# HIGH POINTS IN THE HISTORY OF VALUE DISTRIBUTION THEORY OF SEVERAL COMPLEX VARIABLES 

Wilhelm Stoll<br>Inaugural Lecture

Timothy O'Meara and Frank Castellino thank your for your kind introduction. I am deeply moved by your words and by the appointment to the chair. Foremost I thank the donors Vincent J. Duncan and Annamarie Micus Duncan for their generosity. My colleagues and I are most grateful for this recognition of our work by the donors and the administration of the University.

Ladies and gentlemen, colleagues, speakers and participants! This inaugural address opens the Symposium on Value Distribution Theory in Several Complex Variables sponsored by the University of Notre Dame. Welcome to all of you. An inaugural address, an Antrittsvorlesung, so late in life seems to be out of place and perhaps should be called an Abschiedsvorlesung. Yet, hopefully, this is premature and I can be around a few more years. Taking the hint, I will look backwards and recall some of the high points in the development of the theory. Time permits only a few topics.

Looking backwards, out of the mist of time there emerges not an abstract theory but the lively memory of those who taught me mathematics: Siegfried Kerridge, Wilhelm Germann, Wilhelm Schweizer and later at the University Hellmuth Kneser, Konrad Knopp, Erich Kamke, G. G. Lorentz and Max Müller. Also there appear those who inspired me but who were not directly my teachers: Heinz Hopf, Hermann Weyl, Rolf Nevanlinna and one who is right here with us: Shiing-shen Chern, we all welcome you. Thirty years ago you recruited me for Notre Dame. You supported the growth of this department in many ways. Your work on value distribution in several

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complex variables counts as one of your many marvelous contributions to mathematics. Thank you for coming.

The giants of the 19th century created the theory of entire functions. In this century, in 1925, with a stroke of genius, Rolf Nevanlinna extended this theory to a value distribution theory of meromorphic functions. His two Main Theorems are the foundation upon which Nevanlinna theory rests.

In 1933, Henri Cartan [8] proved Nevanlinna's Second Main Theorem for the case of holomorphic curves. If we view curves belonging to the theory of several dependent variables, then Cartan's paper provides the first theorem in the theory of value distribution in several complex variables. Thus let me outline his result. However, I shall use today's terminology and advancement.

For each $0<r \in \mathbb{R}$ define the discs and circle

$$
\begin{align*}
\mathbb{C}[r] & =\{z \in \mathbb{C}| | z \mid \leq r\} & \mathbb{C}(r) & =\{z \in \mathbb{C}| | z \mid<r\}  \tag{1}\\
\mathbb{C}<r> & =\{z \in \mathbb{C}| | z \mid=r\} & \mathbb{C}_{*} & =\mathbb{C}-\{0\} .
\end{align*}
$$

An integral valued function $\nu: \mathbb{C} \rightarrow \mathbb{Z}$ is said to be a divisor if

$$
\begin{equation*}
S=\operatorname{supp} \nu=\operatorname{clos}\{z \in \mathbb{C} \mid \nu(z) \neq 0\} \tag{3}
\end{equation*}
$$

is a closed set of isolated points in $\mathbb{C}$. For all $r \geq 0$ the counting function $n_{\nu}$ of $\nu$ is defined by the finite sum

$$
\begin{equation*}
n_{\nu}(r)=\sum_{z \in \mathbb{C}[r]} \nu(z) . \tag{4}
\end{equation*}
$$

For $0<s<r \in \mathbb{R}$, the valence function $N_{\nu}$ of $\nu$ is defined by

$$
\begin{equation*}
N_{\nu}(r, s)=\int_{s}^{r} n_{\nu}(t) \frac{d t}{t} . \tag{5}
\end{equation*}
$$

If $h \not \equiv 0$ is an entire function, let $\mu_{h}(z)$ be the zero-multiplicity of $h$ at $z$. Then $\mu_{h}: \mathbb{C} \rightarrow \mathbb{Z}$ is a non-negative divisor called the zero divisor of $h$.

The exterior derivative $d=\partial+\bar{\partial}$ on differential forms twists to

$$
\begin{equation*}
d^{c}=\frac{i}{4 \pi}(\bar{\partial}-\partial) \tag{6}
\end{equation*}
$$

on complex manifolds. Define $\tau_{o}: \mathbb{C} \rightarrow \mathbb{R}$ by $\tau_{o}(z)=|z|^{2}$ for $z \in \mathbb{C}$. Define

$$
\begin{equation*}
\sigma=d^{c} \log \tau_{0} \tag{7}
\end{equation*}
$$

If $r>0$, then

$$
\begin{equation*}
\int_{\mathbb{C}<r>} \sigma=1 . \tag{8}
\end{equation*}
$$

If $h \not \equiv 0$ is an entire function and if $r>o$, the Jensen Formula

$$
\begin{equation*}
N_{\mu_{h}}(r, s)=\int_{\mathbb{C}<r>} \log h \sigma-\int_{\mathbb{C}<s>} \log h \sigma \tag{9}
\end{equation*}
$$

is a forerunner of Nevanlinna's First Main Theorem.
Let $V$ be a normed, complex vector space of finite dimension $n+1>1$. Put $V_{*}=V-\{0\}$. Then the multiplicative group $\mathbb{C}_{*}$ acts on $V_{*}$. The quotient space $\mathbb{P}(V)=V_{*} / \mathbb{C}_{*}$ is the associated projective space. The quotient map $\mathbb{P}: V_{*} \rightarrow \mathbb{P}(V)$ is open and holomorphic. If $M \subseteq V$, put $\mathbb{P}(M)=P\left(M \cap V_{*}\right)$. If $W$ is a linear subspace of $V$ with dimension $p+1$, then $\mathbb{P}(W)$ is called a p-plane of $\mathbb{P}(V)$. If $p=n-1$, then $\mathbb{P}(W)$ is called a hyperplane. The dual complex vector space $V^{*}$ of $V$ consists of all $\mathbb{C}$-linear functions $\mathfrak{a}: V \rightarrow \mathbb{C}$. Here $\|\mathfrak{a}\|$ is the smallest real number such that $|\mathfrak{a}(\mathfrak{x})| \leq\|\mathfrak{a}\|\|\mathfrak{x}\|$ for all $\mathfrak{x} \in V$. Then $\left\|\|\right.$ is a norm on $V^{*}$. Also write $<\mathfrak{x}, \mathfrak{a}>=\mathfrak{a}(\mathfrak{x})$. Here $<\mathfrak{d}, \mathfrak{a}>=<\mathfrak{a}, \mathfrak{x}>$ indicates $\left(V^{*}\right)^{*}=V$. If $a=\mathbb{P}(\mathfrak{a}) \in \mathbb{P}\left(V^{*}\right)$, then $E[a]=\mathbb{P}(\operatorname{ker} \mathfrak{a})$ is a hyperplane in $\mathbb{P}(V)$. The assignment $a \rightarrow E[a]$ parameterizes the set of hyperplanes bijectively. The distance from $x=\mathbb{P}(\mathfrak{x}) \in \mathbb{P}(V)$ to $E[a]$ is measured by

$$
\begin{equation*}
0 \leq \square x, a \square=\frac{|\langle\mathfrak{x}, a\rangle|}{\|\mathfrak{x}\|\|\mathfrak{a}\|} \leq 1 . \tag{10}
\end{equation*}
$$

Let $f: \mathbb{C} \rightarrow \mathbb{P}(V)$ be a holomorphic map. A holomorphic map $\mathfrak{v}: \mathbb{C} \rightarrow V_{*}$ is called a reduced representation of $f$ if $\mathbb{P} \circ \mathfrak{v}=f$. A reduced representation exists. Then $\mathfrak{w}: \mathbb{C} \rightarrow V_{*}$ is a reduced representation of $f$ if and only if there is a holomorphic function $h: \mathbb{C} \rightarrow \mathbb{C}_{*}$ without zeroes such that $\mathfrak{t v}=h \mathfrak{v}$. For $0<s<r \in \mathbb{R}$ the characteristic function of $f$ is defined by

$$
\begin{equation*}
T_{f}(r, s)=\int_{\mathfrak{C}<r>} \log \|\mathfrak{v}\| \sigma-\int_{\mathbb{C}<s>} \log \|\mathfrak{b}\| \sigma . \tag{11}
\end{equation*}
$$

By (9), $T_{f}(r, s)$ does not depend on the choice of $\mathfrak{v}$, Since $\log \|\mathfrak{v}\|$ is subharmonic, $T_{f} \geq 0$. If $f$ is constant, $\mathfrak{b}$ can be taken as a constant. Hence $T_{f}(r, s) \equiv 0$. If $f$ is not constant, then $T_{f}(r, s)>0$ and $T_{f}(r, s) \rightarrow \infty$ for $r \rightarrow \infty$. If \||| and ||| ||| are two norms on $V$, there are constants $C_{2} \geq C_{1}>0$ such that $C_{1}\| \| \mathfrak{x}\|\leq\| \mathfrak{x}\left\|\leq C_{2} \mid\right\| \mathfrak{x} \|$ for all $\mathfrak{x} \in V$. Put $C=\log C_{2} / C_{1} \geq 0$. If $0<s<r$, then

$$
\begin{equation*}
\left|T_{f}(r, s,\| \|)-T_{f}(r, s,|\|\mid\|) \mid \leq C .\right. \tag{12}
\end{equation*}
$$

Let $f: \mathbb{C} \rightarrow \mathbb{P}(V)$ and $g: \mathbb{C} \rightarrow \mathbb{P}\left(V^{*}\right)$ be holomorphic maps. They are called free if $f(z) \notin E[g(z)]$ for some $z \in \mathbb{C}$. Take reduced representations $\mathfrak{v}$ of $f$ and $\mathfrak{w}$ of $g$, then $(f, g)$ is free if and only if $\langle\mathfrak{v}, \mathfrak{w}\rangle=h \not \equiv 0$. If so, the intersection divisor $\mu_{f, g}=\mu_{h} \geq 0$ does not depend on the choices of $\mathfrak{v}$ and $\mathfrak{w}$. Its, counting function and its valence function are abbreviated by $n_{f, g}$ and $N_{f, g}$ respectively. The pair $(f, g)$ is free if and only if $\square f, g, \square \not \equiv 0$. If so, for $r>0$ the compensation function $m_{f, g}$ of $(f, g)$ is defined by

$$
\begin{equation*}
m_{f, g}(r)=\int_{\mathbb{C}<r>} \log \frac{1}{\square f, g \square} \sigma \geq 0 . \tag{13}
\end{equation*}
$$

For $0<s<r$, the identities (9), (11) and (13) imply the First Main Theorem

$$
\begin{equation*}
T_{f}(r, s)+T_{g}(r, s)=N_{f, g}(r, s)+m_{f, g}(r)-m_{f, g}(s) \tag{14}
\end{equation*}
$$

Cartan [8] considered the case of constant $g \equiv a \in \mathbb{P}\left(V^{*}\right)$ only which yields

$$
\begin{equation*}
T_{f}(r, s)=N_{f, a}(r, s)+m_{f, a}(r)-m_{f, a}(s) \tag{15}
\end{equation*}
$$

which Cartan [8] mentions only implicitely. If $n=1$, Rolf Nevanlinna proved (15) in [32] (1925).

If $f$ or $g$ or both are not constant and if $(f, g)$ is free the defect is defined by

$$
\begin{align*}
0 \leq \delta(f, g) & =\varliminf_{r \rightarrow \infty} \frac{m_{f, g}(r)}{T_{f}(r, s)+T_{g}(r, s)} \\
& =1-\overline{\lim _{r \rightarrow \infty}} \frac{N_{f, g}(r, s)}{T_{f}(r, s)+T_{g}(r, s)} \leq 1 \tag{16}
\end{align*}
$$

The map $g$ is said to grow slower than $f$, if $T_{g}(r, s) / T_{f}(r, s) \rightarrow 0$ for $r \rightarrow \infty$. By (12), the defect does not depend on the choice of the norm on $V$. Also the defect is independent of $s$. Observe that $\mu_{f, g}=\mu_{g, f}, n_{f, g}=n_{g, f}, N_{f, g}=N_{g, f}, m_{f, g}=m_{g, f}$ and $\delta(f, g)=$ $\delta(g, f)$. Since most investigators concentrate on constant $g$ or on the case where $g$ grows slower than $f$, this symmetry is little known.

Since the choice of the norm on $V$ does not matter, we can choose a hermitian norm which comes from a positive definite hermitian form $(\cdot \mid \cdot): V \times V \rightarrow \mathbb{C}$ with $\|\mathfrak{x}\|^{2}=(\mathfrak{x} \mid \mathfrak{x})$ for $\mathfrak{x} \in V$. Define $\tau: V \rightarrow \mathbb{C}$ by $\tau(\mathfrak{c})=\|\mathfrak{x}\|^{2}$ for $\mathfrak{x} \in V$. Then $\tau$ is of class $C^{\infty}$. There is one and only one positive form $\Omega$ of bidegree $(1,1)$ on $\mathbb{P}(V)$, called the Fubini Study form such that $d d^{c} \log \tau=\mathbb{P}^{*}(\Omega)$ on $V_{*}$. Let $\mathfrak{v}: \mathbb{C} \rightarrow V_{*}$ be a reduced representation of $f$. Then $f=\mathbb{P} \circ \mathfrak{b}$ implies

$$
\begin{equation*}
d d^{c} \log \|\mathfrak{v}\|^{2}=\mathfrak{v}^{*}\left(\mathbb{P}^{*}(\Omega)\right)=f^{*}(\Omega) \tag{17}
\end{equation*}
$$

If Stokes theorem and fiber integration are applied to (11) we obtain the Ahlfors-Shimizu definition of the characteristic function of $f$

$$
\begin{equation*}
T_{f}(r, s)=\int_{s}^{r} \int_{\mathbb{C}[t]} f^{*}(\Omega) \frac{d t}{t} \quad \text { for } 0<s<r . \tag{18}
\end{equation*}
$$

Here $A_{f}(t)=\int_{\mathbb{C}[t]} f^{*}(\Omega) \geq 0$ increases. Put $A_{f}(\infty)=\lim _{t \rightarrow \infty} A_{f}(t) \leq$ $\infty$. Then

$$
\begin{equation*}
\lim _{r \rightarrow \infty} \frac{T_{f}(r, s)}{\log r}=A_{f}(\infty) \tag{19}
\end{equation*}
$$

Now $f$ is constant if and only if $A_{f}(\infty)=0$ and $f$ is rational if and only if $A_{f}(\infty)<\infty$.

Let $\mathfrak{U}=\left\{a_{j}\right\}_{j \in Q}$ be a family of points $a_{j} \in \mathbb{P}\left(V^{*}\right)$ representating hyperplanes. If $P \subseteq Q$, define $\mathfrak{U}_{P}=\left\{a_{j}\right\}_{j \in P}$. For each $j \in Q$ pick $\mathfrak{a}_{j} \in V_{*}^{*}$ with $a_{j}=\mathbb{P}\left(\mathfrak{a}_{j}\right)$. Our definitions will not depend on the choice of $\mathfrak{a}_{j}$. Put $q=\# Q$. Then $\mathfrak{A}$ is said to be linearly independent if there is a bijective map $\lambda: \mathbb{N}[1, q] \rightarrow Q$ such that $\mathfrak{a}_{\lambda(1)}, \ldots, \mathfrak{a}_{\lambda(q)}$ are linearly independent. If so, then $q \leq n+1$. Moreover $\mathfrak{A}$ is said to be basic if $\mathfrak{A}$ is linearly independent and $q=n+1$. Moreover $\mathfrak{A}$ is said to be in general position if $\mathfrak{U}_{P}$ is linearly independent for each $P \subseteq Q$ with $0<\# P \leq n+1$. If $N$ is an integer and if $q>N \geq n$, then $\mathfrak{U}$ is said to be in $N$-subgeneral position (Chen [9]) if for every subset $S$ of $Q$ with $\# S=N+1$, there is a subset $P$ of $S$ such that $\mathfrak{\vartheta}_{P}$ is basic.

Let $f: \mathbb{C} \rightarrow \mathbb{P}(V)$ be a holomorphic map. Then there is a unique linear subspace $W$ of smallest dimension $k+1$ of $V$ such that $f(\mathbb{C}) \subseteq \mathbb{P}(W)$. Then $f$ is said to be $k$-flat. If $k=n$, then $W=V$ and $f$ is said to be linearly non-degenerated.

Take $0 \leq s \in \mathbb{R}$. Let $G: \mathbb{R}[s,+\infty) \rightarrow \mathbb{R}$ and $H: \mathbb{R}[s,+\infty)$ be functions. Then $G \leq H$ means that there is a subset $E$ of finite measure of $\mathbb{R}_{+}=\mathbb{R}[0,+\infty)$ such that $G(r) \leq H(r)$ for all $r \in \mathbb{R}[s,+\infty)-E$.

## Second Main Theorem (Cartan [8] 1933)

Let $V$ be a hermitian vector space of dimension $n+1>1$. Let $f: \mathbb{C} \rightarrow \mathbb{P}(V)$ be a linearly non-degenerated, holomorphic map. Let $\mathfrak{U}=\left\{a_{j}\right\}_{j \in Q}$ be a finite family of "hyperplanes" $a_{j} \in \mathbb{P}\left(V^{*}\right)$ in general position with $n+1<q=\# Q<\infty$. Take $s>0$ and $\varepsilon>0$. Then

$$
\begin{align*}
\sum_{j \in Q} m_{f, a_{j}}(r) \leqq & (n+1) T_{f}(r, s)  \tag{20}\\
& +\frac{1}{2} n(n+1)(1+\varepsilon) \log T_{f}(r, s)+\varepsilon \log r .
\end{align*}
$$

As a consequence, we obtain trivially

## Defect Relation (Cartan [8] 1933)

Under the assumptions of the Second Main Theorem we have

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq n+1 \tag{21}
\end{equation*}
$$

If $f: \mathbb{C} \rightarrow \mathbb{P}(V)$ is only k -flat, and if $\mathfrak{N}$ is in general position such that $\left(f, a_{j}\right)$ is free for each $j \in Q$, Henri Cartan conjectured in 1933 that

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq 2 n-k+1 \tag{22}
\end{equation*}
$$

which was proven by Nochka [35] in 1982. Thus if $\# Q \geq 2 n+1$ and $f(\mathbb{C}) \cap E\left[a_{j}\right]=0$ for all $j \in Q$, then $2 n+1 \leq 2 n-k+1$. Therefore $k=0$ and $f$ is constant. Hence

$$
\begin{equation*}
\mathbb{P}(V)-\bigcup_{j \in Q} E\left[a_{j}\right] \tag{23}
\end{equation*}
$$

is Brody-hyperbolic. In fact by a theorem of Chen [9] (22) can be improved:

## Defect Relation of Cartan-Nochka-Chen

Let $V$ be a hermitian vector space of dimension $n+1>1$. Let $f: \mathbb{C} \rightarrow \mathbb{P}(V)$ be a k-flat, holomorphic map. Let $\mathfrak{N}=\left\{a_{j}\right\}_{j \in Q}$ be a finite family of "hyperplanes" $a_{j} \in \mathbb{P}\left(V^{*}\right)$ in $N$-subgeneral position with $N \geq n$ and $N+1 \leq \# Q=q<\infty$. Assume that $\left(f, a_{j}\right)$ is free for each $j \in Q$. Then

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq 2 N-k+1 \tag{24}
\end{equation*}
$$

An alternative proof of the defect relation (21) was given by Ahlfors [1] in 1941. Also he proves a defect relation for associated maps. His proof is very powerful and works in more general situations.

Hermann and Joachim Weyl [90] lifted Ahlfors's proof to Riemann surfaces. It was simplified by H. Wu [92] in 1970, Cowen and Griffiths [17] in 1976 and Pit-Mann Wong [93] in 1976. I extended this AhlforsWeyl theory to non-compact Kaehler manifolds [65]. However first we have to inquire how value distribution was extended to functions and maps of several independent complex variables.

Hellmuth Kneser created such an extension in two fundamental papers [23] in 1936 and [24] in 1938. Although these papers are little remembered today, they still influence the present research in value distribution of several independent complex variables. Therefore let me explain his fundamental ideas. Again I will cast them in modern terminology and perspective.

Let $M$ be a connected, complex manifold of dimension $m$. Let $f \not \equiv 0$ be a holomorphic function on $M$. Take $p \in M$. Let $\alpha: U^{\prime} \rightarrow U$ be a biholomorphic map of an open ball $U^{\prime}$ in $\mathbb{C}^{m}$ centered at 0 onto an open subset $U$ of $M$ with $\alpha(0)=p$. Then for each integer $\lambda \geq 0$ there is a unique homogeneous polynomial $P_{\lambda}$ of degree $\lambda$ such that

$$
\begin{equation*}
f \circ \alpha=\sum_{\lambda=0}^{\infty} P_{\lambda} \tag{25}
\end{equation*}
$$

where the convergence is uniform on every compact subset of $U^{\prime}$. Since $f \mid U \not \equiv 0$, there is a unique number $\mu=\mu_{f}(p) \geq 0$ depending on $f$ and $p$ only such that $P_{\mu} \not \equiv 0$ and $P_{\lambda} \equiv 0$ for all $\lambda \in \mathbb{Z}$ with $0 \leq \lambda<\mu$. The number $\mu_{f}(p)$ is called the zero-multiplicity of $f$ at $p$ and the function $\mu_{f}: M \rightarrow \mathbb{Z}$ is called the zero-divisor of $f$.

An integral valued function $\nu: M \rightarrow \mathbb{Z}$ is said to be a divisor on $M$ if and only if for every point $p \in M$ there is an open, connected neighborhood $U$ of $p$ with holomorphic functions $g \not \equiv 0$ and $h \not \equiv 0$ on $U$ such that

$$
\begin{equation*}
\nu \mid U=\mu_{g}-\mu_{h} . \tag{26}
\end{equation*}
$$

Let $S$ be the support of $\nu$. Then $S=\emptyset$ if and only if $\nu \equiv 0$. If $S \neq \emptyset$, then $S$ is a pure ( $\mathrm{m}-1$ )-dimensional analytic subset of $M$. Let $\mathfrak{R}(S)$ be the set of regular points of $S$ and let $\sum(S)=S-\Re(S)$ be the set of singular points of $S$. Then $\nu \mid \Re(S)$ is locally constant.

Let $\tau: M \rightarrow \mathbb{R}_{+}$be an unbounded, non-negative function of class $C^{\infty}$ on $M$. If $B \subseteq M$ and $0 \leq r \in \mathbb{R}$, abbreviate

$$
\begin{array}{rll}
B[r] & =\left\{x \in B \mid \tau(x) \leq r^{2}\right\} &  \tag{27}\\
B(r)=\left\{x \in B \mid \tau(x)<r^{2}\right\} \\
B<r> & =\left\{x \in B \mid \tau(x)=r^{2}\right\} & \\
B_{*}=\{x \in B \mid \tau(x)>0\}
\end{array}
$$

Here $\tau$ is called an exhaustion of $M$ if and only if $M[r]$ is compact for each $r>0$. Abbreviate

$$
\begin{equation*}
v=d d^{c} \tau \quad \omega=d d^{c} \log \tau \quad \sigma=d^{c} \log \tau \wedge \omega^{m-1} \tag{29}
\end{equation*}
$$

Then $d \sigma=\omega^{m}$. The function $\tau$ is said to be parabolic if and only if

$$
\begin{equation*}
\omega \geq 0 \quad \omega^{m} \equiv 0 \quad v^{m} \not \equiv 0 . \tag{30}
\end{equation*}
$$

If so, then $v \geq 0$. More over $\tau$ is said to be strictly parabolic if and only if $\tau$ is parabolic and $v>0$ on $M$. If $\tau$ is an exhaustion and parabolic, then $(M, \tau)$ is said to be a parabolic manifold. If so, there is a constant $\varsigma>0$ such that

$$
\begin{equation*}
\int_{M[r]} v^{m}=\varsigma r^{2 m} \tag{31}
\end{equation*}
$$

for all $r \geq 0$. Then for almost all $r>0$ we have

$$
\begin{equation*}
\int_{M<r>} \sigma=\varsigma . \tag{32}
\end{equation*}
$$

In 1973 Griffiths and King [19] introduced parabolic manifolds. The concept was expanded in [75]. If $\tau$ is an exhaustion and strictly parabolic function, $(M, \tau)$ is said to be a strictly parabolic manifold. In [77] 1980 I showed that $(M, \tau)$ is strictly parabolic if and only if there is a hermitian vector space $W$ of dimension $m$ and a biholomorphic map $h: M \rightarrow W$ such that $\tau=\|h\|^{2}$. We assume that $(M, \tau)$ is strictly parabolic and we identify $M=W$ such that $h$ becomes the identity. In this case $\varsigma=1$ and $M$ is a hermitian vector
space, which was Kneser's starting point. We assume that $m>1$. If $u: M<1>\rightarrow \mathbb{C}$ is a function such that $u \sigma$ is integrable over the unit sphere $M<1\rangle$, the mean value of $u$ is defined by

$$
\begin{equation*}
\mathfrak{M}(u)=\int_{M<1>} u \sigma \tag{33}
\end{equation*}
$$

Let $V$ be a hermitian vector space of dimension $n+1>1$. Let $f: M \rightarrow \mathbb{P}(V)$ and $g: M \rightarrow \mathbb{P}\left(V^{*}\right)$ be meromorphic maps. Let $I_{f}$ and $I_{g}$ be the indeterminacies of $f$ and $g$ respectively. Then $(f, g)$ is said to be free if there exists $z \in M-I_{f} \cup I_{g}$ such that $f(z) \notin E[g(z)]$. For each "unit" vector $\mathfrak{b} \in M<1>$ an isometric embedding $j_{\mathfrak{b}}: \mathbb{C} \rightarrow M$ is defined by $j_{\mathfrak{b}}(z)=z \mathfrak{b}$ for $z \in \mathbb{C}$. If $j_{\mathfrak{b}}(\mathbb{C}) \nsubseteq I_{f} \cup I_{g}$ the pull back holomorphic maps $f_{\mathfrak{b}}=j_{\mathfrak{b}}^{*}(f): \mathbb{C} \rightarrow$ $\mathbb{P}(V)$ and $g_{\mathfrak{b}}=j_{\mathfrak{b}}^{*}(g): \mathbb{C} \rightarrow \mathbb{P}\left(V^{*}\right)$ exist and $\left(f_{\mathfrak{b}}, g_{\mathfrak{b}}\right)$ is free for almost all $\mathfrak{b} \in M<1>$. If $0<s<r$ the First Main Theorem holds
(34) $T_{f_{\mathfrak{b}}}(r, s)+T_{g_{\mathfrak{b}}}(r, s)=\bar{N}_{f_{b}, g_{\mathfrak{b}}}(r, s)+m_{f_{6}, g_{\mathfrak{b}}}(r)-m_{f_{6}, g_{\mathfrak{b}}}(s)$

Now Kneser [24] applied the operator $\mathfrak{M}$ termwise in (34) to obtain the respective value distribution functions and the First Main Theorem

$$
\begin{equation*}
T_{f}(r, s)+T_{g}(r, s)=N_{f, g}(r, s)+m_{f, g}(r)-m_{f, g}(s) \tag{35}
\end{equation*}
$$

Of course Kneser considered the case $n=1$ only. Then $f$ is a meromorphic function. Also he assumed that $g \equiv a \in \mathbb{P}_{1}$ is constant. Had he stopped with the above derivation of (35), his result would have been worthless. He proceeded and expressed the value distribution functions in meaningful analytic and geometric terms. This made the paper successful.

Let $\Omega$ be the Fubini Study form on $\mathbb{P}(V)$. For $t>0$ define $A_{f}$ by

$$
\begin{equation*}
A_{f}(t)=\frac{1}{t^{2 m-2}} \int_{M[t]} f^{*}(\Omega) \wedge v^{m-1} \geq 0 \tag{36}
\end{equation*}
$$

He showed that $A_{f}$ increases. Hence the limits

$$
\begin{align*}
& 0 \leq \lim _{t \rightarrow 0} A_{f}(t)=A_{f}(0)<\infty \\
& 0 \leq \lim _{t \rightarrow \infty} A_{f}(t)=A_{f}(\infty) \leq \infty \tag{37}
\end{align*}
$$

exist. Kneser obtained the identity

$$
\begin{equation*}
A_{f}(t)=\int_{M[t]} f^{*}(\Omega) \wedge \omega^{m-1}+A_{f}(0) \tag{38}
\end{equation*}
$$

for $t>0$. Here $f$ is constant if and only if $A_{f}(\infty)=0$ and $f$ is rational if and only if $A_{f}(\infty)<\infty$. Kneser proved

$$
\begin{equation*}
T_{f}(r, s)=\int_{s}^{r} A_{f}(t) \frac{d t}{t} \tag{39}
\end{equation*}
$$

for $0<s<r$. Moreover we have

$$
\begin{equation*}
\lim _{r \rightarrow \infty} \frac{T_{f}(r, s)}{\log r}=A_{f}(\infty) \tag{40}
\end{equation*}
$$

A holomorphic map $\mathfrak{v}: M \rightarrow V$ is said to be a reduced representation of $f$ if and only if $\operatorname{dim} \mathfrak{v}^{-1}(0) \leq m-2$ and $f(z)=\mathbb{P}(\mathfrak{v}(z))$ for all $z \in M-I_{f}$ with $\mathfrak{v}(z) \neq 0$. In fact $I_{f}=\mathfrak{v}^{-1}(0)$. Reduced representations exist since $M$ is a vector space. If $\mathfrak{v}$ is a reduced representation of $f$, any other reduced representation is given by $h \mathfrak{v}$, where $h: M \rightarrow \mathbb{C}_{*}$ is an entire function without zeroes. If $0<s<r$, then

$$
\begin{equation*}
T_{f}(r, s)=\int_{M<r>} \log \|\mathfrak{v}\| \sigma-\int_{M<s>} \log \|\mathfrak{v}\| \sigma . \tag{41}
\end{equation*}
$$

Since $(f, g)$ is free, $\square f, g \square \not \equiv 0$, and for $r>0$ the compensation function $m_{f, g}$ of $f, g$ is defined

$$
\begin{equation*}
m_{f, g}(r)=\int_{M<r>} \log \frac{1}{\square f, g \square} \sigma \geq 0 . \tag{42}
\end{equation*}
$$

Let $\nu: M \rightarrow \mathbb{Z}$ be a divisor with support $S$. Fot $t>0$ the
counting function $n_{\nu}$ of $\nu$ is defined by

$$
\begin{equation*}
n_{\nu}(t)=\frac{1}{t^{2 m-2}} \int_{S[t]} \nu v^{m-1}=\int_{S[t]} \nu \omega^{m-1}+n_{\nu}(0), \tag{43}
\end{equation*}
$$

where the limit $n_{\nu}(0)=\lim _{t \rightarrow 0} n_{\nu}(t)$ exists. Actually since $M$ is a vector space, $n_{\nu}(0)=\nu(0)$ (see Stoll [62]). For each $\mathfrak{b} \in M<1>$ with $j_{\mathfrak{b}}(\mathbb{C}) \nsubseteq S$, the pullback divisor $\nu_{\mathfrak{b}}=j_{\mathfrak{b}}^{*}(\nu)$ exists. If $t>0$ then

$$
\begin{equation*}
n_{\nu}(t)=\int_{\mathfrak{b} \in M<1>} n_{\nu_{\mathfrak{b}}}(t) \sigma(\mathfrak{b}) . \tag{44}
\end{equation*}
$$

Thus for $0<s<r$ the valence function $n_{\nu}$ of $\nu$ is given by

$$
\begin{equation*}
N_{\nu}(r, s)=\int_{\mathfrak{b} \in M<1>} N_{\nu_{\mathfrak{b}}}(r, s) \sigma(\mathfrak{b})=\int_{s}^{r} n_{\nu}(t) \frac{d t}{t} \tag{45}
\end{equation*}
$$

Take reduced resprentations $\mathfrak{v}: M \rightarrow \mathbb{P}(V)$ of $f$ and $\mathfrak{w}: M \rightarrow$ $\mathbb{P}\left(V^{*}\right)$ of $g$. Since $(f, g)$ is free, $h=<\mathfrak{v}, \mathfrak{w}>\not \equiv 0$. Then $\mu_{f, g}=\mu_{h}$ depends on $f$ and $g$ only. Put $S=h^{-1}(0)$. If $\mathfrak{b} \in M<1>$ with $j_{\mathfrak{b}}(\mathbb{C}) \nsubseteq S$, then $\mu_{f_{b}, g_{\mathfrak{b}}}=j_{\mathfrak{b}}^{*}\left(\mu_{f, g}\right)$. Hence

$$
\begin{equation*}
N_{f, g}(r, s)=\int_{\mathfrak{b} \in M<1>} N_{f_{6}, \mathfrak{g}_{\mathfrak{b}}}(r, s) \sigma(\mathfrak{b})=N_{\mu_{f, g}}(r, s) . \tag{46}
\end{equation*}
$$

Thus each term in (35) is explicitely expressed.
Actually, Kneser [24] provided a more general version of (42).
For $t>0$ the counting function of a pure p-dimensional analytic set $S$ in $M$ is defined by

$$
\begin{equation*}
n_{S}(t)=\frac{1}{t^{2 p}} \int_{S[t]} \bar{v}^{p}=\int_{S[t]} \omega^{p}+n_{S}(0) \tag{47}
\end{equation*}
$$

where

$$
\begin{equation*}
n_{S}(0)=\lim _{t \rightarrow 0} n_{S}(t) \tag{48}
\end{equation*}
$$

exists and is called the Lelong Number of $S$ at 0 . Kneser assumed that $0 \notin S$, then $n_{S}(0)=0$. Pierre Lelong permitted $0 \in S$ and proved (46) in 1957 [26] by the use of currents. Paul Thie [87] (1967) moved that the Lelong number is an integer. This result constituted Paul Thie's theses at Notre Dame and by coincidence Pierre Lelong was present at the defense of the theses. Of course, if $0 \in S$ then $n_{S}(0)>0$. Paul Thie's result proved to be most helpful in estimating volumes from below. Of course the Lelong number of $S$ can be defined for every $x \in M$ and shall be denoted by $L_{S}(x)$. Yum-Tong Siu [56] (1974) proved that the sets $\left\{x \in M \mid L_{S}(x) \geq q\right\}$ is analytic for every $q \in \mathbb{N}$. The proof was simplified by Lelong [28]

Since $n_{S}$ increases, the limit

$$
\begin{equation*}
n_{S}(\infty)=\lim _{t \rightarrow \infty} n_{S}(t) \leq \infty \tag{49}
\end{equation*}
$$

exists. As an application of value distribution theory on complex spaces, I was able to show that $S$ is affine algebraic if and only if $n_{S}(\infty)<\infty([63])$.

This result was localized by Errett Bishop [5] (1964) to extend analytic sets over higher dimensional analytic sets. His result was refined by Shiffman [47], [48], [49].

Hellmuth Kneser did not proceed to a Second Main Theorem and a Defect Relation. Also he did not consider the possible extension of his theory to parabolic manifolds or Kähler manifolds. However, he investigated another problem: the theory of functions of finite order. He solved the two dimensional case and provided the basic ideas in m -dimensions. Later he assigned the completion of these investigations to me as my thesis topic [62], [63].

Again let $(M, \tau)$ a strictly parabolic manifold of dimensions $m>1$. Thus $M$ is a hermitian vector space of dimension $m>1$ and $\tau$ is the square of the norm. If $\mathfrak{x} \in M, \mathfrak{y} \in M$, then $(\mathfrak{x} \mid \mathfrak{y})$ is the hermitian product of $\mathfrak{x}$ and $\mathfrak{y}$. If $u: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$is an increasing function, its order is defined by

$$
\begin{equation*}
0 \leq \operatorname{Ord} u=\underset{r \rightarrow \infty}{\limsup } \frac{\log u(r)}{\log r} \leq \infty \tag{50}
\end{equation*}
$$

If $\nu \geq 0$ is a non-negative divisor, define $\operatorname{Ord} \nu=\operatorname{Ord} n_{\nu}$. Then
$\operatorname{Ord} \nu=\operatorname{Ord} N_{\nu}(\cdot, s)$. If $f: M \rightarrow \mathbb{P}(V)$ is a meromorphic map, define $\operatorname{Ord} f=\operatorname{Ord} T_{f}(\cdot, s)$.

If $q$ is a non-negative integer, the Weierstrass prime factor is defined for all $z \in \mathbb{C}$ by

$$
\begin{equation*}
E(z, q)=(1-z) \exp \left(\sum_{p=1}^{q} \frac{1}{p} z^{p}\right) . \tag{51}
\end{equation*}
$$

For all $z \in \mathbb{C}(1)$ the Kneser Kernel is defined by

$$
\begin{equation*}
e_{m}(z, q)=\frac{1}{(m-1)!} \frac{d^{m-1}}{d z^{m-1}}\left(z^{m-1} \log E(z, q)\right) \tag{52}
\end{equation*}
$$

where $\log E(0, q)=0$.
Let $f: M \rightarrow \mathbb{C}$ be an entire function of finite order with $f(0)=$ 1. Let $S$ be the support of the zero divisor $\nu=\mu_{f}$ of $f$. Trivially $S=f^{-1}(0)$. Assume that $S \neq \emptyset$. Then there exists a largest real number $s>0$ such that $S(s)=0$. Since $f$ has finite order, there is a smallest, non-negative integer $q$ such that

$$
\begin{equation*}
\int_{s}^{\infty} \frac{T_{f}(r, s)}{r^{q+2}} d r<\infty \tag{53}
\end{equation*}
$$

Then $q \leq \operatorname{Ord} f \leq q+1$. Also there exists a holomorphic function $F$ on $M(s)$ such that $F(0)=0$ and $f \mid W(s)=e^{F}$. By the First Main Theorem the following integral converges uniformly on every compact subset of $M(s)$ and defines a holomorphic function $H$ on $M(s)$ with $\mu_{H}(0) \geq q+1$ by

$$
\begin{equation*}
H(\mathfrak{z})=\int_{\mathfrak{y} \in S} \nu(\mathfrak{y}) e_{m}\left(\frac{(\mathfrak{z} \mid \mathfrak{y})}{(\mathfrak{y} \mid \mathfrak{y})}, q\right) \omega^{m-1}(\mathfrak{y}) \tag{54}
\end{equation*}
$$

for $z \in M(s)$. Kneser [24] shows that there is a unique polynomial $P$ of at most degree $q$ with $P(0)=0$ such that

$$
\begin{equation*}
F=P|W(s)+H \quad f| W(s)=e^{P+H} \tag{55}
\end{equation*}
$$

Hence $h=e^{-P} f$ is an entire function with $\mu_{h}=\nu=\mu_{f}$ and
$h \mid W(s)=e^{H}$. Thus $h$ depends on $\nu$ only.
Given a divisor $\nu \geq 0$ on $M$ of finite order with $S=\operatorname{supp} \nu \neq 0$, there is a largest real number $s>0$ such that $S(s)=0$ and a smallest, non-negative integer $q$ such that

$$
\begin{equation*}
\int_{s}^{\infty} n_{\nu}(t) \frac{d t}{t^{q+2}}<\infty \tag{56}
\end{equation*}
$$

Then $q \leq \operatorname{Ord} \nu \leq q+1$. The integral (53) converges uniformly on every compact subset of $W(s)$ and defines a holomorphic function $H$ on $M(s)$ with $H(0)=0$ and $\mu_{H}(0) \geq q+1$ by (53). Does there exist an entire function $h$ on $M$ such that $h \mid W(s)=e^{H}$, such that $\mu_{h}=\nu$ and such that Ord $h=\operatorname{Ord} \nu$ ? In his earlier paper, Kneser [23] (1936) proved the existance of such a canonical function if $m=2$. It was my thesis problem to solve the case $m>2$. His method required to show that a certain closed form was exact. If $m=2$, this lead to a solvable ordinary differential equation. If $m>2$, it took me two weeks to write out the system of partial differential equations to be solved, which I could not do. I asked him for advice. He said he had gone through the same terrible calculation and had been unable to solve the system. Then he threw away his notes. I followed his advice, but I found another proof ([62], [63]). Independently, Pierre Lelong ([25] 1953, [27] 1964) proved the existence of the canonical function $h$ by another integral representation. Both solutions coincide by a uniquenen theorem of Rankin [42] (1968), who provided a third integral representation. In [64] 1953 I showed that the canonical function $h$ of a 2 m -periodic divisor is a theta function for this divisor and that any 2 m -periodic meromorphic function is a quotient of two theta functions (Appell [2] 1891 if $m=2$ and Poincaré [40] 1898 if $m \geq 2$ ). In 1975, Henri Skoda [58] and Gennadi Henkin [20] showed independently, that a non-negative divisor $\nu$ on a strictly pseudoconvex domain $D$ in $M$ with bounded valence $N_{\nu}$ is the zero divisor $\nu=\mu_{h}$ of a holomorphic function $h$ on $D$ with bounded characteristic. Later Henkin [21] (1978) showed, if $\operatorname{Ord} \nu<\infty$ then there is a holomorphic function $h$ on $D$ with $\nu=\mu_{h}$ and $\operatorname{Ord} \nu=\operatorname{Ord} h$. Recently, Polyakov [41] (1987) extended this result to the polydisc. Skoda [60] (1972) solved
the problem for analytic sets of higher codimension in a complex vector space. For more details see [73].

The integral means method of Kneser fails on complex manifolds. Also he did not attempt to prove a Second Main Theorem and a Defect Relation. From the theory of holomorphic curves there are available the method of Cartan [8] and the method of Ahlfors [1] which was extended to Riemann surfaces by Hermann and Joachim Weyl [90], improved later by $\mathrm{H} . \mathrm{Wu}$ [92].

In 1953/54 I extended the theory of Ahlfors-Weyl to meromorphic maps $f: M \rightarrow \mathbb{P}(V)$, where $M$ is a m-dimensional, connected, complex manifold of dimension $m>1$ endowed with a positive form $\chi$ of bidegree $(m-1, m+1)$ such that $d \chi=0$. Here $V$ is a hermitian vector space of dimension $n+1$. Again the targets are the hyperplanes in $\mathbb{P}(V)$ and $f$ is linearly non-degenerated. Let $\mathfrak{X}=\left\{a_{j}\right\}_{j \in Q}$ be a family of hyperplanes $a_{j} \in \mathbb{P}\left(V^{*}\right)$ in general position. Then, under suitable assumptions a defect relation

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq n+1 \tag{57}
\end{equation*}
$$

was obtained. Also a defect relation for associated maps was proved [65]. I cannot go into details here. The extension to $m>1$ is based on two ideas:
(1) Let $\mathbb{S}$ be a set of open, relative compact subsets $G$ of $M$ with $C^{\infty}$-boundary such that $\bar{g} \in G$ for all $G \in \mathscr{S}$, where $g$ is open with a $C^{\infty}$-boundary. Assume that for each compact subset $K$ of $M$ there is $G \in \mathbb{S}$ with $G \supset K$. There the Dirichlet problem $d d^{c} \Psi \wedge \chi=0$ is solved for $\bar{G}-g$ with $\Psi \mid \partial G=0$ and $\Psi \mid \partial \bar{g}=1$.
(2) The associated maps are defined by the use of a holomorphic differential form $B$ of bidegree ( $m-1,0$ ) such that

$$
\begin{equation*}
0 \leq m i_{m-1} B \wedge \bar{B} \leq Y(G) \chi \quad \text { on } \bar{G} \tag{58}
\end{equation*}
$$

where $Y(G)$ is the smallest possible constant.
On parabolic manifolds the proof has been greatly simplified by Cowen-Griffiths [17] (1976), Pit-Mann Wong [93] (1976), Stoll [80] (1983), [82] (1985), [86] (1992). The definitions and identities (34)-(46) also hold on parabolic manifolds except, of course, for the
slicing $j_{\mathfrak{b}}$ and the equality $n_{\nu}(0)=\nu(0)$ and (41) may be vacuous, since $f$ may not have a global, reduced representation on $M$. The defect of $(f, g)$ is defined as in (16). For an exact statement of the defect relation I refer to the papers mentioned before, but I will state the defect relation in a special case with a new variation:

Let $M$ be a connected, complex manifold of dimension $m>1$. Let $W$ be a hermitian vector space of dimension $m$. Let $\pi: M \rightarrow$ $W$ be a surjective, proper, holomorphic map. Then $\tau=\|\pi\|^{2}$ is a parabolic exhaustion of $M$ and $(M, \tau)$ is called a parabolic covering space of $W$. Take any holomorphic form $\zeta$ of bidegree $(m, 0)$ on $W$ without zeroes. Then the zero divisor $\beta \geq 0$ of $\pi^{*}(\zeta)$ does not depend on the choice of $\zeta$ and is called the branching divisor of $\pi$. Put $B=\operatorname{supp} \beta$. Then $\pi$ is locally biholomorphic at $z \in M$ if and only if $z \in M-B$. Since $\pi$ is proper and holomorphic, $B^{\prime}=\pi(B)$ and $\hat{B}=\pi^{-1}\left(B^{\prime}\right)$ are analytic and $\pi: M-\hat{B}=W-B^{\prime}$ is a covering space in the sense of topology. Its sheet number $\varsigma$ is given by (31).

Let $V$ be a hermitian vector space of dimension $n+1>1$. Let $f: M \rightarrow \mathbb{P}(V)$ be a linearly non-degenerated meromorphic map of transcendental growth (i.e. $A_{f}(\infty)=\infty$ ). Assume that the Ricci defect

$$
\begin{equation*}
R_{f}=\lim _{r \rightarrow \infty} \frac{N_{\beta}(r, s)}{T_{p}(r, s)}<\infty \tag{59}
\end{equation*}
$$

Let $\mathfrak{Z}=\left\{a_{j}\right\}_{j \in Q}$ a finite family of hyperplanes $a_{j} \in \mathbb{P}\left(V^{*}\right)$ in general position. Then we have the Defect Relation

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq n+1+\frac{1}{2} n(n+1) R_{f} . \tag{60}
\end{equation*}
$$

A meromorphic map $h: M \rightarrow \mathbb{P}(V)$ is said to separate the fibers of $\pi$, if there is a point $x \in W-B^{\prime}$ such that $\pi^{-1}(x) \wedge I_{h}=\emptyset$ and such that $h \mid \pi^{-1}(x)$ is injective. If so, and if $s>0$, there is a constant $C(s)>0$ such that

$$
\begin{equation*}
N_{\beta}(r, s) \leq 2(\varsigma-1) T_{h}(r, s)+C(s) \tag{61}
\end{equation*}
$$

for all $r>0$ (Noguchi [38], Stoll [83]). Define

$$
\begin{equation*}
\mathfrak{G}=\bigcup_{k \in \mathbb{N}}\left\{h \mid h: M \rightarrow \mathbb{P}_{k} \text { meromorphic, separates fibers of } \pi\right\} \tag{62}
\end{equation*}
$$

Then the separation index of $f$ is defined by

$$
\begin{equation*}
\gamma=\inf _{h \in \mathfrak{h}} \limsup _{r \rightarrow \infty} \frac{T_{h}(r, s)}{T_{f}(r, s)} \tag{63}
\end{equation*}
$$

If $f$ separates the fibers of $\pi$, then $\gamma \leq 1$. We obtain the Defect Relation

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq n+1+n(n+1)(\varsigma-1) \gamma \tag{64}
\end{equation*}
$$

If $n=1$, that is, if $f$ is a meromorphic function with transcendental growth separating the fibers of $\pi$, then

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, a_{j}\right) \leq 2 \varsigma \tag{65}
\end{equation*}
$$

which, in the case $m=1$, was already proved by H. Cartan [8] (1933).
In 1977 Al Vitter [89] proved the Lemma of the logarithmic derivative for meromorphic functions on a hermitian vector space $W$ and derived the defect relation for meromorphic maps $f: W \rightarrow \mathbb{P}(V)$ by Cartan's original method. For a detailed account see also Stoll [79], 1982. E. Bardis [3] (1990) extended the result to parabolic covering spaces of $W$.

In 1973-74, Carlson and Griffiths [16] and Griffiths and King [19] invented a new method to prove the defect relation. In keeping within [19], the result shall be stated only in the case of a parabolic covering space $(M, \tau)$ of a hermitian vector space of dimension $m>$ 1 . The advantage of the new method is, that it applies to holomorphic maps $f: M \rightarrow N$, where $N$ is a connected, n-dimensional, compact, complex manifold. A positive holomorphic line bundle $L$ spanned by its holomorphic sections is given on $N$. Then $N$ is projective algebraic. The disadvantage of the new method is, that we have to assume that the map $f$ is dominant which means that rank $f=n$. The vector space $Y^{*}$ of all holomorphic sections of $L$ have finite dimension $k+1>1$. If $0 \neq \mathfrak{a} \in Y^{*}$, the zero set $E_{L}[a]=\{x \in N \mid \mathfrak{a}(x)=0\}$ depends on
$a=\mathbb{P}(\mathfrak{a}) \in \mathbb{P}\left(Y^{*}\right)$ only. Let $Y=\left(Y^{*}\right)^{*}$ be the dual vector space of $Y$. If $x \in N$, the linear subspace $\Phi(x)=\left\{\mathfrak{a} \in V^{*} \mid \mathfrak{a}(x)=0\right\}$ has dimension $k$. Thus one and only one $\varphi(x) \in \mathbb{P}(Y)$ exists such that $E[\varphi(x)]=\mathbb{P}(\Phi(x))$. The holomorphic map $\varphi: N \rightarrow \mathbb{P}(Y)$ is called the dual classification map of $L$. The value distribution functions of $f$ are defined as those of $\varphi \circ f$. First Main Theorem holds but the defect relation so obtained is not optimal. As before we assume that $f$ has transcendental growth and that there is given a finite family $\mathfrak{A}=\left\{a_{j}\right\}_{j \in Q}$ of points $a_{j} \in \mathbb{P}\left(Y^{*}\right)$. However we have to consider the geometry of $\left\{E_{L}\left[a_{j}\right]\right\}_{j \in Q}$ and not the geometry of $\left\{E\left[a_{j}\right]\right\}_{j \in Q}$. Define

$$
\begin{equation*}
E_{L}[\mathscr{2}]=\bigcup_{j \in Q} E_{L}\left[a_{j}\right] . \tag{66}
\end{equation*}
$$

For each $j \in Q$ take $\mathfrak{a}_{j} \in V_{*}^{*}$ with $a_{j}=\mathbb{P}\left(\mathfrak{a}_{j}\right)$. Take $x \in E_{L}[\mathfrak{X}]$. Then

$$
\begin{equation*}
P=\left\{j \in Q \mid x \in E_{L}\left[a_{j}\right]\right\}=\left\{j \in Q \mid \mathfrak{a}_{j}(x)=0\right\} \neq \emptyset \tag{67}
\end{equation*}
$$

Put $p=\# P$. Take a bijective map $\lambda: \mathbb{N}[1, p] \rightarrow P$. There is an open, connected neighborhood $U$ of $x$ and a holomorphic section $\mathfrak{b}: U \rightarrow L$ such that $\mathfrak{b}(z) \neq 0$ for all $z \in U$. For each $j \in \mathbb{N}[1, p]$, there is one and only one holomorphic function $h_{j}$ on $U$ such that $\mathfrak{a}_{\lambda(j)} \mid U=h_{j} \mathfrak{b}$. Then $\mathfrak{2 l}$ is said to have strictly normal crossings at $x$ if and only if

$$
\begin{equation*}
d h_{1}(x) \wedge \ldots \wedge d h_{p}(x) \neq 0 \tag{68}
\end{equation*}
$$

The definition is independent of the choices which were made. $\mathfrak{W}$ is said to have strictly normal crossings if $\mathfrak{Q}$ has strictly normal crossings at every $x \in E_{L}[\mathfrak{2}]$, which we assume now.

Let $K$ be the canonical bundle of $N$ Let $K^{*}$ be the dual bundle to $K$. Define

$$
\begin{equation*}
\left[\frac{K^{*}}{L}\right]=\inf \left\{\left.\frac{v}{w} \right\rvert\, v \in \mathbb{N}, w \in \mathbb{N}, L^{v} \otimes K^{w} \text { positive }\right\} \tag{69}
\end{equation*}
$$

Define $R_{f}$ by (59) and $\gamma$ by (63). With these assumptions and definitions, the

## Defect Relation of Griffiths-King

$$
\begin{align*}
& \sum_{j \in Q} \delta\left(f, a_{j}\right) \leq\left[\frac{K^{*}}{L}\right]+R_{f}  \tag{70}\\
& \sum_{j \in Q} \delta\left(f, a_{j}\right) \leq\left[\frac{K^{*}}{L}\right]+2(\varsigma-1) \gamma \tag{71}
\end{align*}
$$

holds. In [75] (1977) the theory was refined and extended to general parabolic manifolds.

A difficult, major, unsolved problem is the question if "dominant" can be replaced by another assumption which does not imply $m \geq n$. For instance does (70) hold if $f(M)$ is not contained in any proper analytic subset of $N$ ? As Biancofiore [4] has shown the assumption $f(M) \nsubseteq E_{L}[a]$ for all $a \in \mathbb{P}\left(Y^{*}\right)$ does not suffice. Can the condition "strictly normal crossings" be relaxed?

Let $V$ be a hermitian vector space of dimension $n+1>1$. Apply the previous theory to $N=\mathbb{P}(V)$. Let $H$ be the hyperplane section bundle on $\mathbb{P}(V)$. Take $p \in \mathbb{N}$ and choose $L=H^{p}$. Then $K=H^{-n-1}$ and $L^{v} \otimes K^{w}=H^{p v-w(n+1)}$. Thus $\left[\frac{K^{*}}{L}\right]=\frac{n+1}{p}$. Thus (70) and (71) reads

$$
\begin{align*}
& \sum_{j \in Q} \delta\left(f, a_{j}\right) \leq \frac{n+1}{p}+R_{f}  \tag{72}\\
& \sum_{j \in Q} \delta\left(f, a_{j}\right) \leq \frac{n+1}{p}+2(\varsigma-1) \gamma . \tag{73}
\end{align*}
$$

If $p=1$, this is sharper than (60) which is due to the dominance of $f$.

Until now, target families of codimension 1 only where considered. Does there exist a value distribution theory for codimension $\ell>1$. In 1958, H. Levine [30] proved an unintegrated First Main Theorem for projective planes of codimension $\ell>1$ in $\mathbb{P}(V)$. At the 1958 Summer School at the University of Chicago, S. S. Chern asked me to find the integrated version. When I left, I told him that there is no such thing. I was much surprised when he published an integrated version [10] (1960) shortly afterwards. I failed, since I insisted on an
old version to be obtained and because I had forgotten one of Max Planck's admonitions in one of his textbooks: "The energy principle is not a law of nature, but of man. Each time it fails in nature, man invents a new type of energy to restore the principle." The First Main Theorem is such a principle. In order to retain it, S. S. Chern had to admit a new, nasty term, later called the deficit, into the equation.

In 1965, Bott and Chern [6] extended the First Main Theorem to the equidistribution of the zeroes of holomorphic sections in hermitian vector bundles. Thus differential geometry was brought into value distribution theory. Later the theory was expanded to include all Schubert varieties associated to holomorphic vector bundles. With the work of H. Wu [91] (1968-70), F. Hirschfelder [22] (1969), L. Dektjarev [18] (1970), Michael Cowen [16] (1973), Chia-Chi Tung [88] (1973), and myself [67] (1967) [68] (1969) [69] (1970) and [76] (1978) a wide range of First Main Theorems for codimension $\ell>1$ was established.

Mostly, they can be brought under the following scheme


Where $M, N$ and $E$ are connected, complex manifolds of dimensions $m, n$ and $k$ respectively. Here $E$ is a compact Kähler manifold and $S$ is an analytic subset of $N \times E$. The projections $\varrho$ and $\pi$ are surjective, open and of pure fiber dimensions $q$ and $p$ respectively with $n-p=$ $\ell \geq 1$ and $m-\ell \geq 0$. The map $\varrho$ is locally a product at every point of $S$. Since $E$ is compact, $\varrho$ is proper. Thus

$$
\begin{equation*}
\operatorname{dim} S=p+k=n+q \quad k-q=n-p=\ell . \tag{75}
\end{equation*}
$$

The diagram is completed as a pull back by the holomorphic map $f$ :

$$
\begin{equation*}
Q=\{(x, z) \mid f(x)=\varrho(z)\} \tag{76}
\end{equation*}
$$

$$
\begin{array}{lc}
\tilde{\varrho}(x, z)=x & \tilde{f}(x, z)=1 \quad \hat{f}(x, z)=\pi(z) \\
\varrho \circ \tilde{f}=f \circ \tilde{\varrho} & \hat{f}=\tilde{f} \circ \pi . \tag{78}
\end{array}
$$

The map $\tilde{\varrho}$ has pure fiber dimensions $q$ and is locally a product at every point of $Q$. Hence $Q$ has pure dimension $m+q$.

For each $y \in E$, the analytic subset $S_{y}=\varrho\left(\pi^{-1}(y)\right)$ is a pure p-dimensional analytic subset of $N$. The family $\mathbb{S}=\left\{S_{y}\right\}_{y \in E}$ is the target family for the holomorphic map $f$. We assume that $E_{y}=$ $f^{-1}\left(S_{y}\right)$ is either empty or has generically the dimension $m-\ell$. Let $\xi>0$ be the Kähler volume for of $E$ with

$$
\begin{equation*}
\int_{E} \xi=1 \tag{79}
\end{equation*}
$$

Let $\varrho_{*}$ be the fiber integration operator. Then $\Omega=\varrho_{*} \pi^{*}(\xi)$ is a nonnegative closed form of bidegree ( $\ell, \ell$ ) and class $C^{\infty}$ on $N$. Here $\Omega$ is the Poincaré dual of the homology class defined by $\mathbb{S}$. Take $y \in E$, by Hodge theory or construction (H. Wu [91], Stoll [69]) there is a non-negative form $\lambda_{y} \geq 0$ on $E-\{y\}$ of bidegree $(k-1, k-1)$ with residue 1 at $y$ such that

$$
\begin{equation*}
d d^{c} \lambda_{y}=\xi \quad \text { on } E-\{y\} \tag{80}
\end{equation*}
$$

Then $\Lambda_{y}=\varrho_{*} \pi^{*}\left(\lambda_{y}\right) \geq 0$ is a form of bidegree $(\ell-1, \ell-1)$ on $N-S_{y}$ with

$$
\begin{equation*}
d d^{c} \Lambda_{y}=\Omega \quad \text { on } N-S_{y} \tag{81}
\end{equation*}
$$

Let $\varphi$ be a form of bidegree ( $m-\ell, m-\ell$ ) and of class $C^{\infty}$ with compact support in $M$. With proper multiplicities $\nu_{y}$, the Stokes Theorem, the Residue Theorem and fiber integration imply

$$
\begin{align*}
\int_{M} f^{*}\left(\Lambda_{y}\right) \wedge d d^{c} \varphi & =-\int_{M} d f^{*}\left(\Lambda_{y}\right) \wedge d^{c} \varphi  \tag{82}\\
& =-\int_{M} d \varphi \wedge d^{c} f^{*}\left(\Lambda_{y}\right)
\end{align*}
$$

$$
=-\int_{F_{y}} \nu_{y} \varphi+\int_{M} \varphi \wedge d d^{c} f^{*}\left(\Lambda_{y}\right),
$$

if $E_{y}$ has pure dimension $m-\ell$. As a generalization of the PoincareLelong formula we obtain the Unintegrated First Main Theorem

$$
\begin{equation*}
\int_{M} f^{*}(\Omega) \wedge \varphi=\int_{F_{y}} \nu_{y} \varphi+\int_{M} f^{*}\left(\Lambda_{y}\right) \wedge d d^{c} \varphi . \tag{83}
\end{equation*}
$$

For the integration, we assume that an exhaustion $\tau: M \rightarrow \mathbb{R}_{+}$is given with
(84) $w=d d^{c} \log \tau \geq 0 \quad v=d d^{c} \tau \geq 0 \quad \sigma_{\ell}=d^{c} \log \tau \wedge w^{m-\ell}$,

Then $d \sigma_{\ell}=w^{m-\ell+1}$. We keep the notations (27) and (28), but do not require that $\tau$ is parabolic. For $t>0$ the spherical image function is defined by

$$
\begin{equation*}
A_{f}(t)=\frac{1}{t^{2 m-2 \ell}} \int_{M[t]} f^{*}(\Omega) \wedge v^{m-\ell} \geq 0 \tag{85}
\end{equation*}
$$

For $0<s<r$ the characteristic function is defined by

$$
\begin{equation*}
T_{f}(r, s)=\int_{s}^{r} A_{f}(t) \frac{d t}{t} \geq 0 \tag{89}
\end{equation*}
$$

Take $y \in E$ such that $E_{y}$ has pure codimension $\ell$ or is empty. For all $t>0$ the counting function is defined by

$$
\begin{equation*}
n_{f, y}(t)=\frac{1}{t^{2 m-2 \ell}} \int_{E_{y}[t]} \nu_{y} v^{m-\ell} \geq 0 \tag{90}
\end{equation*}
$$

and for $0<s<r$ the valence function is defined by

$$
\begin{equation*}
N_{f, y}(r, s)=\int_{s}^{r} n_{f, y}(t) \frac{d t}{t} \geq 0 \tag{91}
\end{equation*}
$$

For almost all $r>0$ the compensation function is defined by

$$
\begin{equation*}
m_{f, y}(r)=\frac{1}{2} \int_{M<r>} f^{*}\left(\Lambda_{y}\right) \wedge \sigma_{\ell} \geq 0 \tag{92}
\end{equation*}
$$

For $0<s<r$ the deficit is defined by

$$
\begin{equation*}
D_{f, y}(r, s)=\frac{1}{2} \int_{M[r]-M[s]} f^{*}\left(\Lambda_{y}\right) \wedge \omega^{m-\ell+1} . \tag{93}
\end{equation*}
$$

If $\ell=1$ and $\tau$ is parabolic, then $\omega^{m} \equiv 0$ which implies $D_{f, y} \equiv 0$. However if $\ell>1$, then this is false even if $\tau$ is parabolic. The same calculation as in (82) but respecting boundary terms yields the First Main Theorem

$$
\begin{equation*}
T_{f}(r, s)=N_{f, y}(r, s)+m_{f, y}(r)-m_{f, y}(s)-D_{f, y}(r, s) \tag{94}
\end{equation*}
$$

A continuous form $\hat{\lambda} \geq 0$ bidegree $(k-1, k-1)$ on $E$ exists such that $x \in E$ implies

$$
\begin{equation*}
\hat{\lambda}(x)=\int_{y \in E} \lambda_{y}(x) \otimes \xi(y) \geq 0 \tag{95}
\end{equation*}
$$

The $\hat{\Lambda}=\varphi_{*} \pi^{*}(\hat{\lambda}) \geq 0$ is a continuous form of bidegree $(\ell-1, \ell-1)$ on $N$. For all $x \in N$, fiber integration yields

$$
\begin{equation*}
\hat{\Lambda}(z)=\int_{y \in E} \Lambda_{y}(z) \otimes \xi(y) \geq 0 \tag{96}
\end{equation*}
$$

Thus we obtain
(97) $\quad \mu_{f}(r)=\int_{y \in E} m_{f, y}(r) \xi(y)=\frac{1}{2} \int_{M<r>} f^{*}(\hat{\Lambda}) \wedge \sigma_{\ell} \geq 0$
(98) $\Delta_{f}(r, s)=\int_{y \in E} D_{f, y}(r, s) \xi(y)=\frac{1}{2} \int_{M[r]-M[s]} f^{*}(\hat{\Lambda}) \wedge \omega^{m-\ell-1} \geq 0$

$$
\begin{equation*}
T_{f}(r, s)=\int_{y \in E} N_{f, y}(r, s) \xi(y) \geq 0 \tag{99}
\end{equation*}
$$

which implies

$$
\begin{equation*}
\Delta_{f}(r, s)=\mu_{f}(r)-\mu_{f}(s) . \tag{100}
\end{equation*}
$$

For $r>0$ define

$$
\begin{align*}
B(r) & =\left\{y \in E \mid E_{y} \cap M[r] \neq \emptyset\right\}  \tag{101}\\
0 \leq b_{f}(r) & =\int_{B(r)} \xi \leq 1 \\
B & =\left\{y \in E \mid E_{y} \neq \emptyset\right\} \\
0 \leq b_{f} & =\int_{B} \xi \leq 1 .
\end{align*}
$$

Then $B=\bigcup_{r>0} B(r)$ and $b_{f}(r) \rightarrow b_{f}$ for $r \rightarrow \infty$ increasingly. Now (94) implies

$$
\begin{equation*}
N_{f, y}(r, s) \leq T_{f}(r, s)+m_{f, y}(s)+D_{f, y}(r, s) . \tag{103}
\end{equation*}
$$

If $y \in E-B(r)$, then $N_{f, y}(r, s)=0$ and (99) implies

$$
\begin{align*}
T_{f}(r, s) & =\int_{y \in E} N_{f, y}(r, s) \xi(y)=\int_{y \in B(r)} N_{f, y}(r, s) \xi(y) \\
& \leq \int_{y \in B(r)}\left(T_{f}(r, s)+m_{f, y}(s)+D_{f, y}(r, s)\right) \xi(y)  \tag{104}\\
& \leq b_{f}(r) T_{f}(r, s)+\int_{y \in E}\left(m_{f, y}(s)+D_{f, y}(r, s)\right) \xi(y) \\
& =b_{f}(r) T_{f}(r, s)+\mu_{f}(s)+\Delta_{f}(r, s) .
\end{align*}
$$

Therefore
(105) $0 \leq\left(1-b_{f}(r)\right) \leq \frac{\mu_{f}(s)+\Delta_{f}(r, s)}{T_{f}(r, s)} \quad$ if $r>s>0$.

Assume that $T_{f}(r, s) \rightarrow \infty$ and $\Delta_{f}(r, s) / T_{f}(r, s) \rightarrow 0$ for $r \rightarrow \infty$. Then $b_{f}=1$. Thus $f(M)$ intersects almost all targets $S_{y}$. Even for holomorphic curves on $\mathbb{C}$ surprising results can be obtained:

## Proposition

A holomorphic map $f: \mathbb{C} \rightarrow \mathbb{P}_{6}$ is defined for all $z \in \mathbb{C}$ by

$$
\begin{equation*}
f(z)=\mathbb{P}\left(1, e^{\zeta z}, e^{\zeta^{2} z}, \ldots, e^{\zeta^{6} z}\right) \tag{106}
\end{equation*}
$$

where $\zeta=e^{\frac{\pi i}{3}}$. If $r \geq \frac{\pi}{6}(245+\log 7) \approx 129.3006$, then $f(\mathbb{C}[r])$ intersects at least $99 \%$ of all hyperplanes in $\mathbb{P}_{6}$.

Proof. A reduced representation $\mathfrak{b}$ of $f$ is defined for all $z \in \mathbb{C}$ by

$$
\mathfrak{v}(z)=\left(1, e^{\zeta z}, e^{\zeta^{2} z}, \ldots, e^{\zeta^{6} z}\right)
$$

with $\mathfrak{v}(0)=(1, \cdots, 1)$. Thus $\|\mathfrak{v}(0)\|=\sqrt{7}$. We can take $s=0$. Thus

$$
T_{f}(r, 0)=\int_{\mathbb{C}<r>} \log \|\mathfrak{b}\| \sigma-\frac{1}{2} \log 7
$$

Observe that

$$
L=\sum_{j=1}^{6}\left|\zeta^{j}-\zeta^{j-1}\right|=6
$$

By Stoll [80] Proposition 15.5 page 201 we have

$$
0 \leq \int_{\mathbb{C}\langle r\rangle} \log \|\mathfrak{v}\| \sigma-\frac{L}{2 \pi} r \leq \frac{1}{2} \log 7
$$

Thus

$$
\frac{6 r-\pi \log 7}{2 \pi} \leq T_{f}(r, 0)
$$

By Stoll [80] (6.66) page 140 we have $\mu_{f}(s)=\frac{1}{2} \sum_{p=1}^{6} \frac{1}{p}=\frac{49}{40}$ for all $s>0$. If $r>(\pi / 6) \log 7$, then

$$
0 \leq 1-b_{f}(r) \leq \frac{49}{40} \frac{2 \pi}{6 r-\pi \log 7}=\frac{49}{20} \frac{\pi}{6 r-\pi \log 7}
$$

Define $r_{0}=\frac{\pi}{6}(245+\log 7)$. Take $r \geq r_{0}$ Then

$$
0 \leq 1-b_{f}(r) \leq \frac{49}{20} \frac{\pi}{6 r_{0}-\pi \log 7}=\frac{49}{20 \times 245}=\frac{1}{100}
$$

Hence $b_{f}(r) \geq \frac{99}{100}$, q.e.d.
This calculation was made possible by a theorem of ShiffmanWeyl. The method can be greatly improved, see Molzon, Shiffman, and Sibony [31] (1981), and Lelong and Gruman [29] (1986).

In 1929 Rolf Nevanlinna [33] conjectured that his defect relation remains valid, if the constant target points $a_{j} \in \mathbb{P}_{1}$ are replaced by "target" functions $g_{j}: \mathbb{C} \rightarrow \mathbb{P}_{1}$ which move slower than the "hunter" function $f: \mathbb{C} \rightarrow \mathbb{P}_{1}$, that is, if

$$
\begin{equation*}
T_{g_{j}}(r, s) / T_{f}(r, s) \rightarrow 0 \quad \text { for } r \rightarrow \infty \tag{107}
\end{equation*}
$$

In 1964 Chi-Tai Chuang [14] proved the conjecture for entire functions $f: \mathbb{C} \rightarrow \mathbb{C}$ and created the basis for the solution of the problem. In 1986, Norbet Steinmetz [61] proved Nevanlinna's conjecture. In 1991, Ru Min and I [43] [44] [85] proved the conjecture for holomorphic curves and solved the case of the Cartan conjecture for moving targets [46]. In 1985, Charles F. Osgood [39] claimed that these theorems are a consequence of his results in diophantine approximation, but to me this implication is not self evident and still has to be established.

At the end let me state a result at Notre Dame on this subject matter, combining the work of Emmanuel Bardis [3], and Ru Min and myself [44].

At first some concepts have to be explained. Let $M$ be a connected, complex manifold of dimension $m$. Let $V$ be a hermitian vector space of finite dimension $n+1>1$. Let $f: M \rightarrow \mathbb{P}(V)$ be a meromorphic map. Take $\mathfrak{a} \in V^{*}$ and $0 \neq \mathfrak{b} \in V^{*}$. Put $b=\mathbb{P}(\mathfrak{b})$. Assume that $(f, b)$ is free. Then there exists one and only one meromorphic function $f_{a, b}$ on $M$, called a coordinate function, such that for each point $p \in M$ there exists an open, connected neighborhood $U$ of $p$ and a reduced representation $\mathfrak{v}: U \rightarrow V$ such that

$$
\begin{equation*}
f \left\lvert\, U=\frac{\langle\mathfrak{b}, \mathfrak{a}\rangle}{\langle\mathfrak{b}, \mathfrak{b}\rangle} .\right. \tag{108}
\end{equation*}
$$

Here $\langle\mathfrak{b}, \mathfrak{b}\rangle \not \equiv 0$ since $(f, b)$ is free. Let $\mathfrak{C}_{f}$ be the set of all those coordinate functions of $f$. Trivially $\mathbb{C} \subseteq \mathfrak{C}_{f}$. Let $\mathfrak{M}$ be the field of meromorphic functions on $M$. Let $\Re$ be a subfield of $M$. The $f$ is said to be defined over $\mathfrak{\Re}$ if and only if $\mathfrak{C}_{f} \subseteq \mathfrak{\Re}$. The meromorphic map $f$ is said to be linearly non-degenerated over $\Re$ if and only if $(f, g)$ is free for every meromorphic map $g: M \rightarrow \mathbb{P}\left(V^{*}\right)$ defined over $\Re$. Let $\mathbb{S}=\left\{g_{j}\right\}_{j \in Q}$ be a finite family of meromorphic maps $g_{j}: M \rightarrow \mathbb{P}\left(V^{*}\right)$ with indeterminacy $I_{g_{j}}$. Define

$$
\begin{equation*}
I_{\mathfrak{G}}=\bigcup_{j \in Q} I_{g_{j}} \quad \mathfrak{C}_{\mathfrak{G}}=\bigcup_{j \in Q} \mathfrak{C}_{g_{j}} . \tag{109}
\end{equation*}
$$

Let $\Re_{\mathscr{G}}=\mathbb{C}\left(\mathfrak{C}_{\mathfrak{G}}\right)$ be the extension field of $\mathfrak{C}$ in $M$ generated by $\mathfrak{C}_{\mathbb{G}}$. The family $\mathbb{E S}_{5}$ is said to be in general position if and only if there is a point $z \in M-I_{\circledast}$ such that $\left(\mathbb{S}(z)=\left\{g_{j}(z)\right\}_{j \in Q}\right.$ is in general position.

## Theorem: Defect relation for moving target.

Let $M$ be a connected, complex manifold of dimension M. Let $W$ be a hermitian vector space of dimension $m$. Let $\pi: M \rightarrow W$ be a surjective, proper holomorphic map. Then $\tau=\|\pi\|^{2}$ is a parabolic exhaustion of $M$. Let $V$ be a hermitian vector space of finite dimension $n+1>1$. Let $\mathbb{G} S=\left\{g_{j}\right\}_{j \in Q}$ be a finite family of meromorphic maps $g: M \rightarrow \mathbb{P}\left(V^{*}\right)$ in general position. Assume at least on $k \in Q$ exists such that $g_{k}$ is not constant and separates the fibers of $\pi$. Let $f: M \rightarrow$ $\mathbb{P}(V)$ be a meromorphic map which is linearly non-degenerated over $\Re_{\mathfrak{G}}$. Assume that $g_{j}$ grows slower than for each $j \in Q$. Then

$$
\begin{equation*}
\sum_{j \in Q} \delta\left(f, g_{j}\right) \leq n+1 \tag{110}
\end{equation*}
$$

During the time from 1933 to 1960 the foundation was laid. The $1960^{\text {th }}$ was the decade of the First Main Theorem. The $1970^{\text {th }}$ was the decade of the Second Main Theorem. The $1980^{\text {th }}$ was the decade of the moving targets. Perhaps the $1990^{\text {th }}$ will be a decade of refinement and of
value distribution over function fields in conjunction with diophantine approximation.

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