligned}
$$

Clearly $A_{k}$ is semi-algebraic and $\left(x, v_{1}, \ldots, v_{k}\right) \in A_{k}$ if and only if $\left(v_{1}, \ldots, v_{k}\right)$ is a basis of $E_{x}$.

Now, the projection $\varphi: A_{k} \rightarrow \mathbb{R}^{n},\left(x, v_{1}, \ldots, v_{k}\right) \mapsto x$, maps $A_{k}$ onto $F_{k}$. By Lemma 6.12, $\varphi$ has a semi-algebraic section $\sigma: F_{k} \rightarrow A_{k}$ (i.e. $\sigma$ is semi-algebraic and $\varphi(\sigma(x))=x$ for $\left.x \in F_{k}\right)$. If we write

$$
\sigma(x)=\left(x, v_{1}(x), \ldots, v_{k}(x)\right), x \in F_{k}
$$

then the $v_{\alpha}: F_{k} \rightarrow \mathbb{R}^{m}, 1 \leq \alpha \leq k$, are semi-algebraic and $\left(v_{1}(x), \ldots, v_{k}(x)\right)$ forms a basis of $E_{x}$ for $x \in F_{k}$ (since $\left.\left(x, v_{1}(x), \ldots, v_{k}(x)\right) \in A_{k}\right)$.

Lemma 6.14. - Let $E$ be a semi-algebraic subset of $\mathbb{R}^{n} \times \mathbb{R}^{m}$ and define $E^{\perp}=\left\{(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{m} \mid\langle y, u\rangle=0\right.$ for all $u \in \mathbb{R}^{m}$ with $\left.(x, u) \in E\right\}$; here $\langle y, u\rangle=\sum y_{j} u_{j}, y=\left(y_{1}, \ldots, y_{m}\right), u=\left(u_{1}, \ldots, u_{m}\right)$, is the usual inner product on $\mathbb{R}^{m}$.

Then $E^{\perp}$ is again semi-algebraic.

$$
\begin{aligned}
& \text { Proof. - Let } F \subset \mathbb{R}^{n} \times \mathbb{R}^{m} \times \mathbb{R}^{m} \text { be the set } \\
& \quad F=\{(x, y, u) \mid(x, u) \in E \text { and }\langle y, u\rangle=0\}
\end{aligned}
$$

Then $F$ is semi-algebraic and we have

$$
E^{\perp}=\left\{(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{m} \mid(x, y, u) \in F \text { for all } u \text { with }(x, u) \in E\right\}
$$

This latter set is semi-algebraic since its complement is the projection in $\mathbb{R}^{n} \times \mathbb{R}^{m}$ of the semi-algebraic set

$$
\left\{(x, y, u) \in \mathbb{R}^{n} \times \mathbb{R}^{m} \times \mathbb{R}^{m} \mid(x, u) \in E,(x, y, u) \notin F\right\}
$$

Finally, we note two results proved in [FN] (Lemma 5.2 and Cor. 3 to Theorem 5.1 in that paper).

Lemma 6.15. - Let $S$ be a compact semi-algebraic set in $\mathbb{R}^{n}$ and let $E \subset S$ be a semi-algebraic subset. Let $f: E \rightarrow \mathbb{R}$ be a semi-algebraic function which is locally bounded on $E$.

Then, there exist constants $m, C>0$ such that

$$
|f(x)| \leq C(\operatorname{dist}(x, S \backslash E))^{-m} \text { for } x \in E
$$

(dist $(x, S \backslash E)$ is, of course, the distance of $x$ from $S \backslash E$.)
Lemma 6.16. - Let $E \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ be semi-algebraic, and let $f: E \rightarrow \mathbb{R}$ be a semi-algebraic function. For $x \in \mathbb{R}^{n}$, define $\tilde{f}(x)=$ $\sup \{f(x, y) \mid(x, y) \in E\}$ (the sup over the empty set being defined to be $-\infty$ ).

Then $\tilde{f}$ is an extended semi-algebraic function.

## 7. Some special semi-algebraic sets and maps.

Lemma 7.1. - Let $E_{1}, E_{2} \subset \mathbb{R}^{n} \times \mathbb{R}^{m}$ be semi-algebraic sets, and let

$$
E=\left\{(x, y) \in \mathbb{R}^{n} \times \mathbb{R}^{m} \mid y=y_{1}+y_{2} \text { where }\left(x, y_{1}\right) \in E_{1},\left(x, y_{2}\right) \in E_{2}\right\}
$$

Then $E$ is semi-algebraic.

Proof. - The set $\left\{\left(x, y_{1}, y_{2}\right) \in \mathbb{R}^{n} \times \mathbb{R}^{m} \times \mathbb{R}^{m} \mid\left(x, y_{1}\right) \in E_{1}\right.$, $\left.\left(x, y_{2}\right) \in E_{2}\right\}$ is semi-algebraic and $E$ is the image of this set under the $\operatorname{map}\left(x, y_{1}, y_{2}\right) \mapsto\left(x, y_{1}+y_{2}\right)$.

For our next lemma, recall that $W \subset H^{D} \times \cdots \times H^{D}$ ( $r$ times) is the space of $r$-tuples $P=\left(p_{1}, \ldots, p_{r}\right)$ of polynomials $p_{j} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$, $\operatorname{deg} p_{j} \leq D,\|P\| \leq C_{1}, P(0)=0, J_{P}(0)=1$ (where $J_{P}(x)=$ $\left.\operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}\right)_{1 \leq j, k \leq r}\right)$. Clearly $W$ is a semi-algebraic subset of $H^{D} \times \cdots \times H^{D}$.

We have seen that there is $\rho_{1}>0$ such that whenever $P \in W$, $f \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ and $f=0$ on some neighborhood of 0 in $Z(P)=\{x \in$ $\left.\mathbb{R}^{n} \mid P(x)=0\right\}$ (i.e. $f=0$ as a germ in $\mathcal{G}(P)$ ), then $f(x)=0$ for any $x$ with $\left|x_{j}\right| \leq \rho_{1}, P(x)=0$.

Lemma 7.2. - Let $d \geq 1$ be a given integer. The set

$$
Z_{d}=\left\{(P, f) \in W \times H^{d} \mid f=0 \text { as a germ in } \mathcal{G}(P)\right\}
$$

is semi-algebraic.

Proof. - Let $\rho_{1}>0$ be as above. We first remark that the set

$$
X=\left\{(P, f, x) \in W \times H^{d} \times Q_{\rho_{1}} \mid f(x)=0 \text { or } P(x) \neq 0\right\}
$$

is semi-algebraic. In fact, $X$ is the inverse image of the semi-algebraic set

$$
\begin{aligned}
& S=\left\{\left(x_{1}, \ldots, x_{n}, t\right) \in \mathbb{R}^{n} \times \mathbb{R} \mid t=0\right\} \cup\left\{\left(x_{1}, \ldots, x_{n}, t\right) \in \mathbb{R}^{n} \times \mathbb{R} \mid\right. \\
&\left.x_{j} \neq 0 \text { for some } j, 1 \leq j \leq n\right\}
\end{aligned}
$$

- under the semi-algebraic map $W \times H^{d} \times Q_{\rho_{1}} \rightarrow \mathbb{R}^{r} \times \mathbb{R}$ given by $(P, f, x) \mapsto$ $(P(x), f(x))$ and so is semi-algebraic (Prop. 6.2 (a)). Since $Z_{d}=\{(P, f) \in$ $\left.W \times H^{d} \mid(P, f, x) \in X, \forall x \in Q_{\rho_{1}}\right\}, Z_{d}$ is semi-algebraic by Proposition 6.2 (d).

We next define two norms.

## Definition 7.3.

(a) Let $d \geq 1$ be an integer, let $P \in W, q \in H^{d}, 0<\rho \leq \rho_{1}$. We define the norm $N$ by $N(P, d, q, \rho)=$ infimum of the real numbers $C \geq 0$ such that, for any $f \in H^{D}$, there exists $F \in H^{d}$ for which $F-q f=0$ as a germ in $\mathcal{G}(P)$ and $\|F\| \leq C \sup _{V_{\rho}(P)}|f|$; if no such $C$ exists, we set $N(P, d, q, \rho)=+\infty$.
(b) Let $d \geq 1$ be an integer, let $C \geq 0$ and $0<\rho \leq \rho_{1}$. If $P \in W$, we define the norm $N_{*}$ by

$$
N_{*}(P, d, C, \rho)=\inf \left\{N(P, d, q, \rho) \mid q \in H^{d}, q(0)=1,\|q\| \leq C\right\}
$$

Lemma 7.4. - The two norms defined above are extended semialgebraic functions of their arguments. More precisely, given $d \geq 1$, the maps
(a)

$$
\begin{aligned}
N: W \times H^{d} \times\left(0, \rho_{1}\right] & \rightarrow \mathbb{R} \cup\{\infty\} \\
(P, q, \rho) & \mapsto N(P, d, q, \rho)
\end{aligned}
$$

and
(b)

$$
\begin{aligned}
N_{*}: W \times \mathbb{R}^{+} \times\left(0, \rho_{1}\right] & \rightarrow \mathbb{R} \cup\{\infty\}, \mathbb{R}^{+}=\{C \in \mathbb{R} \mid C \geq 0\} \\
(P, C, \rho) & \mapsto N_{*}(P, d, C, \rho)
\end{aligned}
$$

are extended semi-algebraic functions.

## Proof.

(a) We first note that the set

$$
\begin{align*}
& E_{1}=\left\{(P, F, f, q) \in W \times H^{d} \times H^{D} \times H^{d} \mid F-q f=0\right.  \tag{1}\\
&\text { as a germ in } \mathcal{G}(P)\}
\end{align*}
$$

is semi-algebraic. In fact, it is the inverse image of the semi-algebraic set $Z_{d+D} \subset W \times H^{d+D}$ (Lemma 7.2) under the semi-algebraic map $(P, F, f, q) \mapsto(P, F-q f)$ of $W \times H^{d} \times H^{D} \times H^{d}$ into $W \times H^{d+D}$.

Next, we show that the map

$$
\begin{equation*}
\varphi: W \times H^{D} \times\left(0, \rho_{1}\right] \rightarrow \mathbb{R},(P, f, \rho) \mapsto \sup _{V_{\rho}(P)}|f| \tag{2}
\end{equation*}
$$

is semi-algebraic. To see this, remark that the set

$$
\begin{aligned}
& E_{2}=\left\{(P, f, x, \rho) \in W \times H^{D} \times Q_{\rho_{1}} \times\left(0, \rho_{1}\right] \mid P(x)=0\right. \\
&\left.\left|x_{j}\right| \leq \rho, j=1, \ldots, n\right\}
\end{aligned}
$$

is clearly semi-algebraic, as is the map $\psi: W \times H^{D} \times Q_{\rho_{1}} \times\left(0, \rho_{1}\right] \rightarrow \mathbb{R}$ given by

$$
\psi(P, f, x, \rho)=|f(x)|
$$

Now, given $(P, f, \rho)$, we have

$$
\sup _{V_{\rho}(P)}|f|=\sup \left\{\psi(P, f, x, \rho) \text { over all } x \text { such that }(P, f, x, \rho) \in E_{2}\right\}
$$

Further, since $V_{\rho}(P)$ is compact, this sup is finite. Hence (2) follows from Lemma 6.16.

It follows from (2) that the set

$$
E_{3}=\left\{(P, F, f, \rho, C) \in W \times H^{d} \times H^{D} \times\left(0, \rho_{1}\right] \times \mathbb{R}^{+}\left|\|F\| \leq C \sup _{V_{\rho}(P)}\right| f \mid\right\}
$$

is semi-algebraic, so that (see (1)), so is the set

$$
\begin{aligned}
E_{4}=\left\{(P, F, f, \rho, C, q) \in W \times H^{d}\right. & \times H^{D} \times\left(0, \rho_{1}\right] \times \mathbb{R}^{+} \times H^{d} \mid \\
& \left.(P, F, f, \rho, C) \in E_{3} \text { and }(P, F, f, q) \in E_{1}\right\}
\end{aligned}
$$

Hence, by Proposition 6.2 (c), (d), the set

$$
\begin{aligned}
& E_{5}=\left\{(P, \rho, C, q) \in W \times\left(0, \rho_{1}\right] \times \mathbb{R}^{+} \times H^{d} \mid \forall f \in H^{D}, \exists F \in H^{d}\right. \\
&\text { such that } \left.(P, F, f, \rho, C, q) \in E_{4}\right\}
\end{aligned}
$$

is semi-algebraic. Finally,

$$
N(P, d, q, \rho)=\inf \left\{C \in \mathbb{R}^{+} \mid(P, \rho, C, q) \in E_{5}\right\}
$$

so that $N$ is an extended semi-algebraic function by Lemma 6.16.
(b) The set

$$
E_{6}=\left\{(P, C, \rho, q) \in W \times \mathbb{R}^{+} \times\left(0, \rho_{1}\right] \times H^{d} \mid q \in H^{d}, q(0)=1,\|q\| \leq C\right\}
$$

is clearly semi-algebraic. By (a), so is the set

$$
E_{7}=\left\{(P, C, \rho, q) \in E_{6} \mid N(P, d, q, \rho)<\infty\right\}
$$

If $E_{0}$ is the image of $E_{7}$ under the projection $(P, C, \rho, q) \mapsto(P, C, \rho)$, then, for $(P, C, \rho) \in E_{0}$,

$$
N_{*}(P, d, C, \rho)=\inf \left\{N(P, d, q, \rho) \mid(P, C, \rho, q) \in E_{7}\right\}
$$

while $N_{*}(P, d, C, \rho)=\infty$ if $(P, C, \rho) \notin E_{0}$. Hence $N_{*}$ is an extended semialgebraic function by Lemma 6.16.

## 8. The induction scheme and some technical lemmas.

We recall the notation that we have been using. Let $n \geq 2,1 \leq r \leq$ $n-1, D \geq 1$ be given integers. Let $C_{1}>0$. We have set:

$$
H^{D}=\left\{p \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right] \mid \operatorname{deg} p \leq D\right\}
$$

$W=\left\{P=\left(p_{1}, \ldots, p_{r}\right) \mid p_{j} \in H^{D}, P(0)=0,\|P\| \leq C_{1}\right.$,

$$
\left.\operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}(0)\right)_{1 \leq j, k \leq r}=1\right\}
$$

where, if $p=\sum p_{\alpha} x^{\alpha} \in H^{d}, d \geq 1,\|p\|^{2}=\sum\left|p_{\alpha}\right|^{2}$ and $\|P\|^{2}=\sum_{j}\left\|p_{j}\right\|^{2}$.
Fix $\rho_{1}>0$ such that, for any $P \in W$, the set

$$
\left\{x \in \mathbb{R}^{n}| | x_{j} \mid \leq \rho_{1}, p_{1}(x)=\cdots=p_{r}(x)=0\right\}
$$

is contained in the connected component of $Z(P) \cap\left\{x \in \mathbb{R}^{n} \mid J_{P}(x) \neq\right.$ $0\}=\left\{x \in \mathbb{R}^{n} \mid P(x)=0, \operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}\right) \neq 0\right\}$ which contains the origin (see $\S 1,(4)$ and the remark following Lemma 2.1).

For $0<\rho \leq \rho_{1}$, and $P \in W$, we define

$$
V_{\rho}(P)=Q_{\rho} \cap Z(P)=\left\{x \in \mathbb{R}^{n}| | x_{j} \mid \leq \rho, p_{1}(x)=\cdots=p_{r}(x)=0\right\}
$$

Recall that $\mathcal{G}(P)$ is the space of germs at 0 of functions defined on some neighborhood of 0 in $Z(P)$ [or $V(P)$ or $\left.V_{\rho_{1}}(P)\right]$.

We shall use, both in this section and the next, the following lemma proved in [FN], and so state it here.

Lemma 8.1. - Let $d \geq 1, E \subset \mathbb{R}^{n}$ and $\varphi_{\alpha} \in H^{d}, \alpha=1, \ldots, s$. Assume that $\varphi_{1}\left|E, \ldots, \varphi_{s}\right| E$ are linearly independent over $\mathbb{R}$.

Then, there exist constants $\epsilon>0, K>0, K^{\prime}>0$ (depending only on $E$ and the $\varphi_{\alpha}$ ) and points $x_{\beta} \in E, \beta=1, \ldots, s$, such that the following holds.

$$
\text { If } \tilde{x}_{\beta} \in \mathbb{R}^{n},\left|\tilde{x}_{\beta}-x_{\beta}\right|<\epsilon, \beta=1, \ldots, s, \text { and } \tilde{\varphi}_{\alpha} \in H^{d},\left\|\tilde{\varphi}_{\alpha}-\varphi_{\alpha}\right\|<\epsilon
$$ $\alpha=1, \ldots, s$, then, for any $A_{\alpha} \in \mathbb{R}, \alpha=1, \ldots, s$, we have

$$
\sum_{1 \leq j \leq s}\left|A_{j}\right| \leq K \max _{\beta}\left|\sum_{\alpha=1}^{s} A_{\alpha} \tilde{\varphi}_{\alpha}\left(\tilde{x}_{\beta}\right)\right|
$$

In particular,

$$
\left\|\sum_{1 \leq j \leq s} A_{j} \tilde{\varphi}_{j}\right\| \leq K^{\prime} \max _{\beta}\left|\sum_{\alpha=1}^{s} A_{\alpha} \tilde{\varphi}_{\alpha}\left(\tilde{x}_{\beta}\right)\right|
$$

The extension theorem stated in $\S 1$ will be proved by proving the following:

Main Lemma 8.2. - Fix an integer $t$ with $-1 \leq t \leq \operatorname{dim} W$.
There exist constants $C_{t}, D_{t}>0$ depending only on $t$ and the given constants $n, r, D, C_{1}$, and a closed semi-algebraic set $W^{t} \subset W$ with the following properties:
(A) $\operatorname{dim} W^{t} \leq t$ (so that $W^{-1}=\emptyset$ ).
(B) Given $P_{0} \in W \backslash W^{t}$ and $\rho>0, \rho \leq \rho_{1}$, there exist constants $\delta>0$ and $K>0$ depending only on $P_{0}$ and $\rho$ such that: for any $P \in W$ with $\left\|P-P_{0}\right\|<\delta$, we can find $q \in H^{D_{t}}$ satisfying the following requirements:
(i) $q(0)=1,\|q\| \leq C_{t}$;
(ii) For any $f \in H^{D}$, there exists $F \in H^{D_{t}}$ so that $q f=F$ as germs in $\mathcal{G}(P)$ and $\|F\| \leq K \sup |f|$.

$$
V_{\rho}(P)
$$

If $t=\operatorname{dim} W$, this statement is trivial; we have only to take $W^{t}=W$.
In the rest of this section and in all of $\S 9$, we assume that $t$ is fixed and that the main lemma holds for the number $t$.

Thus, we assume that the set $W^{t}$ and the constants $C_{t}, D_{t}$ are given.
We use the following:
Convention. - Constants written $C, C^{\prime}, C_{\#}, \underline{C}, \ldots, D, D^{\prime}$ (not involving auxiliary parameters) will be understood to depend only on $n, r, D, C_{1}$ unless otherwise stated. Constants written as $C(d), \bar{C}(m), \ldots$ will depend only on $n, r, D, C_{1}$ and the auxiliary parameters indicated (such as $d, m)$. Constants written as $\delta, \tilde{\delta}_{1}, \ldots, K, K^{\prime}, \ldots$ may depend on other data which will be indicated explicitly.

For $d \geq 1$, consider the set

$$
\begin{equation*}
Z_{d}=\left\{(P, f) \in W \times H^{d} \mid f=0 \text { as a germ in } \mathcal{G}(P)\right\} \tag{1}
\end{equation*}
$$

The set $Z_{d}$ is semi-algebraic by Lemma 7.2.
For $C>0, d \geq 1$, define

$$
\begin{equation*}
S(C, d)=\left\{q \in H^{d} \mid q(0)=1,\|q\| \leq C\right\} \tag{2}
\end{equation*}
$$

This set is clearly semi-algebraic for any $C, d$.

Lemma 8.3. - There exist constants $C_{*}>0, D_{*} \geq 1$, and, for $d \geq 1$, a function $d_{*} \geq 1$ of $d, n, r, D, C_{1}$ such that the following holds.

There exists a semi-algebraic map

$$
Q: W \rightarrow S\left(C_{*}, D_{*}\right)
$$

(independent of $d$ ) and, for $d \geq 1$, semi-algebraic maps

$$
g_{j}: Z_{d} \rightarrow H^{d_{*}}, j=1, \ldots, r
$$

such that

$$
Q(P) f=\sum_{j=1}^{r} g_{j}(P, f) p_{j}\left(\text { in } \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]\right) \text { for all }(P, f) \in Z_{d}
$$

Proof. - Let $D_{0}>0$. Let $E$ be the following subset of $W \times H^{D_{0}}$. (3) $E$ is the set of pairs $(P, q) \in W \times H^{D_{0}}$ such that for any $d \geq 1$, if $f \in H^{d}$ and $f=0$ as a germ in $\mathcal{G}(P)$, then $q f=\sum g_{j} p_{j}$ for some $g_{1}, \ldots, g_{r} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$.

We claim that $E$ is semi-algebraic. To prove this claim, we use Theorem $5.6^{\prime}$, which shows that if $P \in W$, the ideal in $\mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ of those $f$ which $=0$ as a germ in $\mathcal{G}(P)$ is generated by $f_{1}, \ldots, f_{r_{*}}$ with $\operatorname{deg} f_{\nu} \leq D^{\prime}$, $1 \leq \nu \leq r_{*}$, where $D^{\prime}$ depends only on $n, r, D$. Clearly, $E$ is the set of pairs $(P, q)$ such that $q f_{\nu}=\sum_{j} g_{j}^{(\nu)} p_{j}, \nu=1, \ldots, r_{*}, g_{j}^{(\nu)} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$. By Theorem 5.4 if the $g_{j}^{(\nu)}$ exist, we may assume that $\operatorname{deg}\left(g_{j}^{(\nu)}\right) \leq D^{\prime \prime}$, where $D^{\prime \prime}$ depends only on $D_{0}, D^{\prime}, n, r, D$.

Hence, we have:
(4) $E=\left\{(P, q) \in W \times H^{D_{0}} \mid\right.$ if $f \in H^{D^{\prime}}$ and $f=0$ as a germ in $\mathcal{G}(P)$, then there exist $g_{1}, \ldots, g_{r} \in H^{D^{\prime \prime}}$ such that $\left.q f=\sum g_{j} p_{j}\right\}$.

The set $E_{1} \subset W \times H^{D_{0}} \times H^{D^{\prime}} \times\left(H^{D^{\prime \prime}}\right)^{r}$ defined by

$$
E_{1}=\left\{\left(P, q, f,\left(g_{1}, \ldots, g_{r}\right)\right) \mid q f=\sum g_{j} p_{j}\right\}
$$

is clearly semi-algebraic. Hence its projection on $W \times H^{D_{0}} \times H^{D^{\prime}}: E_{2}=$ $\left\{(P, q, f) \in W \times H^{D_{0}} \times H^{D^{\prime}} \mid \exists g \in\left(H^{D^{\prime \prime}}\right)^{r}\right.$ with $\left.q f=\sum g_{j} p_{j}\right\}$ is semialgebraic.

We have:
$E=\left\{(P, q) \in W \times H^{D_{0}} \mid \forall f\right.$ with $(P, f) \in Z_{D^{\prime}}$, we have $\left.(P, q, f) \in E_{2}\right\}$.
This set is semi-algebraic by Prop. 6.2, (d), thus proving our claim.

Now, by Theorem 5.5 , there exist $C_{*}, D_{*} \geq 1$ such that if $E$ is the set (3) defined with $D_{0}=D_{*}$, then the projection

$$
E \cap W \times S\left(C_{*}, D_{*}\right) \rightarrow W
$$

is onto $W$. By Lemma 6.12, we can find a semi-algebraic section of this projection, say $\widetilde{Q}: W \rightarrow E \cap W \times S\left(C_{*}, D_{*}\right)$; set $\widetilde{Q}(P)=(P, Q(P))$. By the defining property of $E$, we have:
(5) If $P \in W, f \in H^{d}$ and $f=0$ as a germ in $\mathcal{G}(P)$, there exist $g_{1}, \ldots, g_{r} \in \mathbb{R}\left[x_{1}, \ldots, x_{n}\right]$ such that $Q(P) f=\sum g_{j} p_{j}$.

By Theorem 5.4, there is a function $d_{*}$ of $d, n, r, D, C_{1}$ such that:
(5') If $P \in W, f \in H^{d}$ and $(P, f) \in Z_{d}$, then there exist $g_{1}, \ldots, g_{r} \in H^{d_{*}}$ so that $Q(P) f=\sum g_{j} p_{j}$.

In other words, if $\tilde{E} \subset W \times H^{d} \times\left(H^{d_{*}}\right)^{r}$ is the semi-algebraic set defined by

$$
\tilde{E}=\left\{\left(P, f,\left(g_{1}, \ldots, g_{r}\right)\right) \mid Q(P) f=\sum g_{j} p_{j}\right\}
$$

then projection onto $W \times H^{d}$ maps $\tilde{E}$ onto $Z_{d}$. Again by Lemma 6.12 , this map has a semi-algebraic section. This means precisely that there are semi-algebraic maps $g_{j}: Z_{d} \rightarrow H^{d_{*}}, j=1, \ldots, r$, such that $\left(P, f,\left(g_{1}(P, f), \ldots, g_{r}(P, f)\right)\right) \in \tilde{E}$ for any $(P, f) \in Z_{d}$, i.e. such that $Q(P) f=\sum_{j} g_{j}(P, f) p_{j}$. This is Lemma 8.3.

We now proceed to the principal technical step which enables us to prove the Main Lemma 8.2.

Technical Lemma 8.4. - Let $W^{t}, C_{t}, D_{t}$ be as in the Main Lemma, and let $Q: W \rightarrow S\left(C_{*}, D_{*}\right)$ be the semi-algebraic map constructed in Lemma 8.3

Let $m \geq 0$ be a given integer.
We can find a semi-algebraic partition $\left\{Y_{\nu}\right\}$ of $W^{t}$, a closed semialgebraic set $Z \subset W^{t}$ of dimension $\leq t-1(Z=\emptyset$ if $t=0)$, a constant $D(m) \geq 1$, and, for each $\nu$, a finite number of semi-algebraic maps $\Phi_{\alpha}: Y_{\nu} \rightarrow H^{D(m)}$ (we do not indicate their dependence on $\nu$ in the notation) with the following properties:
( $\mathrm{I}_{m}$ ) If $P \in Y_{\nu}$, the polynomials $\Phi_{\alpha}(P)$ are linearly independent as elements of $\mathcal{G}(P)$.
$\left(\mathrm{II}_{m}\right)$ Let $P_{0} \in Y_{\nu} \backslash Z$. Then, there exist constants $\delta>0, K>0$ (depending only on $P_{0}$ and the above data) such that whenever we are
given:
$P \in Y_{\nu}$ with $\left\|P-P_{0}\right\|<\delta ; g_{1}, \ldots, g_{r} \in H^{D}$ with $\left\|g_{i}\right\| \leq 1 ;$
a real number $\tau, 0<\tau<1$; and $f \in H^{D_{t}}$,
then, we can find real numbers $A_{\alpha}$ and a polynomial $\tilde{f} \in H^{D(m)}$ for which we have

$$
\begin{align*}
(Q(P))^{m} f=\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)+\tau^{m} \tilde{f} \bmod \left(p_{1}+\tau g_{1}, \ldots, p_{r}+\tau g_{r}\right)  \tag{i}\\
P=\left(p_{1}, \ldots, p_{r}\right)
\end{align*}
$$

and

$$
\begin{equation*}
\|\tilde{f}\| \leq K\|f\| \tag{ii}
\end{equation*}
$$

Proof of the Technical Lemma. - The proof is by induction on $m$. For $m=0$, the result is obvious: take $\left\{Y_{\nu}\right\}$ to consist of $W^{t}$ alone, $f=\tilde{f}$, and $\left\{\Phi_{\alpha}\right\}$ to be the empty family.

Assume therefore that $m,\left\{Y_{\nu}\right\}, Z, \Phi_{\alpha}: Y_{\nu} \rightarrow H^{D(m)}$ are all given with the properties ( $\mathrm{I}_{m}$ ) and ( $\mathrm{II}_{m}$ ) above.

For $P \in Y_{\nu}$, let $V_{P}^{\prime}=\left\{f \in H^{D(m)} \mid f=0\right.$ as a germ in $\mathcal{G}(P)\}=\left\{f \in H^{D(m)} \mid(P, f) \in Z_{D(m)}\right\}$, and let $V_{P}^{\prime \prime}$ be the orthogonal complement of $V_{P}^{\prime} \oplus \bigoplus_{\alpha} \mathbb{R} \Phi_{\alpha}(P)\left[\bigoplus_{\alpha}\right.$ is the linear span of the $\left.\Phi_{\alpha}(P)\right]$. The set

$$
\left\{(P, f) \in Y_{\nu} \times H^{D(m)} \mid f \in V_{P}^{\prime}\right\}=Z_{D(m)} \cap\left(Y_{\nu} \times H^{D(m)}\right)
$$

is semi-algebraic. By Lemma 7.2 and Lemma 6.14, the set

$$
\left\{(P, f) \in Y_{\nu} \times H^{D(m)} \mid f \in V_{P}^{\prime \prime}\right\}
$$

is also semi-algebraic. We apply Lemma 6.13 to these sets to obtain the following:
(6) There is a semi-algebraic partition $\left\{Y_{\mu}^{\prime}\right\}$ of $W^{t}$, refining the partition $\left\{Y_{\nu}\right\}$, such that, for each $\mu$, we can find finitely many semi-algebraic mappings $h_{\beta}: Y_{\mu}^{\prime} \rightarrow H^{D(m)}, \psi_{\gamma}: Y_{\mu}^{\prime} \rightarrow H^{D(m)}$ with the property that for $P \in Y_{\mu}^{\prime}$, the $\left\{h_{\beta}(P)\right\}$ form a basis of $V_{P}^{\prime}$, the $\left\{\psi_{\gamma}(P)\right\}$ form a basis of $V_{P}^{\prime \prime}$.

Moreover, if we restrict the $\Phi_{\alpha}$ from $Y_{\nu}$ to the $Y_{\mu}^{\prime}$ contained in $Y_{\nu}$, properties ( $\mathrm{I}_{m}$ ) and ( $\mathrm{II}_{m}$ ) continue to hold.

Note that for $P \in Y_{\mu}^{\prime}$, the polynomials $\left\{h_{\beta}(P), \psi_{\gamma}(P), \Phi_{\alpha}(P)\right\}$ form a basis of $H^{D(m)}$. In fact, since $h_{\beta}(P)=0$ in $\mathcal{G}(P)$, and the images
of the $\Phi_{\alpha}(P)$ are linearly independent in $\mathcal{G}(P),\left\{h_{\beta}(P), \Phi_{\alpha}(P)\right\}$ form a basis of $V_{P}^{\prime} \oplus \bigoplus_{\alpha} \mathbb{R} \Phi_{\alpha}(P)$, while the $\psi_{\gamma}(P)$ form a basis of the orthogonal complement of this space in $H^{D(m)}$.

Since $\left(P, h_{\beta}(P)\right) \in Z_{D(m)}$ for $P \in Y_{\mu}^{\prime}$ and any $\beta$, Lemma 8.3 can be applied, and provides semi-algebraic maps

$$
g_{\beta}^{j}: Y_{\mu}^{\prime} \rightarrow H^{d(m)}
$$

[where $d(m)$ depends only on $D(m)$ and the constants defining $W$ ] such that

$$
\begin{equation*}
Q(P) h_{\beta}=\sum_{j=1}^{r} g_{\beta}^{j}(P) p_{j}, P \in Y_{\mu}^{\prime} \tag{7}
\end{equation*}
$$

( $Q$ being again the map constructed in Lemma 8.3).
Given $P \in Y_{\mu}^{\prime}$ and $F \in H^{D(m)}$, we can write

$$
\begin{equation*}
F=\sum_{\alpha} a_{\alpha} \Phi_{\alpha}(P)+\sum_{\beta} b_{\beta} h_{\beta}(P)+\sum_{\gamma} e_{\gamma} \psi_{\gamma}(P) \tag{8}
\end{equation*}
$$

where $a_{\alpha}, b_{\beta}, e_{\gamma}$ are real numbers depending on $P$ and $F$ (they are semialgebraic functions of $P$ and $F$, as is easily proved, but we do not need this fact).

Multiplying (8) by $Q(P)$ and using (7), we obtain

$$
\begin{align*}
& Q(P) F=\sum_{\alpha} a_{\alpha} Q(P) \Phi_{\alpha}(P)+\sum_{\gamma} e_{\gamma} Q(P) \psi_{\gamma}(P)  \tag{9}\\
&+\sum_{j=1}^{r}\left(\sum_{\beta} b_{\beta} g_{\beta}^{j}(P)\right) p_{j}
\end{align*}
$$

Suppose that in addition to $P \in Y_{\mu}^{\prime}$ and $F \in H^{D(m)}$, we are given $\tau \in(0,1)$ and $g_{1}, \ldots, g_{r} \in H^{D}$ with $\left\|g_{j}\right\| \leq 1$. Then (9) implies

$$
\begin{align*}
Q(P) F \equiv \sum_{\alpha} a_{\alpha} Q(P) \Phi_{\alpha}(P)+ & \sum_{\gamma} e_{\gamma} Q(P) \psi_{\gamma}(P)  \tag{10}\\
& +\tau \underline{F} \bmod \left(p_{1}+\tau g_{1}, \ldots, p_{r}+\tau g_{r}\right)
\end{align*}
$$

where

$$
\underline{F}=-\sum_{j=1}^{r}\left(\sum_{\beta} b_{\beta} g_{\beta}^{j}(P)\right) g_{j}
$$

We can now define sets $\left\{\underline{Y_{\nu}}\right\}, \underline{Z}$, a constant $D(m+1)$, and, for each $\nu$, maps $\underline{\Phi_{\alpha}}: \underline{Y_{\nu}} \rightarrow H^{D(m+1)}$ which have properties $\left(\mathrm{I}_{m+1}\right)$ and $\left(\mathrm{II}_{m+1}\right)$ of Lemma 8.4.
(a) We take the partition $\left\{\underline{Y_{\nu}}\right\}$ of $W^{t}$ be the partition $\left\{Y_{\mu}^{\prime}\right\}$ above.
(b) $\underline{Z}$ will be any closed semi-algebraic set of dimension $\leq t-1$ such that $\underline{Z} \supset Z$ and the maps $Q, \Phi_{\alpha}, h_{\beta}, \psi_{\gamma}, g_{\beta}^{j}$ are continuous on $\underline{Y_{\nu}} \backslash \underline{Z}$.

Note that we can find such a $\underline{Z}$ by Lemma 6.11 and Lemmas 6.6, 6.7 since $\operatorname{dim} \underline{Y}_{\nu} \leq \operatorname{dim} W^{t} \leq t$ and $\operatorname{dim} Z \leq t-1$ by assumption.
(c) $D(m+1)=\max \left(D(m)+D_{*}, d(m)+D\right)$ with $D_{*}$ as in Lemma 8.3, $\operatorname{deg} g_{\beta}^{j}(P) \leq d(m)$ and $\left\{\underline{\Phi}_{\alpha}\right\}$ is an enumeration of the maps

$$
P \mapsto Q(P) \Phi_{\alpha}(P) \text { and } P \mapsto Q(P) \psi_{\gamma}(P)
$$

This completes our choices. It remains to show that with these choices, properties ( $\mathrm{I}_{m+1}$ ) and ( $\mathrm{II}_{m+1}$ ) hold.

Verification of Property $\left(\mathrm{I}_{m+1}\right)$. - We have to check that the polynomials $Q(P) \Phi_{\alpha}(P), Q(P) \psi_{\gamma}(P)$ are linearly independent in $\mathcal{G}(P)$ for $P \in \underline{Y}_{\nu}$. Since $Q(P)=1$ at the origin, it is enough to show that $\Phi_{\alpha}(P), \psi_{\gamma}(P)$ are linearly independent in $\mathcal{G}(P)$. If $\sum u_{\alpha} \Phi_{\alpha}(P)+$ $\sum w_{\gamma} \psi_{\gamma}(P)=0$ in $\mathcal{G}(P)$, this sum lies in $V_{P}^{\prime}$ by definition of $V_{P}^{\prime}$, so that there are constants $v_{\beta} \in \mathbb{R}$ so that $\sum u_{\alpha} \Phi_{\alpha}(P)+\sum w_{\gamma} \psi_{\gamma}(P)=\sum v_{\beta} h_{\beta}(P)$ (since the $h_{\beta}(P)$ span $V_{P}^{\prime}$ ). But this can only happen if $u_{\alpha}=v_{\beta}=w_{\gamma}=0$ since, as noted earlier, $\Phi_{\alpha}(P), h_{\beta}(P), \psi_{\gamma}(P)$ form a basis of $H^{D(m)}$ for $P \in Y_{\mu}^{\prime}=\underline{Y_{\nu}}$.

Verification of Property $\left(\mathrm{II}_{m+1}\right)$. - Recalling the choices made in (a), (b), (c) above, we can formulate property ( $\mathrm{II}_{m+1}$ ) as follows.
(11) Let $P_{0} \in Y_{\mu}^{\prime} \backslash \underline{Z}$. There exist $\underline{\delta}>0$ and $\underline{K}>0$ (depending only on $P_{0}$ and $m+1$ ) such that, whenever we are given:

$$
\begin{gathered}
P \in Y_{\mu}^{\prime} \text { with }\left\|P-P_{0}\right\|<\underline{\delta} ; \\
g_{1}, \ldots, g_{r} \in H^{D} \text { with }\left\|g_{j}\right\| \leq 1 ; \\
\tau \in(0,1) ; \text { and } f \in H^{D_{t}},
\end{gathered}
$$

we can find real numbers $\underline{A_{\alpha}}, \underline{E_{\gamma}}$ and a polynomial $\underline{f} \in H^{D(m+1)}$ such that

$$
\begin{array}{r}
(\mathrm{i}, m+1) \quad(Q(P))^{m+1} f \equiv \sum_{\alpha} \underline{A_{\alpha}} Q(P) \Phi_{\alpha}(P)+\sum_{\gamma} \underline{E_{\gamma}} Q(P) \psi_{\gamma}(P)+\tau^{m+1} \underline{f} \\
\bmod \left(p_{1}+\tau g_{1}, \ldots, p_{r}+\tau g_{r}\right)
\end{array}
$$

and

$$
\|\underline{f}\| \leq \underline{K}\|f\| .
$$

To verify this, recall that by inductive hypothesis (Technical Lemma for the value $m$ ), we have the following:

Given $P_{0} \in Y_{\mu}^{\prime} \backslash \underline{Z}$, there exist $\delta, K>0$ such that whenever we are given: $P \in Y_{\mu}^{\prime}$ with $\left\|P-P_{0}\right\|<\delta ; g_{1}, \ldots, g_{r} \in H^{D}$ with $\left\|g_{j}\right\| \leq 1 ; \tau \in(0,1)$; and $f \in H^{D_{t}}$; then there exist $A_{\alpha} \in \mathbb{R}$ and $F \in H^{D(m)}$ such that
(a) $(Q(P))^{m} f \equiv \sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)+\tau^{m} F \bmod \left(p_{1}+\tau g_{1}, \ldots, p_{r}+\tau g_{r}\right)$,
and
(b) $\quad\|F\| \leq K\|f\|$.

To prove (11), given $P$ and $f$, let $F$ be as in (a) above; we have $F \in H^{D(m)}$ and we can apply (10). We multiply (a) by $Q(P)$ and substitute into (10). This gives

$$
\begin{align*}
(Q(P))^{m+1} f \equiv & \sum_{\alpha} A_{\alpha} Q(P) \Phi_{\alpha}(P)  \tag{12}\\
& +\tau^{m}\left\{\sum_{\alpha} a_{\alpha} Q(P) \Phi_{\alpha}(P)+\sum_{\gamma} e_{\gamma} Q(P) \psi_{\gamma}(P)\right\} \\
& +\tau^{m+1} \underline{F} \bmod \left(p_{1}+\tau g_{1}, \ldots, p_{r}+\tau g_{r}\right)
\end{align*}
$$

where

$$
\begin{equation*}
\underline{F}=-\sum_{j=1}^{r}\left(\sum_{\beta} b_{\beta} g_{\beta}^{j}(P)\right) g_{j} \tag{13}
\end{equation*}
$$

Equation (12) is of the form (i, $m+1$ ) if we set

$$
\underline{A_{\alpha}}=A_{\alpha}+\tau^{m} a_{\alpha}, \underline{E_{\gamma}}=\tau^{m} e_{\gamma}, \underline{f}=\underline{F} .
$$

Thus, to complete the proof of the technical lemma, it remains to prove the estimate (ii, $m+1$ ) for $\underline{F}=\underline{f}$. To do this, we have only to prove the following:
(14) Given $P_{0} \in Y_{\mu}^{\prime} \backslash \underline{Z}$, there exist $\bar{\delta}, \bar{K}>0$ (depending only on $P_{0}$ ) such that: if $P \in Y_{\mu}^{\prime}$ and $\left\|P-P_{0}\right\|<\bar{\delta}$; if

$$
\begin{equation*}
F=\sum a_{\alpha} \Phi_{\alpha}(P)+\sum b_{\beta} h_{\beta}(P)+\sum_{\gamma} e_{\gamma} \psi_{\gamma}(P) \tag{15}
\end{equation*}
$$

and we are given $g_{1}, \ldots, g_{r} \in H^{D}$ with $\left\|g_{j}\right\| \leq 1$; then we have

$$
\begin{equation*}
\left\|\sum_{j=1}^{r} \sum_{\beta} b_{\beta} g_{\beta}^{j}(P) g_{j}\right\| \leq \bar{K}\|F\| \tag{16}
\end{equation*}
$$

Now, by the choice of $\underline{Z}$, the $\operatorname{map} P \mapsto\left(\Phi_{\alpha}(P), h_{\beta}(P), \psi_{\gamma}(P)\right)$ is continuous on $Y_{\mu}^{\prime}$ at $P_{0}$ (since its restriction to $Y_{\mu}^{\prime} \backslash \underline{Z}$ is continuous and $\underline{Z}$
is closed). Moreover, the triple in brackets is a basis of $H^{D(m)}$ close to the basis $\left(\Phi_{\alpha}\left(P_{0}\right), h_{\beta}\left(P_{0}\right), \psi_{\gamma}\left(P_{0}\right)\right)$ if $P$ is close to $P_{0}$. Hence, if $\left\|P-P_{0}\right\| \leq \tilde{\delta}$, $P \in Y_{\mu}^{\prime}\left(\tilde{\delta}\right.$ depending only on $P_{0}$ and $\left.D(m)\right)$, there is $\tilde{K}_{1}>0$ (depending only on $P_{0}$ and $D(m)$ ) so that

$$
\begin{equation*}
\left|a_{\alpha}\right|,\left|b_{\beta}\right|,\left|e_{\gamma}\right| \leq \tilde{K}_{1}\|F\| \tag{17}
\end{equation*}
$$

whenever (15) holds.
In addition, the maps $g_{\beta}^{j}(P)$ are continuous on $Y_{\mu}^{\prime}$ at $P_{0}$ (again because $\underline{Z}$ is closed). Thus, there exist $\tilde{\delta}_{2}, \tilde{K}_{2}>0$ (depending only on $P_{0}$ and $D(m)$ ) such that

$$
\begin{equation*}
\left\|g_{\beta}^{j}(P)\right\| \leq \tilde{K}_{2} \text { if }\left\|P-P_{0}\right\|<\tilde{\delta}_{2}, P \in Y_{\mu}^{\prime} \tag{18}
\end{equation*}
$$

Since $\left\|g_{j}\right\| \leq 1$ by assumption, (17) and (18) show that whenever (15) holds, we have, for $\left\|P-P_{0}\right\|<\bar{\delta}=\min \left(\tilde{\delta}_{1}, \tilde{\delta}_{2}\right), P \in Y_{\mu}^{\prime}$,

$$
\left\|\sum_{j} \sum_{\beta} b_{\beta} g_{\beta}^{j}(P) g_{j}\right\| \leq \bar{K}\|F\|
$$

This completes the verification of property $\left(\mathrm{II}_{m+1}\right)$.
The induction step being complete, Technical Lemma 8.4 is proved.
Lemma 8.5. - Let $Y$ be a semi-algebraic subset of $W$ with $\operatorname{dim} Y \leq t$. Then, there exists a closed semi-algebraic subset $Z \subset W$ with $\operatorname{dim} Z \leq t-1$ (so that $Z=\emptyset$ if $t=0$ ) with the following property:

If $P_{0} \in Y \backslash Z$ and $\rho>0$, there exist constants $\delta, K>0$ (depending only on $P_{0}$ and $\rho$ ) such that: given $P \in Y$ with $\left\|P-P_{0}\right\|<\delta$ and $f \in H^{D}$, there exists $F \in H^{D}$ with $F=f$ in $\mathcal{G}(P)$ and $\|F\| \leq K \sup _{V_{\rho}(P)}|f|$.

Proof. - If $P \in Y$, denote by $V_{P}^{\prime \prime}$ the orthogonal complement in $H^{D}$ of the space of polynomials in $H^{D}$ which are 0 as germs in $\mathcal{G}(P)$. By Lemma 7.2 and Lemmas $6.14,6.13$, we can find a semi-algebraic partition $\left\{Y_{\nu}\right\}$ of $Y$ and, on each $Y_{\nu}$, semi-algebraic maps $\Phi_{\alpha}: Y_{\nu} \rightarrow H^{D}$ such that $\left\{\Phi_{\alpha}(P)\right\}$ is a basis of $V_{P}^{\prime \prime}$ for any $P \in Y_{\nu}$.

Let $Z$ be a closed semi-algebraic set in $W$ such that $\operatorname{dim} Z \leq t-1$, $Z \supset \bar{Y}_{\nu} \backslash Y_{\nu}$ for each $\nu\left(\bar{Y}_{\nu}\right.$ being the closure of $\left.Y_{\nu}\right)$, and such that for any $\nu$ and any $\alpha, \Phi_{\alpha} \mid Y_{\nu} \backslash Z$ is continuous; there is such a $Z$ by Lemmas 6.11, 6.7. Since $Z \supset \bar{Y}_{\nu} \backslash Y_{\nu}$ for all $\nu$, if $P_{0} \in Y_{\nu} \backslash Z$, then $P_{0} \notin \bar{Y}_{\mu}$ for $\mu \neq \nu$. Hence, given $P_{0} \in Y_{\nu} \backslash Z$, there is $\delta_{1}>0$ such that if $P \in Y$ and $\left\|P-P_{0}\right\|<\delta_{1}$, then $P \in Y_{\nu}$.

Now, if $P \in Y_{\nu}$ and $f \in H^{D}$, we can find real numbers $A_{\alpha}$ such that $f=\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)$ as germs in $\mathcal{G}(P)$. Since, for $0<\rho \leq \rho_{1}$ (see $\S 1,(4)), V_{\rho}(P)$ is contained in the connected component through 0 of $Z(P) \cap\left\{x \in \mathbb{R}^{n} \left\lvert\, \operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}\right) \neq 0\right.\right\}$ (see remark after Lemma 2.1), we have

$$
\begin{equation*}
f=\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P) \text { on } V_{\rho}(P) \tag{19}
\end{equation*}
$$

We now make the following remarks:
(a) If $x_{0} \in V_{\rho / 2}\left(P_{0}\right)$ and $\epsilon>0$ are given, there is $\delta_{2}>0$ such that if $P \in W$ and $\left\|P-P_{0}\right\|<\delta_{2}$, then, there exists $x \in V_{\rho}(P)$ such that $\left|x-x_{0}\right|<\epsilon$.

This follows from the real version of Lemma 2.1.
(b) Given $P_{0} \in Y_{\nu} \backslash Z$ and $\epsilon>0$, there is $\delta_{3}>0$ such that if $P \in Y_{\nu}$ and $\left\|P-P_{0}\right\|<\delta_{3}$, then $\left\|\Phi_{\alpha}(P)-\Phi_{\alpha}\left(P_{0}\right)\right\|<\epsilon$ for any $\alpha$.

This is simply the statement that the $\Phi_{\alpha}$ are continuous on $Y_{\nu}$ at $P_{0} \notin Z$.
(c) $\left\{\Phi_{\alpha}(P)\right\}$ are linearly independent functions on $V_{\rho / 2}\left(P_{0}\right)$. (They form a basis of $V_{P_{0}}^{\prime \prime}$ and so are linearly independent in $\left.\mathcal{G}(P)\right)$.

Because of (c), we can apply Lemma 8.1. Thus, there exist finitely many points $x_{\beta} \in V_{\rho / 2}\left(P_{0}\right), \epsilon>0$ and $K$ (depending only on the $\Phi_{\alpha}\left(P_{0}\right)$, thus only on $P_{0}$ ) such that if $\tilde{x}_{\beta} \in \mathbb{R}^{n},\left|\tilde{x}_{\beta}-x_{\beta}\right|<\epsilon$ and $\varphi_{\alpha} \in H^{D}$, $\left\|\varphi_{\alpha}-\Phi_{\alpha}\left(P_{0}\right)\right\|<\epsilon$, then

$$
\sum_{\alpha}\left|A_{\alpha}\right| \leq K \max _{\beta}\left|\sum_{\alpha} A_{\alpha} \varphi_{\alpha}\left(\tilde{x}_{\beta}\right)\right|, A_{\alpha} \in \mathbb{R}
$$

We can apply this with $\varphi_{\alpha}=\Phi_{\alpha}(P), \tilde{x}_{\beta} \in V_{\rho}(P)$ if $\left\|P-P_{0}\right\|<\min \left(\delta_{2}, \delta_{3}\right)$ because of (a) and (b) above. This gives

$$
\begin{aligned}
\left\|\sum A_{\alpha} \Phi_{\alpha}(P)\right\| & \leq \max _{\alpha}\left(\left\|\Phi_{\alpha}\left(P_{0}\right)\right\|+\epsilon\right) \sum_{\alpha}\left|A_{\alpha}\right| \\
& \leq K^{\prime} \max _{\beta}\left|\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)\left(\tilde{x}_{\beta}\right)\right| \\
& \leq K^{\prime} \sup _{V_{\rho}(P)}\left|\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)\right|
\end{aligned}
$$

Taking $F=\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)$ and using (19), we obtain the lemma.

Remark 8.6. - We have used Theorem 5.6 in Lemma 8.3 and thus in the statement and proof of the Technical Lemma 8.4. However, to prove the Main Lemma, a weaker version of the Technical Lemma would be sufficient; this weaker version can be proved using only Theorem 5.2 (not 5.6). This version of the Technical Lemma is the following.

Given $W^{t}, C_{t}, D_{t}$ and $m \geq 0$ as in 8.4, there exist constants $C_{0}(m), D_{0}(m)$ (depending only on these data and the constants in $W$ ) and a semi-algebraic $\operatorname{map} \stackrel{\vee}{q}$ : $W^{t} \rightarrow S\left(C_{0}(m), D_{0}(m)\right.$ ) (so that $\operatorname{deg} \stackrel{\vee}{q}(P) \leq$ $D_{0}(m), \stackrel{\vee}{q}(P)(0)=1,\|\stackrel{\vee}{q}(P)\| \leq C_{0}(m)$ for $\left.P \in W^{t}\right)$, with the property that the statement of the Technical Lemma holds if, in ( $\mathrm{II}_{m}$ ), we replace $(Q(P))^{m}$ by $\stackrel{\vee}{q}(P)$.

The sharper form in the text is the exact analogue of the argument given in [FN]. Moreover, Theorem 5.5 and Theorem 5.6 on which it is based, are clearly of interest in themselves.

## 9. The induction step in the proof of the Main Lemma.

We are assuming that the Main Lemma has been proved for the value $t$ of the parameter; in this section we carry out the induction step proving it for the value $t-1$ under this assumption, thus completing the proof of the Main Lemma 8.2.

If we recall the definition of the norms $N, N_{*}$ defined in $\S 7$, we can reformulate the Main Lemma for $t$ as follows.

There exists a closed semi-algebraic set $W^{t} \subset W$ of dimension $\leq t$, and constants $C_{t}, D_{t} \geq 1$ such that

Given $P_{0} \in W \backslash W^{t}$ and $\rho>0, \rho \leq \rho_{1}$, there exist $\delta>0, K_{*}>0$ so that

$$
N_{*}\left(P, D_{t}, C_{t}, \rho\right) \leq K_{*} \text { whenever }\left\|P-P_{0}\right\|<\delta, P \in W \backslash W^{t}
$$

If we note that $W^{t}$ is closed in $W$ and that, for fixed values of the other parameters, $N_{*}$ is a decreasing function of $\rho$, we obtain the following:
(1) With $C_{t}, D_{t}$ as above, the function

$$
(P, \rho) \mapsto N_{*}\left(P, D_{t}, C_{t}, \rho\right)
$$

is locally bounded on the set $\left(W \backslash W^{t}\right) \times\left(0, \rho_{1}\right]$.

Since $N_{*}$ is semi-algebraic where it is finite (Lemma 7.4, (b)) and $W \times\left[0, \rho_{1}\right]$ is compact, we can apply Lemma 6.15 to $N_{*}$ and the pair $\left(W \times\left[0, \rho_{1}\right],\left(W \backslash W^{t}\right) \times\left(0, \rho_{1}\right]\right)$.

Note that the distance of $(P, \rho) \in\left(W \backslash W^{t}\right) \times\left(0, \rho_{1}\right]$ from $\left(W \times\left[0, \rho_{1}\right]\right) \backslash\left(\left(W \backslash W^{t}\right) \times\left(0, \rho_{1}\right]\right)=W^{t} \times\left[0, \rho_{1}\right] \cup W \times\{0\}$ is at least $\min \left(\operatorname{dist}\left(P, W^{t}\right), \rho\right) \geq \rho \operatorname{dist}\left(P, W^{t}\right)$ (since $\rho_{1} \leq 1$; see $\S 1$, (4)). From this, and Lemma 6.15, we obtain:

Lemma 9.1 (Quantitative Form of Main Lemma). - Let $W^{t}, C_{t}$, $D_{t}$ be as in the Main Lemma 8.2. There exist constants $m_{t}, K_{t}>0$ (depending only on these data) such that the following holds.

Given $P \in W \backslash W^{t}$ and $\rho>0,0<\rho \leq \rho_{1}$, we can find $q \in H^{D_{t}}$ with $q(0)=1,\|q\| \leq C_{t}$ such that, for any $f \in H^{D}$, there exists $F \in H^{D_{t}}$ for which
(i) $F=q f$ on $V_{\rho}(P)$
and
(ii) $\|F\| \leq K_{t} \cdot \rho^{-m_{t}}\left(\operatorname{dist}\left(P, W^{t}\right)\right)^{-m_{t}} \sup _{V_{\rho}(Z)}|f|$.

We now apply the technical lemma with $m=m_{t}$, where $m_{t}$ is as in Lemma 9.1. This gives us: a semi-algebraic partition $\left\{Y_{\nu}\right\}$ of $W^{t}$; a closed semi-algebraic set $Z \subset W^{t}$ with $\operatorname{dim} Z \leq t-1$; a constant $D\left(m_{t}\right) \geq 1$; and, for each $\nu$, semi-algebraic maps $\Phi_{\alpha}: Y_{\nu} \rightarrow H^{D\left(m_{t}\right)}$, having the properties $\left(\mathrm{I}_{m_{t}}\right)$ and ( $\mathrm{II}_{m_{t}}$ ) of Lemma $8.4\left(Q: W \rightarrow S\left(C_{*}, D_{*}\right)\right.$ is the semi-algebraic map constructed in Lemma 8.3).

Let $Z^{\prime}$ be a closed semi-algebraic set of dimension $\leq t-1$ such that, for any $\nu, \alpha, Q, \Phi_{\alpha}$ restricted to $Y_{\nu} \backslash Z^{\prime}$ are continuous; there is such a $Z^{\prime}$ by Lemma 6.11.

Let $Z^{\prime \prime} \subset W$ be a closed semi-algebraic set of dimension $\leq t-1$ such that, for any $\nu$, if $P_{0} \in Y_{\nu} \backslash Z^{\prime \prime}$ and $0<\rho \leq \rho_{1}$, if $P \in Y_{\nu}$ is close to $P_{0}$, and if $f \in H^{D}$, then there exists $F \in H^{D}$ with $F=f$ on $V_{\rho}(P)$ and $\|F\| \leq K \sup _{V_{\rho}(P)}|f|, K$ being a constant depending only on $P_{0}$ and $\rho$. The existence of $Z^{\prime \prime}$ is guaranteed by Lemma 8.5 (and Lemma 6.6).

We define $W^{t-1}$ by:
(2) $W^{t-1}=$ closure in $W$ of $Z \cup Z^{\prime} \cup Z^{\prime \prime} \cup \bigcup_{\nu}\left(\bar{Y}_{\nu} \backslash Y_{\nu}\right), \bar{Y}_{\nu}$ being the closure of $Y_{\nu}$.

Since $Z, Z^{\prime}, Z^{\prime \prime}$ have dimension $\leq t-1$ by construction, and $\operatorname{dim} Y_{\nu} \leq$ $\operatorname{dim} W^{t} \leq t$, we have, by Lemmas 6.6, 6.7,

$$
\operatorname{dim} W^{t-1} \leq t-1
$$

By construction, we have the following properties:
(3) If $P_{0} \in Y_{\nu} \backslash W^{t-1}$ and $0<\rho \leq \rho_{1}(\leq 1)$, there exist $\delta_{0}, K_{0}>0$ depending only on $P_{0}$ and $\rho$ such that the following assertions hold:
(3.i) $Q, \Phi_{\alpha}$ are continuous on $Y_{\nu}$ at $P_{0}$.
(3.ii) If $P \in W^{t}$ and $\left\|P-P_{0}\right\|<\delta_{0}$, then $P \in Y_{\nu}$.
(3.iii) If $P \in Y_{\nu}$ and $\left\|P-P_{0}\right\|<\delta_{0}$, then given $f \in H^{D}$, there exists $F \in H^{D}$ with $F=f$ on $V_{\rho}(P)$ and $\|F\| \leq K_{0} \sup |f|$.

$$
V_{\rho}(P)
$$

To complete the induction step, we have to show that there exist constants $C_{t-1}, D_{t-1} \geq 1$ (depending only on the constants defining $W$ ) having the following property:
(4) Let $P_{0} \in W \backslash W^{t-1}$ and let $0<\rho \leq \rho_{1}$. There exist constants $\underline{\delta}$, $\underline{K}>0$ (depending only on $P_{0}$ and $\rho$ ) such that whenever $\hat{P} \in W$ and $\left\|\hat{P}-P_{0}\right\|<\underline{\delta}$, we can find $\hat{q} \in H^{D_{t-1}}$ with $\hat{q}(0)=1,\|\hat{q}\| \leq C_{t-1}$ with the property that for any $f \in H^{D}$, there is $F \in H^{D_{t-1}}$ for which $F=\hat{q} f$ on $V_{\rho}(\hat{P})$ and $\|F\| \leq \underline{K} \sup _{V_{\rho}(\hat{P})}|f|$.

If $P_{0} \in W \backslash W^{t}$, this follows from our inductive hypothesis that the Main Lemma is true for the value $t$.

If $P_{0} \in W^{t} \backslash W^{t-1}$ and $\hat{P} \in W^{t}$, then, if $\nu$ is such that $P_{0} \in Y_{\nu}$ and $\left\|\hat{P}-P_{0}\right\|$ is small, (3.ii) implies that $\hat{P} \in Y_{\nu}$ (same $\nu$ ) and (3.iii) then shows that (4) holds with $\hat{q} \equiv 1$.

Thus, to prove (4) and hence complete the induction step, we have only to prove

Lemma 9.2. - Let $\left\{Y_{\nu}\right\}, W^{t-1}$ be as above. There exist constants $C_{t-1}>0, D_{t-1} \geq 1$ (depending only on the constants defining $W$ ) for which we have the following:

Given $P_{0} \in Y_{\nu} \backslash W^{t-1}$ and $0<\rho \leq \rho_{1}$, there exist $\underline{\delta}, \underline{K}>0$ (depending only on $P_{0}$ and $\rho$ ) such that if $\hat{P} \in W \backslash W^{t}$ and $\left\|\hat{P}-P_{0}\right\|<\underline{\delta}$, then there is $\underline{q} \in H^{D_{t-1}}$ with the following properties:
(a) $\underline{q}(0)=1,\|\underline{q}\| \leq C_{t-1} ;$
(b) For any $f \in H^{D}$, we can find $F \in H^{D_{t-1}}$ so that $F=\underline{q} f$ on $V_{\rho}(\hat{P})$ and $\|F\| \leq \underline{K} \sup _{V_{\rho}(\hat{P})}|f|$.

$$
V_{\rho}(\hat{P})
$$

Proof. - Let $P_{0} \in Y_{\nu} \backslash W^{t-1}$ and $\rho>0$ be given. Let $\hat{P} \in W \backslash W^{t}$ and suppose that $\left\|\hat{P}-P_{0}\right\|<\underline{\delta}<1 / 2$; we shall list the conditions that $\underline{\delta}$ must satisfy for the lemma to hold.

Since $W^{t}$ is closed, we have $\tau=\operatorname{dist}\left(\hat{P}, W^{t}\right)>0$; further $\tau<\underline{\delta}$ since $P_{0} \in Y_{\nu} \subset W^{t}$. We choose $P \in W^{t}$ such that $\|\hat{P}-P\|=\tau$; we then have

$$
\left\|P-P_{0}\right\| \leq\|P-\hat{P}\|+\left\|\hat{P}-P_{0}\right\| \leq \tau+\underline{\delta}<2 \underline{\delta} .
$$

Hence:
(5) If $\delta_{1}>0$ is the constant in (3.i), (3.ii), (3.iii) above, and if $2 \underline{\delta} \leq \delta_{1}$, then $P \in Y_{\nu}$.

Write $\hat{P}=P+\tau\left(g_{1}, \ldots, g_{r}\right)$ with $g_{j} \in H^{D}$. Then $\sum_{\hat{P}}\left\|g_{j}\right\|^{2}=1$. We apply Lemma 9.1 (quantitative form of Main Lemma) to $\hat{P} \in W \backslash W^{t}$. This gives:
(6) Given $\rho>0$, there is $q \in H^{D_{t}}$ with $q(0)=1,\|q\| \leq C_{t}$ such that for any $f \in H^{D}$, we can find $F_{0} \in H^{D_{t}}$ for which $F_{0}=q f$ on $V_{\rho}(\hat{P})$ and $\left\|F_{0}\right\| \leq K_{t} \tau^{-m_{t}} \rho^{-m_{t}} \sup _{V_{\rho}(\hat{P})}|f|$.

We wish to apply the Technical Lemma 8.4 to $F_{0}$ with $m=m_{t}$. To do this, we must make sure that
(7) $2 \underline{\delta} \leq \delta, \delta$ being the constant in Property $\left(\mathrm{II}_{m_{t}}\right)$ of 8.4.

If (5) and (7) hold, we have $\left\|P-P_{0}\right\|<\delta, P_{0} \in Y_{\nu} \backslash W^{t-1}, P \in Y_{\nu}$, $g_{1}, \ldots, g_{r} \in H^{D},\left\|g_{j}\right\| \leq 1,0<\tau<1, F_{0} \in H^{D_{t}}$, so that the conditions imposed in Lemma 8.4 are verified. Hence, if the $\Phi_{\alpha}$ are as in Lemma 8.4 with $m=m_{t}$, we have:
(8) There exist real numbers $A_{\alpha}$ such that

$$
(Q(P))^{m_{t}} F_{0} \equiv \sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)+\tau^{m_{t}} \tilde{F} \bmod \left(\hat{p}_{1}, \ldots, \hat{p}_{r}\right), \hat{P}=\left(\hat{p}_{1}, \ldots, \hat{p}_{r}\right)
$$

and

$$
\|\tilde{F}\| \leq \tilde{K}\left\|F_{0}\right\|
$$

note that $\hat{p}_{j}=p_{j}+\tau g_{j}$ by definition of the $g_{j}$. Here $\tilde{K}$ is a constant depending only on $\tau, P_{0}$ and $\rho$, and $\Phi_{\alpha}(P), \tilde{F} \in H^{D\left(m_{t}\right)}$.

From (6) and (8), we obtain:
(9) Set $\underline{q}=q \cdot Q^{m_{t}}(P)$ (with $q$ as in (6) and $Q$ as in 8.3). Then, if $f \in H^{D}$, we have

$$
\begin{align*}
& \underline{q} \cdot f=\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)+\tau^{m_{t}} \tilde{F} \text { on } V_{\rho}(\hat{P})  \tag{9.i}\\
& \left\|\tau^{m_{t}} \tilde{F}\right\| \leq \tilde{K} \cdot K_{t} \rho^{-m_{t}} \sup _{V_{\rho}(\hat{P})}|f| \tag{9.ii}
\end{align*}
$$

and, because of Lemma 8.3 and (6) above, we also have

$$
\begin{align*}
& \underline{q}(0)=1, \underline{q} \in H^{D_{t-1}^{\prime}} \text { with } D_{t-1}^{\prime}=D_{t}+m_{t} D_{*}  \tag{9.iii}\\
& \|\underline{q}\| \leq \tilde{C} \cdot C_{t} \cdot C_{*}^{m_{t}}=C_{t-1}
\end{align*}
$$

where $C_{*}$ is as in Lemma 8.3, $C_{t}$ is as in (6) above, and $\tilde{C}$ is a constant depending only on $n$ and the degrees of $q, Q(P)^{m_{t}}$, i.e. only on $n, D_{t}, D_{*}, m_{t}$.

Since, by (9.iii), $\sup _{V_{\rho}(\hat{P})}|\underline{q} \cdot f| \leq K_{t-1}^{\prime} \sup _{V_{\rho}(\hat{P})}|f|$, where $K_{t-1}^{\prime}$ depends only on $P_{0}, \rho$ and $t-1,(9 . i)$ and (9.ii) imply that

$$
\begin{equation*}
\sup _{V_{\rho}(\hat{P})}\left|\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)\right| \leq K^{\prime \prime} \sup _{V_{\rho}(\hat{P})}|f|, \tag{10}
\end{equation*}
$$

where $K^{\prime \prime}$ depends only on $P_{0}, \rho$ and $t$.
Since the functions $\Phi_{\alpha}\left(P_{0}\right) \mid V_{\rho / 2}\left(P_{0}\right)$ are linearly independent, we can apply Lemma 8.1 to obtain the following:
(11) There are finitely many points $x_{\beta} \in V_{\rho / 2}\left(P_{0}\right)$, and constants $\epsilon>0$, $K_{*}>0$ such that if $\tilde{x}_{\beta} \in \mathbb{R}^{n},\left|\tilde{x}_{\beta}-x_{\beta}\right|<\epsilon$, and $\varphi_{\alpha} \in H^{D\left(m_{t}\right)}$, $\left\|\varphi_{\alpha}-\Phi_{\alpha}\left(P_{0}\right)\right\|<\epsilon$, then, for any real numbers $A_{\alpha}$, we have

$$
\left\|\sum_{\alpha} A_{\alpha} \varphi_{\alpha}\right\| \leq K_{*} \max _{\beta}\left|\sum_{\alpha} A_{\alpha} \varphi_{\alpha}\left(\tilde{x}_{\beta}\right)\right| .
$$

Here, $\epsilon, K_{*}$ depend only on $P_{0}, \rho$ and $m_{t}$. With $\epsilon$ as in (11), there is $\delta_{2}>0$ depending only on $P_{0}, \rho, t$ such that if $\left\|\hat{P}-P_{0}\right\|<\delta_{2}$, then, for any $x_{0} \in V_{\rho / 2}\left(P_{0}\right)$, there exists $x \in V_{\rho}(\hat{P})$ such that $\left|x-x_{0}\right|<\epsilon$.

We assume that

$$
\begin{equation*}
\underline{\delta}<\delta_{2} . \tag{12}
\end{equation*}
$$

Now note that the maps $P^{\prime} \mapsto \Phi_{\alpha}\left(P^{\prime}\right)$ are continuous on $Y_{\nu}$ at the point $P^{\prime}=P_{0}$ (since $P_{0} \in Y_{\nu} \backslash W^{t-1}$ and $W^{t-1}$ contains the points of discontinuity of the $\Phi_{\alpha}$ on $Y_{\nu}$ by construction). This gives
(13) There is $\delta_{3}>0$ depending only on $P_{0}$ and $\rho$ such that if $\underline{\delta}<\delta_{3}$, then $\left\|\Phi_{\alpha}\left(P^{\prime}\right)-\Phi_{\alpha}\left(P_{0}\right)\right\|<\epsilon$ for $\left\|P^{\prime}-P_{0}\right\|<2 \underline{\delta}, P^{\prime} \in Y_{\nu}$.

We now pick $\underline{\delta}$ to satisfy the conditions (5), (7), (12) and (13). We may then apply (11) with $\varphi_{\alpha}=\Phi_{\alpha}(P)$ and $\tilde{x}_{\beta} \in V_{\rho}(\hat{P}),\left|\tilde{x}_{\beta}-x_{\beta}\right|<\epsilon$, to obtain

$$
\begin{align*}
\left\|\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)\right\| & \leq K_{*} \max _{\beta}\left|\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)\left(\tilde{x}_{\beta}\right)\right|  \tag{14}\\
& \leq K_{*} \sup _{V_{\rho}(\hat{P})}\left|\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)\right| \\
& \leq K^{\prime \prime} \cdot K_{*} \sup _{V_{\rho}(\hat{P})}|f| \text { by }(10)
\end{align*}
$$

Thus, if we define

$$
\begin{equation*}
F=\sum_{\alpha} A_{\alpha} \Phi_{\alpha}(P)+\tau^{m_{t}} \tilde{F}, \text { with } \tilde{F} \text { as in (8) } \tag{15}
\end{equation*}
$$

we obtain, from (6), (9) and (14),

$$
\begin{align*}
\underline{q} f & =F \text { on } V_{\rho}(\hat{P}), F \in H^{D\left(m_{t}\right)}  \tag{16}\\
\|F\| & \leq \sup _{V_{\rho}(\hat{P})}|f|, \text { where } \underline{K}=K^{\prime \prime} K_{*}+\tilde{K} \cdot K_{t} \cdot \rho^{-m_{t}}
\end{align*}
$$

If we set $D_{t-1}=\max \left(D\left(m_{t}\right), D_{t}+m_{t} D_{*}\right.$ ), we also have, by (9.iii), $\underline{q} \in H^{D_{t-1}}, \underline{q}(0)=1,\|\underline{q}\| \leq C_{t-1}, F \in H^{D_{t-1}}$.

Thus, if $D_{t-1}=\max \left(D\left(m_{t}\right), D_{t}+m_{t} D_{*}\right)$, and $C_{t-1}=\tilde{C} \cdot C_{t} \cdot C_{*}^{m_{t}}$ (as in (9.iii)), then choosing $\underline{\delta}>0$ to satisfy (5), (7), (12), (13) and $\underline{K}>0$ as in (16), we have shown that if $\hat{P} \in W \backslash W^{t},\left\|\hat{P}-P_{0}\right\|<\underline{\delta}$ and $\underline{q}$ is defined as in (9) (where $P \in Y_{\nu}$ satisfies $\|\hat{P}-P\|=\tau=\left(\hat{P}, W^{t}\right)$ ), then the conclusions of Lemma 9.2 hold.

This proves the lemma, and with it the Main Lemma 8.2 as remarked earlier.

## 10. Proof of the Extension Theorem.

Proof of Part 1. - We take $t=-1$ in the Main Lemma 8.2, so that $W^{t}=\emptyset$. Using the definition of the norm $N_{*}(8.3,(\mathrm{~b}))$, we can formulate the conclusion of the Main Lemma as follows.
(1) There exist constants $C_{-1}>0, D_{-1} \geq 1$ depending only on $n, r, D, C_{1}$ such that if $\rho>0$ and $P_{0} \in W$, then there exist $\delta_{0}>0, K_{0}>0$ depending only on $P_{0}$ and $\rho$ such that

$$
N_{*}\left(P, D_{-1}, C_{-1}, \rho\right) \leq K_{0} \text { for } P \in W,\left\|P-P_{0}\right\|<\delta_{0}
$$

Since $W$ is compact, for each $\rho>0$, there is a constant $K_{*}$ depending only on $\rho$ and $W$ such that

$$
N_{*}\left(P, D_{-1}, C_{-1}, \rho\right) \leq K_{*} \text { for } P \in W
$$

Since $N_{*}$ is a decreasing function of $\rho\left(0<\rho \leq \rho_{1}\right)$, we conclude, in particular, that
(2) The function $(P, \rho) \mapsto N_{*}\left(P, D_{-1}, C_{-1}, \rho\right)$ is locally bounded on $W \times\left(0, \rho_{1}\right]$.

Because of Lemma 7.4, (b), this function is semi-algebraic, and we can apply Lemma 6.15 with $S=W \times\left[0, \rho_{1}\right], E=W \times\left(0, \rho_{1}\right]$. This gives:
(3) There exist constants $m>0, \bar{C}>0$ (depending only on $\rho_{1}, D_{-1}, C_{-1}$ and $W$ ) such that

$$
N_{*}\left(P, D_{-1}, C_{-1}, \rho\right) \leq \bar{C} \rho^{-m} \text { for all } P \in W, 0<\rho \leq \rho_{1}
$$

If we take $D^{\prime}=D_{-1}$ and $C^{\prime}=\max \left(C_{-1}, \bar{C}\right)$, this is simply a restatement of Part 1 of the Extension Theorem. [Note that if $P \in W$ and a polynomial $g$ is 0 as a germ in $\mathcal{G}(P)$, then $g=0$ on $V_{\rho_{1}}(P)$.]

Proof of Part 2. - Choose $\rho_{0}>0$ such that $2 \rho_{0} \leq \rho_{1}(\leq 1)$ and such that if $q \in H^{D^{\prime}}, q(0)=1,\|q\| \leq C^{\prime}$, then

$$
\frac{1}{2}<q<2 \text { on } Q_{2 \rho_{0}}
$$

By Part 1 of the Extension Theorem, there exists $q \in H^{D^{\prime}}, q(0)=1$, $\|q\| \leq C^{\prime}$ (hence $\frac{1}{2}<q<2$ on $Q_{2 \rho_{0}}$ ) such that for any $f \in H^{D}$, we can find $F \in H^{D^{\prime}}$ with

$$
f=F / q \text { on } V_{2 \rho_{0}}(P)
$$

and

$$
\|F\| \leq C^{\prime} \rho_{0}^{-m} \sup _{V_{\rho_{0}}(P)}|f|
$$

Since clearly $\sup _{Q_{2 \rho_{0}}}|F| \leq C\left(D^{\prime}, n\right)\|F\|$, we obtain Part 2 of the Extension Theorem for the value $\rho=\rho_{0}$.

If $0<\rho \leq \rho_{0}$, given $P=\left(p_{1}, \ldots, p_{r}\right) \in W$, consider the polynomials $\tilde{p}_{j}(x)=\frac{\rho_{0}}{\rho} p_{j}\left(\frac{\rho}{\rho_{0}} x\right)$. Since $p_{j}(0)=0$, we have $\left\|\tilde{p}_{j}\right\| \leq\left\|p_{j}\right\|$. Moreover,
$\frac{\partial \tilde{p}_{j}}{\partial x_{k}}(0)=\frac{\partial p_{j}}{\partial x_{k}}(0)$. Thus we can apply Part 2 for the value $\rho_{0}$ to $\tilde{f}(x)=$ $f\left(\frac{\rho}{\rho_{0}} x\right), f \in H^{D}$, to write

$$
\begin{gathered}
\tilde{f}=\tilde{F} / \tilde{q} \text { on } \quad V_{2 \rho_{0}}(\tilde{P}), \tilde{P}=\left(\tilde{p}_{1}, \ldots, \tilde{p}_{r}\right), \\
\frac{1}{2}<\tilde{q}<2 \text { on } Q_{2 \rho_{0}}, \sup _{Q_{2 \rho_{0}}}|\tilde{F}| \leq C^{\prime \prime} \sup _{V_{\rho_{0}}(\tilde{P})}|\tilde{f}| .
\end{gathered}
$$

Writing $\tilde{F}(x)=F\left(\frac{\rho x}{\rho_{0}}\right), \tilde{q}=q\left(\frac{\rho x}{\rho_{0}}\right)$ gives us Part 2 for the value $\rho$, $0<\rho \leq \rho_{0}$.

The proof of the Extension Theorem is complete.

## 11. Polynomial behaviour of algebraic functions.

In this section, we shall use the extension theorem to prove Theorem 1 stated in the Introduction.

Given $n \geq 2,1 \leq r \leq n-1, D \geq 1$ and $C_{1}>0$, we consider our basic space $W$ defined by these constants.

If $\theta>0$, and $s \geq 1$ is an integer, we denote by $B_{s}(\theta)$ the ball in $\mathbb{R}^{s}: B_{s}(\theta)=\left\{\left(x_{1}, \ldots, x_{s}\right) \in \mathbb{R}^{s} \mid \sum x_{j}^{2}<\theta^{2}\right\}$ (see $\left.\S 1\right)$.

By Lemma 2.1, there exist arbitrarily small constants $\theta_{1}, \theta_{2}$ with $0<\theta_{1}, \theta_{2}<1$ depending only on $n, r, D, C_{1}$ such that if $P \in W$ and $\pi=\pi_{P}$ is the restriction to $Z(P) \cap B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right)$ of the projection $\left(x_{1}, \ldots, x_{n}\right) \mapsto\left(x_{r+1}, \ldots, x_{n}\right)\left(Z(P)\right.$ is, as usual, the set $\left\{x \in \mathbb{R}^{n} \mid P(x)=\right.$ $0\}$ ), then we have:
(1) $\pi$ is a real analytic isomorphism of $Z(P) \cap B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right)$ onto $B_{n-r}\left(\theta_{2}\right)$. Moreover, there is a real analytic map $\varphi=\left(\varphi_{1}, \ldots, \varphi_{r}\right)$ : $B_{n-r}\left(\theta_{2}\right) \rightarrow B_{r}\left(\theta_{1}\right)$ such that

$$
Z(P) \cap B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right)=\left\{\left(\varphi_{1}(y), \ldots, \varphi_{r}(y), y\right) \mid y \in B\left(\theta_{2}\right)\right\}
$$

thus $\pi^{-1}$ is the map $y \mapsto(\varphi(y), y), y \in B\left(\theta_{2}\right)$.
If $\theta_{1}, \theta_{2}$ are chosen sufficiently small (depending on $n, D, r, C_{1}$ ), we have
(2) There is a constant $C_{2}>0$ (depending only on $\left.n, D, r, C_{1}\right)$ such that if $P=\left(p_{1}, \ldots, p_{r}\right) \in W$, then

$$
\left|\frac{\partial p_{j}}{\partial x_{k}}\right|<C_{2}, 1 \leq j \leq r, 1 \leq k \leq n, x \in B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right)
$$

Further, if $J_{P}(x)=\operatorname{det}\left(\frac{\partial p_{j}}{\partial x_{k}}\right)_{1 \leq j, k \leq r}$, we have

$$
J_{P}(x) \geq \frac{1}{2} \text { for } x \in B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right)
$$

If (1) and (2) hold, we also have the following estimate:
(3) There is a constant $C_{3}>0$ depending only on $n, r, D, C_{1}$ such that if $P \in W$ and $\varphi_{\nu}\left(x_{r+1}, \ldots, x_{n}\right), \nu=1, \ldots, r$, are the functions defined in (1), then

$$
\left|\frac{\partial \varphi_{\nu}}{\partial x_{i}}\right|<C_{3}, r<i \leq n, 1 \leq \nu \leq r,\left(x_{r+1}, \ldots, x_{n}\right) \in B_{n-r}\left(\theta_{2}\right)
$$

In fact, since, if we set $y=\left(x_{r+1}, \ldots, x_{n}\right) \in B_{n-r}\left(\theta_{2}\right)$, we have

$$
p_{j}\left(\varphi_{1}(y), \ldots, \varphi_{r}(y), y\right) \equiv 0,1 \leq j \leq r
$$

we see that, for $r<i \leq n, 1 \leq j \leq r$,

$$
\begin{equation*}
\frac{\partial p_{j}}{\partial x_{i}}+\sum_{\ell=1}^{r} \frac{\partial p_{j}}{\partial x_{\ell}} \cdot \frac{\partial \varphi_{\ell}}{\partial x_{i}} \equiv 0 \tag{4}
\end{equation*}
$$

the derivatives $\partial p_{j} / \partial x_{\nu}$ being evaluated at $\left(\varphi_{1}(y), \ldots, \varphi_{r}(y), y\right), y=$ $\left(x_{r+1}, \ldots, x_{n}\right)$.

Now, by (2), the matrix $\left(\frac{\partial p_{j}}{\partial x_{\ell}}\right)_{1 \leq j, \ell \leq r}$ is invertible on $B_{r}\left(\theta_{1}\right) \times$ $B_{n-r}\left(\theta_{2}\right)$ and its inverse, considered as a linear map of $\mathbb{R}^{n}$ into itself, has norm $\leq C_{4}$, where $C_{4}$ depends only on $n, r, D, C_{1}$. Hence (3) follows from (4) and (2).

We now formulate Theorem 2 of the Introduction with the present notation.

THEOREM 11.1. - There exist constants $\delta_{*}, C_{*}>0$ depending only on $n, r, D, C_{1}$ with $0<\delta_{*} \leq \theta_{2}$ such that, if $P \in W, f \in H^{D}$ and we set $F=f \circ \pi^{-1} \in C^{\infty}\left(B_{n-r}\left(\theta_{2}\right)\right)\left[F\left(x_{r+1}, \ldots, x_{n}\right)=f\left(\varphi_{1}(y), \ldots, \varphi_{r}(y), y\right)\right.$ with $\left.y=\pi(x)=\left(x_{r+1}, \ldots, x_{n}\right)\right]$, then the following inequalities hold:
(A) Polynomial Growth. If $0<\delta \leq \delta_{*}$, then

$$
\sup _{B_{n-r}(2 \delta)}|F| \leq C_{*} \sup _{B_{n-r}(\delta)}|F|
$$

(B) Bernstein's Inequality. If $\nabla F=\left(\frac{\partial F}{\partial x_{r+1}}, \ldots, \frac{\partial F}{\partial x_{n}}\right)$ is the gradient of $F$, then, for $0<\delta \leq \delta_{*}$,

$$
\sup _{B_{n-r}(\delta)}|\nabla F| \leq \frac{C_{*}}{\delta} \cdot \sup _{B_{n-r}(\delta)}|F| .
$$

(C) Equivalence of Norms. If $0<\delta \leq \delta_{*}$, then

$$
\sup _{B_{n-r}(\delta)}|F| \leq \frac{C_{*}}{\delta^{n-r}} \int_{B_{n-r}(\delta)}|F| d \lambda
$$

where $d \lambda$ is Lebesgue measure in $\mathbb{R}^{n-r}$.

Proof. - Let $C^{\prime}, D^{\prime}$ be as in the Extension Theorem, Part 1. Choose constants $\rho_{*}, \delta_{*}>0$ and $\rho_{0}>0$ depending only on $n, r, D, C_{1}$ such that $0<\rho_{0}, \rho_{*}<\rho_{1}$ and such that
(5.i) If $q \in H^{D^{\prime}}, q(0)=1,\|q\| \leq C^{\prime}$, then $\frac{1}{2}<q<2$ on $Q_{\rho_{*}}$.
(5.ii) If $P \in W$, then $\pi_{P}^{-1}\left(B_{n-r}\left(2 \delta_{*}\right)\right) \subset Q_{\rho_{*}}$.
(5.iii) If $P \in W$, then $\pi_{P}\left(Q_{\rho_{0}}\right) \subset B_{n-r}\left(\delta_{*}\right)$.

By the extension theorem (Part 1), given $P \in W$ and $f \in H^{D}$, we can find $G, q \in H^{D^{\prime}}$ such that

$$
q(0)=1,\|q\| \leq C^{\prime}, q \cdot f=G \text { on } V_{\rho_{1}}(P)
$$

and

$$
\begin{equation*}
\|G\| \leq C^{\prime} \rho_{0}^{-m} \sup _{V_{\rho_{0}}(P)}|f| \tag{7}
\end{equation*}
$$

By (5.iii), we have

$$
\begin{equation*}
\sup _{V_{\rho_{0}}(P)}|f| \leq \sup _{B_{n-r}\left(\delta_{*}\right)}\left|f \circ \pi^{-1}\right|, \pi=\pi_{P} \tag{8}
\end{equation*}
$$

Moreover,

$$
\begin{equation*}
\sup _{Q_{\rho_{*}}}|G| \leq \bar{C}\|G\| \tag{9}
\end{equation*}
$$

where $\bar{C}$ depends only on $\rho_{*}, n, D^{\prime}$. Using (5), (6), (7), (8) and (9), we obtain

$$
\begin{align*}
\sup _{B_{n-r}\left(2 \delta_{*}\right)}\left|f \circ \pi^{-1}\right| & \leq \frac{\sup _{Q_{\rho_{*}}}|G|}{\inf _{Q_{*}}|g|} \leq 2 \bar{C}\|G\|  \tag{10}\\
& \leq 2 \bar{C} C^{\prime} \rho_{0}^{-m} \sup _{V_{\rho_{0}}(P)}|f| \\
& \leq 2 \bar{C} C^{\prime} \rho_{0}^{-m} \sup _{B_{n-r}\left(\delta_{*}\right)}\left|f \circ \pi^{-1}\right| .
\end{align*}
$$

Next, if $F=f \circ \pi^{-1}$ and $r<\nu \leq n$, and we set $y=\left(x_{r+1}, \ldots, x_{n}\right)$, we have

$$
\begin{aligned}
&\left.\frac{\partial F}{\partial x_{\nu}}=\frac{\partial}{\partial x_{\nu}} G(\varphi(y), y)\right) \cdot(q(\varphi(y), y))^{-1} \\
&-(q(\varphi(y), y))^{-2} G(\varphi(y), y) \frac{\partial}{\partial x_{\nu}} q(\varphi(y), y)
\end{aligned}
$$

if we remark that $\left\|\frac{\partial G}{\partial x_{k}}\right\| \leq D^{\prime}\|G\|$ and $\left\|\frac{\partial q}{\partial x_{k}}\right\| \leq D^{\prime}\|q\| \leq D^{\prime} C^{\prime}(1 \leq k \leq n)$ and use (3) and (5.i), we obtain

$$
\sup _{B_{n-r}\left(\delta_{*}\right)}\left|\frac{\partial F}{\partial x_{\nu}}\right| \leq C^{\prime \prime}\|G\| \leq C_{\#} \sup _{B_{n-r}\left(\delta_{*}\right)}|F|
$$

by (7) and (8); here $C^{\prime \prime}, C_{\#}$ depend only on $n, r, D, C_{1}$. This gives

$$
\begin{equation*}
\sup _{B_{n-r}\left(\delta_{*}\right)}|\nabla F| \leq \sqrt{n} C_{\#} \sup _{B_{n-r}\left(\delta_{*}\right)}|F| . \tag{11}
\end{equation*}
$$

Next we make the following remark. Let $U$ be a bounded convex open set in $\mathbb{R}^{N}, U \neq \emptyset$. Let $\varphi=C^{1}(\bar{U}), \bar{U}$ being the closure of $U$. Then

$$
\begin{equation*}
\sup _{U}|\varphi| \leq \operatorname{diam}(U) \sup _{U}|\nabla \varphi|+\frac{1}{\operatorname{vol}(U)} \int_{U}|\varphi| d \lambda . \tag{12}
\end{equation*}
$$

To prove this, we first remark that it is sufficient to prove (12) when $\varphi$ is real-valued. In fact, if $x \in U$ and $\alpha \in \mathbb{R}$ is so chosen that $e^{i \alpha} \varphi(x) \in \mathbb{R}$, we can apply (12) to the function $\varphi_{1}=\operatorname{Re}\left(e^{i \alpha} \varphi\right)$ to obtain

$$
|\varphi(x)|=\left|\varphi_{1}(x)\right| \leq \operatorname{diam}(U)\left|\nabla \varphi_{1}\right|+\frac{1}{\operatorname{vol}(U)} \int_{U}\left|\varphi_{1}\right| d \lambda
$$

since $x \in U$ is arbitrary and $\left|\varphi_{1}\right| \leq|\varphi|,\left|\nabla \varphi_{1}\right| \leq|\nabla \varphi|$, we obtain (12) for $\varphi$.

Now, if $\psi \in C^{1}(\bar{U})$ is real-valued and $\int_{U} \psi d \lambda=0$, then, there is a point $a \in U$ with $\psi(a)=0$. Hence, if $x \in U$, there is a point $\xi$ on the line segment from $a$ to $x$ such that

$$
\psi(x)=\psi(x)-\psi(a)=\sum_{\nu=1}^{N}\left(x_{\nu}-a_{\nu}\right) \frac{\partial \psi}{\partial x_{\nu}}(\xi)
$$

which gives $|\psi(x)| \leq|x-a||\nabla \psi(\xi)| \leq \operatorname{diam}(U) \sup _{U}|\nabla \psi|$. If $\varphi \in C^{1}(\bar{U})$ and we apply this to $\psi(x)=\varphi(x)-\frac{1}{\operatorname{vol}(U)} \int_{U} \varphi d \lambda$, we obtain (12) when $\varphi$ is real-valued, hence in general as remarked above.

Let $\epsilon=\frac{1}{4 \sqrt{n} C_{\#}}$ where $C_{\#}$ is as in (11) and let $y_{0} \in \overline{B_{n-r}\left(\delta_{*}\right)}$ be such that $\left|F\left(y_{0}\right)\right|=\sup _{B_{n-r}\left(\delta_{*}\right)}|F|$. We apply (12) with $N=n-r, \varphi=F$ and $U=B_{n-r}\left(\delta_{*}\right) \cap\left\{y \in \mathbb{R}^{n-r}| | y-y_{0} \mid<\epsilon\right\}$. If we remark that there is a constant $\gamma>0$ depending only on $\epsilon, \delta_{*}$ and $n-r$ such that $\operatorname{vol}(U) \geq \gamma$, we obtain

$$
\begin{aligned}
\sup _{B_{n-r}\left(\delta_{*}\right)}|F|=\left|F\left(y_{0}\right)\right|=\sup _{U}|F| & \leq 2 \epsilon \sup _{U}|\nabla F|+\frac{1}{\gamma} \int_{U}|F| d \lambda \\
& \leq \frac{1}{2} \sup _{B_{n-r}\left(\delta_{*}\right)}|F|+\frac{1}{\gamma} \int_{U}|F| d \lambda \quad(\text { by }(11)) .
\end{aligned}
$$

Thus

$$
\begin{equation*}
\sup _{B_{n-r}\left(\delta_{*}\right)}|F| \leq \frac{2}{\gamma} \int_{U}|F| d \lambda . \tag{13}
\end{equation*}
$$

If we take $C_{*}=\max \left(\frac{2}{\gamma} \delta^{n-r}, \sqrt{n} C_{\#} \delta_{*}, 2 \bar{C} C^{\prime} \rho_{0}^{-m}\right),(10),(11),(13)$ give, respectively, statements (A), (B), (C) of Theorem 11.1 for the value $\delta=\delta_{*}$.

To prove these statements for $0<\delta \leq \delta_{*}$, we simply rescale as in the proof of Part 2 of the extension theorem. To do this, given $P=\left(p_{1}, \ldots, p_{r}\right) \in W$, let $\tilde{P}=\left(\tilde{p}, \ldots, \tilde{p}_{r}\right)$ where $\tilde{p}_{j}(x)=\frac{\delta_{*}}{\delta} p_{j}\left(\frac{\delta x}{\delta_{*}}\right)$. Then $\tilde{P} \in W$. We denote $\left(x_{r+1}, \ldots, x_{n}\right)=\pi_{P}\left(x_{1}, \ldots, x_{n}\right)$ by $y$. Let $\varphi=\left(\varphi_{1}, \ldots, \varphi_{r}\right): B_{n-r}\left(\theta_{2}\right) \rightarrow B_{r}\left(\theta_{1}\right)$ be the map given by (1), so that $Z(P) \cap B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right)=\left\{(\varphi(y), y) \mid y \in B_{n-r}\left(\theta_{2}\right)\right\}$. We claim that the $\operatorname{map}\left(\tilde{\varphi}_{1}, \ldots, \tilde{\varphi}_{r}\right)=\tilde{\varphi}: B_{n-r}\left(\theta_{2}\right) \rightarrow B_{r}\left(\theta_{1}\right)$ corresponding to $\tilde{P}$ is given by $\tilde{\varphi}_{j}(y)=\frac{\delta_{*}}{\delta} \varphi_{j}\left(\frac{\delta y}{\delta_{*}}\right)$. In fact, if $|y|$ is sufficiently small, this is clear, since, for $1 \leq j \leq r, \tilde{p}_{j}(\tilde{\varphi}(y), y)=\tilde{p}_{j}\left(\frac{\delta_{*}}{\delta} \varphi\left(\frac{\delta y}{\delta_{*}}\right), y\right)=0$ (by the definition of $\tilde{p}_{j}$ and $\varphi$ ); the above formula then holds for $|y|<\theta_{2}$ since both functions are real analytic.

If $f \in H^{D}$, set $\tilde{f}(x)=f\left(\frac{\delta x}{\delta_{*}}\right)$. Then $\tilde{f} \in H^{D}$. Let $\tilde{\pi}=\pi_{\tilde{P}}$ be the projection $Z(\tilde{P}) \cap B_{r}\left(\theta_{1}\right) \times B_{n-r}\left(\theta_{2}\right) \rightarrow B_{n-r}\left(\theta_{2}\right)$ corresponding to $\tilde{P}$. We have, for $|y|<\theta_{2}$,

$$
\tilde{F}(y)=\tilde{f} \circ \tilde{\pi}^{-1}(y)=f\left(\varphi\left(\frac{\delta y}{\delta_{*}}\right), \frac{\delta y}{\delta_{*}}\right)=\left(f \circ \pi^{-1}\right)\left(\frac{\delta y}{\delta_{*}}\right)=F\left(\frac{\delta y}{\delta_{*}}\right)
$$

It follows that

$$
\begin{gathered}
\sup _{B_{n-r}\left(t \delta_{*}\right)}|\tilde{F}|=\sup _{B_{n-r}(t \delta)}|F|, t=1,2 \\
(\nabla \tilde{F})(y)=\frac{\delta}{\delta_{*}}(\nabla F)\left(\frac{\delta y}{\delta_{*}}\right)
\end{gathered}
$$

and

$$
\int_{B_{n-r}\left(\delta_{*}\right)}|\tilde{F}| d \lambda=\left(\frac{\delta_{*}}{\delta}\right)^{n-r} \int_{B_{n-r}(\delta)}|F| d \lambda
$$

Consequently (A), (B), (C) in Theorem 11.1 for the value $\delta_{*}$ and the pair $(\tilde{F}, \tilde{P})$ are equivalent, respectively, to (A), (B), (C) for the value $\delta$ and the pair $(F, P)$.

This proves Theorem 11.1.

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