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## TEMPEREDNESS OF MEASURES DEFINED BY POLYNOMIAL EQUATIONS OVER LOCAL FIELDS

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Dedicated to the memory of Professor Jun-Ichi Igusa


#### Abstract

We investigate the asymptotic growth of the canonical measures on the fibers of morphisms between vector spaces over local fields of arbitrary characteristic. For a single polynomial over $\mathbb{R}$, this is due to Igusa and Raghavan. For nonarchimedean local fields we use a version of the Lojasiewicz inequality which follows from work of Greenberg, together with the theory of the Brauer group of local fields to construct definite forms of arbitrarily high degree, and to transfer questions at infinity to questions near the origin. We then use these to generalize results of Hörmander on estimating the growth of polynomials at infinity in terms of the distance to their zero loci. Specifically, when a fiber corresponds to a noncritical value which is stable, i.e., remains noncritical under small perturbations, we show that the canonical measure on the fiber is tempered, which generalizes results of Igusa and Raghavan, and Virtanen and Weisbart.


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## 1. Introduction

Let $V$ be a finite dimensional vector space over $\mathbb{R}$ and $f: V \rightarrow \mathbb{R}$ a smooth nonconstant function. In the physics and mathematics literature the measure denoted by $\delta(f-c)$ figures prominently; it is a measure living on the smooth part of the zero locus $Z(f-c)$ of $f-c, c \in \mathbb{R}$ [Gel'fand and Shilov 1964]. Given $f$ and choices of Haar measures on $V$ and $\mathbb{R}, \delta(f-c)$ is uniquely defined for all $c$. Similarly if $\boldsymbol{f}=\left(f_{1}, f_{2}, \ldots, f_{r}\right): V \rightarrow \mathbb{R}^{r}$ is a smooth map with $d f_{1} \wedge \cdots \wedge d f_{r} \neq 0$, for given Haar measures on $V$ and $\mathbb{R}^{r}$, there is a canonical measure on the smooth part of the common zero locus $Z(\boldsymbol{f}-\boldsymbol{c})=Z\left(f_{1}-c_{1}, f_{2}-c_{2}, \ldots, f_{r}-c_{r}\right)$ of the $f_{i}-c_{i}$ for all $\boldsymbol{c}=\left(c_{1}, c_{2}, \ldots, c_{r}\right)$. We denote this measure by $\mu_{f, \boldsymbol{c}}$. In this context, the finiteness of $\mu_{f, c}$ around the singular points of $Z\left(f_{1}-c_{1}, f_{2}-c_{2}, \ldots, f_{r}-c_{r}\right)$, as well as the behavior at infinity of the extended measure, viewed as a Borel measure on $V$, are interesting questions. If the $f_{i}$ are polynomials and $Z(\boldsymbol{f}-\boldsymbol{c})$ is smooth, then it is natural to expect that $\mu_{f, c}$ is tempered. That is,

Definition 1.1 (tempered measure). Let $V$ be any finite dimensional $k$-vector space, $k$ a local field. A Borel measure $\mu$ on $V$ is tempered if

$$
\int_{V}\left(1+\|x\|^{2}\right)^{-\alpha} d \mu(x)<\infty
$$

for some integer $\alpha$ (in any norm).
This is equivalent to saying that there are constants $A>0, b \geq 0$ such that

$$
\begin{equation*}
\mu\left(B_{R}\right) \leq A R^{b} \tag{G}
\end{equation*}
$$

for all $R \geq 1, B_{R}$ being the closed ball in $V$ of radius $R$ and center $\mathbf{0}$ (in any norm).
In [Igusa 1978] Igusa and Raghavan proved that if $k=\mathbb{R}$ and $f$ is a nonconstant polynomial on $V$ and $c \in \mathbb{R}$ is a noncritical value of $f$, i.e., the locus $Z(f-c)$ is smooth, then $\mu_{f, c}$ is tempered, and further that the growth estimate G for the measure is uniform in a neighborhood of $c$; here we must remember that by the algebraic Sard's theorem (Proposition 2.4), $f$ has only finitely many critical values, so that every noncritical value $c$ has neighborhoods consisting only of noncritical values.

The measures $\mu_{f, c}, \mu_{f, c}$ can be defined over any local field. Throughout this paper by local field we mean a locally compact nondiscrete field of any characteristic, other than $\mathbb{C}$; measure theoretic questions over $\mathbb{C}$ usually reduce to $\mathbb{R}$, and so we do not treat the case of $\mathbb{C}$ separately. In [Igusa 1978] Igusa and Raghavan define the measures $\mu_{f, c}$ for any local field but do not consider their behavior at infinity, the reason being that over a nonarchimedean field they were concerned only with integrating Schwartz-Bruhat functions (i.e., compactly supported complex-valued locally constant functions). However the work of Harish-Chandra [1973] shows the necessity as well as utility of working with locally constant functions that do not
vanish outside a compact set. The question of extending the results of [Igusa 1978] to the nonarchimedean case and for $r>1$ is certainly a natural one. In [Virtanen and Weisbart 2014] the measures $\mu_{f, c}$ were shown to be tempered when $f$ is a nondegenerate quadratic form and $c \neq 0$; moreover for the case $c=0$ the locus $Z(f)$ has 0 as its only singularity, and it was shown that the measure $\mu_{f, 0}$ is finite in the neighborhood of 0 if $\operatorname{dim} V \geq 3$, and the extended measure is tempered in $V$. The work of [Virtanen and Weisbart 2014] was motivated by physical questions arising in the theory of elementary particles over $p$-adic spacetimes. In this paper we generalize the results of [Igusa 1978] and [Virtanen and Weisbart 2014] to the measures $\mu_{f, \boldsymbol{c}}$ where the $f_{i}(1 \leq i \leq r)$ are polynomials on a vector space $V$ over a local field $k$, with $\operatorname{dim}(V)=m$ and $d f_{1} \wedge d f_{2} \wedge \cdots \wedge d f_{r} \not \equiv 0$, so that $m \geq r$. Note that for $r>1$ and $k=\mathbb{R}$ this question is already more general than the one treated in [Igusa 1978].

We now describe our main result using the above notation. Let $f: V \rightarrow k^{r}$ be the polynomial map whose components are the $f_{i}$, with $d f_{1} \wedge \cdots \wedge d f_{r} \not \equiv 0$. A point $x \in V$ is called a critical point $(\mathrm{CP})$ of $\boldsymbol{f}$ if the differentials $d f_{i, x}$ are linearly dependent. We write $C(f)$ for the set of critical points of $f$; the image $f(C(f))$ in $k^{r}$ is called the set of critical values of $\boldsymbol{f}$, and is denoted by $C V(\boldsymbol{f})$. By the algebraic Sard's theorem (Proposition 2.4) one knows that in characteristic zero the Zariski closure in $k^{r}$ of $C V(\boldsymbol{f})$ is a proper algebraic subset of $k^{r}$. A point $c \in k^{r}$ is called stably noncritical if it has an open neighborhood (in the $k$-topology) consisting only of noncritical values. This is the same as saying that the fibers above points sufficiently close to $c$ are smooth. If $k$ has characteristic zero, stably noncritical values exist and form a nonempty open set in $k^{r}$ whose complement in the image of $f$ has measure 0 . Then the following is our main result. For $r=1$ and $k=\mathbb{R}$ it was proved in [Igusa 1978]. Note that in this case the characteristic is 0 and there are only finitely many critical values and so every noncritical value is stably noncritical.

Theorem 1.2. Fix $\boldsymbol{f}$ and write $\mu_{\boldsymbol{c}}=\mu_{\boldsymbol{f}, \boldsymbol{c}}$. Suppose $\boldsymbol{c}$ is stably noncritical. Then $\mu_{c}$ is tempered and there are constants $A>0, \gamma \geq 0$ such that for all $\boldsymbol{d}$ in an open neighborhood of $\boldsymbol{c}$

$$
\mu_{\boldsymbol{d}}\left(B_{R}\right) \leq A R^{m-r+\gamma} \quad(R \geq 1, \boldsymbol{d} \in U)
$$

Suppose $k$ has characteristic 0; then stably noncritical values form a nonempty dense open set whose complement in the image of $\boldsymbol{f}$ has measure 0 ; for $r=1$, the critical set is finite and all noncritical values are stably noncritical.

Remark 1.3. In view of the failure of Sard's theorem over characteristic $p>0$ (see page 233), we do not know if stably noncritical values of $\boldsymbol{c}$ always exist when $k$ is a local field of positive characteristic.

Remark 1.4. The results and ideas in the paper lie at the interface of analysis of geometry over local fields and are motivated by the themes from quantum theory over $p$-adic spacetimes. We do not know what, if any, are the arithmetic consequences of our results.

As an application of our theory we prove that if $k$ has characteristic 0 , the orbits of regular semisimple elements of a semisimple Lie algebra over $k$ are closed, and the invariant measures on them are tempered. For $k=\mathbb{R}$ this is a result of Harish-Chandra [1957].

## 2. Canonical measures on level sets of polynomial maps

Canonical measures on the fibers of submersive maps. The construction below is well known and our treatment is a very mild variant of Harish-Chandra's [1964] for the case $k=\mathbb{R}$ (see also [Varadarajan 1977]). Serre's book [2006] is a good reference for the theory of analytic manifolds and maps over a local field of arbitrary characteristic. (All of our manifolds are second countable.)

Lemma 2.1. Let $V, W$ be vector spaces of finite dimension $m, r$ respectively, and $L: V \rightarrow W$ be a surjective linear map. Let $U=\operatorname{ker} L$. Let $\sigma, \tau$ be exterior forms on $V, W$ of degrees $m, r$ respectively, with $\tau \neq 0$. Then there exists a unique exterior $(m-r)$-form $\rho$ on $U$ such that if $\left\{u_{1}, u_{2}, \ldots, u_{m-r}\right\}$ is a basis for $U$, then

$$
\rho\left(u_{1}, u_{2}, \ldots, u_{m-r}\right)=\frac{\sigma\left(u_{1}, \ldots, u_{m-r}, v_{1}, \ldots, v_{r}\right)}{\tau\left(L v_{1}, \ldots, L v_{r}\right)}
$$

where $v_{i} \in V$ are such that $\left\{u_{1}, \ldots, u_{m-r}, v_{1}, \ldots, v_{r}\right\}$ is a basis for $V$.
Proof. For fixed $v_{i}$ it is obvious that this defines an exterior $(m-r)$-form on $U$. Its independence of the choice of the $v_{i}$ is easy to check.

We write $\rho=\sigma / \tau$. Note that this definition is relative to $L$.
Theorem 2.2. Let $k$ be a local field of arbitrary characteristic and $M, N$ be analytic manifolds over $k$ of dimensions $m, r$ respectively, and $\pi: M \rightarrow N$ be an analytic map, surjective, and submersive everywhere. Let $\sigma_{M}$ (resp. $\tau_{N}$ ) be an analytic exterior $m$-form (resp. $r$-form) on $M$ (resp. $N$ ), with $\tau_{N} \neq 0$ everywhere on $N$. Then there is a unique analytic exterior form $\rho:=\rho_{M / N}$ on $M$ such that for any $y \in N$, the pull back of $\rho$ to the fiber $\pi^{-1}(y)$ is the exterior $(m-r)$-form $x \mapsto \sigma_{x} / \tau_{y}$ relative to $d \pi_{x}: T_{x}(M) \rightarrow T_{y}(N)$.
Proof. The pointwise definition of $\rho$ is clear after the preceding lemma. For analyticity we use local coordinates around $x$ and $y=\pi(x)$, say $x_{1}, \ldots, x_{m}$, such that $\pi$ is the projection $\left(x_{1}, \ldots, x_{m}\right) \mapsto\left(x_{1}, \ldots, x_{r}\right)$. Then

$$
\sigma_{M}=s\left(x_{1}, \ldots, x_{m}\right) d x_{1} \cdots d x_{m}, \quad \tau=t\left(x_{1}, \ldots, x_{r}\right) d x_{1} \cdots d x_{r}
$$

and

$$
\rho=\left(s\left(x_{1}, \ldots, x_{m}\right) / t\left(x_{1}, \ldots, x_{r}\right)\right) d x_{n+1} \cdots d x_{m} .
$$

Remark 2.3. Let $s_{M}$ (resp. $t_{N}$ ) be the measures defined on $M$ (resp. $N$ ) by $\left|\sigma_{M}\right|$ (resp. $\left|\tau_{N}\right|$ ). We denote by $r_{M / N, y}$ the measures defined on $\pi^{-1}(y)$ by $|\rho|$. The smooth functions in the nonarchimedean case are the locally constant functions. Then we have [Harish-Chandra 1964]

$$
\int_{M} \alpha d s_{M}=\int_{N} f_{\alpha} d t_{N}, \quad f_{\alpha}(y)=\int_{\pi^{-1}(y)} \alpha d r_{M / N, y}
$$

for all smooth compactly supported complex-valued functions $\alpha$ on $M$.
It is easy to show, using partitions of unity that the map $\alpha \mapsto f_{\alpha}$ is surjective, and continuous when $k=\mathbb{R}$. This gives rise to an injection of the space of distributions on $N$ into the space of distributions on $M$, say $T \mapsto T^{*}$. Then $r_{M / N, y}=\delta(y)^{*}, \delta(y)$ being the Dirac distribution at $y \in N$. Replacing $\delta(y)$ by its derivatives, we get distributions on $M$, supported by $\pi^{-1}(y)$. If $F$ is a locally integrable function on $N$, it defines a distribution on $N$, say $T_{F}$, and $T_{F}^{*}$ is $T_{F \circ \pi}$ where $F \circ \pi$ is a locally integrable function on $M$. Thus the map $T \mapsto T^{*}$ is the natural extension of the map $F \mapsto F \circ \pi$ from the space of locally integrable functions on $N$ to the corresponding space on $M$. The map $T \mapsto T^{*}$ plays a fundamental role in Harish-Chandra's theory [1964] of characters on real semisimple Lie groups. Finally, in algebrogeometric terminology, $\rho$ above is the top relative exterior form.

We shall now apply this result to polynomial maps $f: V \rightarrow k^{r}$ where $V$ is a vector space of finite dimension $m$ over a local field $k$ of arbitrary characteristic such that $d f_{1} \wedge \cdots \wedge d f_{r} \not \equiv 0$ on $V$, the $f_{i}$ being the components of $f$; let $V^{\times}$be the set of points where this exterior form is nonzero in $V$, so that $V^{\times}$is nonempty Zariski open in $V$; let $N(\boldsymbol{f})=\boldsymbol{f}\left(V^{\times}\right)$. Clearly $m \geq r$ and $N(\boldsymbol{f})$ is nonempty open (in the $k$-topology) in $k^{r}$. Then, by Theorem 2.2 with $M=V^{\times}, N=N(f)$, we have a measure $\mu_{\boldsymbol{c}}$ for $\boldsymbol{c} \in N$ on $L_{\boldsymbol{c}}^{\prime}:=L_{\boldsymbol{c}} \cap V^{\times}$where $L_{\boldsymbol{c}}$ is the level set

$$
\begin{equation*}
L_{c}=Z\left(f_{1}-c_{1}, \ldots, f_{r}-c_{r}\right)=\left\{x \in V \mid f_{1}(x)=c_{1}, \ldots, f_{r}(x)=c_{r}\right\} \tag{2-1}
\end{equation*}
$$

Exactly as before, we may view the $\mu_{f, c}$ as distributions living on $L_{c}^{\prime}$ which is all of $L_{\boldsymbol{c}}$ if $\boldsymbol{c}$ is a noncritical value. The derivatives of $\mu_{f, \boldsymbol{c}}$ with respect to the differential operators of $k^{m}$ (when $k=\mathbb{R}$ ) then yield distributions supported by $L_{\boldsymbol{c}}$. Examples of such distributions have important applications ([Gel'fand and Shilov 1964], [Kolk and Varadarajan 1992]) in analysis and physics.

Fix a noncritical value $\boldsymbol{c}$ of $\boldsymbol{f}$. Let $J=\left\{i_{1}<i_{2}, \ldots,<i_{r}\right\}$ be an ordered subset of $r$ elements in $\{1,2, \ldots, m\}$. Let

$$
\begin{equation*}
\partial_{J}:=\frac{\partial\left(f_{1}, \ldots, f_{r}\right)}{\partial\left(x_{i_{1}}, \ldots, x_{i_{r}}\right)} . \tag{2-2}
\end{equation*}
$$

Then $L_{c}$ is smooth and $L_{c}=\bigcup_{J} L_{c, J}$ where the sum is over all sets $J$ as above and

$$
\begin{equation*}
L_{c, J}:=\left\{x \in L_{\boldsymbol{c}} \mid \partial_{J}(x) \neq 0\right\} \tag{2-3}
\end{equation*}
$$

Locally on $L_{c, J},\left(f_{1}, \ldots, f_{r}, y_{1}, \ldots, y_{m-r}\right)$ is a new coordinate system, the $y_{j}$ being some enumeration of the $x_{i}\left(i \neq i_{v}\right)$. Obviously $d y_{1} \cdots d y_{m}=\varepsilon \partial_{J}(x) d x_{1} \cdots d x_{m}$, where $\varepsilon$ is locally constant and equal to $\pm 1$. Another way of interpreting this formula is the following: if $\pi_{J}$ is the projection map from $L_{c, J}$ that takes $x$ to $\left(y_{1}, \ldots, y_{m-r}\right)$, then $\pi_{J}$ is a local analytic isomorphism and

$$
\begin{equation*}
\rho_{c}=\varepsilon \frac{1}{\partial_{J}(x)} \pi_{J}^{*}\left(d y_{1} \cdots d y_{m-r}\right) \tag{2-4}
\end{equation*}
$$

where $\varepsilon$ is locally constant and $\pm 1$-valued. Hence to control the growth of the measure defined by $|\rho|$ at infinity, we must find lower bounds of the $\left\|\partial_{J}(x)\right\|$ on $L_{c, J}$ for $\|x\| \geq 1$. Let

$$
\nabla_{r}(x)=\left(\partial_{J}(x)\right)
$$

We call $\nabla_{r}$ the generalized gradient of $\left(f_{1}, \ldots, f_{r}\right)$. Then we must find lower bounds for $\left\|\nabla_{r}(x)\right\|:=\max _{J}\left\|\partial_{J}(x)\right\|$ for $\|x\| \geq 1$ on $L_{c, J}$. In this quest we follow [Igusa 1978], and our techniques force us to assume $\boldsymbol{c}$ to be stably noncritical. For $r=1$, $\nabla_{1}$ is just the gradient $\nabla$, and that work reduces the issue of the lower bounds for the gradient field by replacing $\nabla f$ (for $k=\mathbb{R}$ ) by $\sum_{1 \leq j \leq m}\left|\partial_{j} f\right|^{2}$, where $\partial_{j} f=\partial f / \partial x_{j}$. For nonarchimedean local fields and for $r>1$ we have to replace the sum of squares by suitable definite forms whose degrees will grow with $m$. Igusa and Raghavan find lower bounds for $|\nabla|$ using Hörmander's inequalities [1958] over $\mathbb{R}$. We generalize Hörmander's inequalities to any local field and use them with the existence of definite forms of sufficiently high degree to get lower bounds for $\left\|\nabla_{r}\right\|$ on the level sets $L_{\boldsymbol{c}, J}$.

The Hörmander inequalities over $\mathbb{R}$ are of two types: H1 and H2. H1 is local and is essentially the Łojasiewicz inequality [1959]; Hörmander derives H2 from H 1 by inversion. Over nonarchimedean $k, \mathrm{H} 1$ turns out to be a consequence of a Henselization lemma of Greenberg [1966], as observed in [Bollaerts 1990]. The reduction of H 2 to H 1 is more subtle in the nonarchimedean case. We prove it by embedding $V$ in a division algebra $D$, central over $k$, prove H 2 for $D$, and then deduce H 2 for $V$. The descent from $D$ to $V$ is elementary. To prove H 2 in $D$ we use the map $x \mapsto x^{-1}$ on $D \backslash\{0\}$ to reduce H 2 to H1. The existence of central division algebras over $k$ of arbitrarily high dimension is nontrivial and follows from the theory of the Brauer group of $k$. The lower bounds of $\nabla_{r} f$ obtained from these arguments allow us to prove that when $\boldsymbol{c}$ is a stably noncritical value of $\boldsymbol{f}$, $\mu_{f, \boldsymbol{c}}\left(B_{r}\right)=O\left(R^{m-r+\gamma}\right)$ for some $\gamma \geq 0$, uniformly near $\boldsymbol{c}$. We do not know if we can take $\gamma=0$ always. If $\left\|\nabla_{r} \boldsymbol{f}\right\|$ is bounded away from zero at infinity on $L_{\boldsymbol{c}}$, then it is obvious that we may take $\gamma=0$; but $\inf \left\|\nabla_{r} f\right\|$ may be zero on $L_{c}$. (See page 252).

Algebraic Sard's theorem in characteristic 0 for polynomial maps. Let $V$ be a vector space over $k$ of finite dimension $m$. Recall the definitions of $C(f)$ and $C V(f)$.

Proposition 2.4. Let $k$ be of characteristic 0 . The Zariski closure, $C l(C V(f))$ is a proper subset of $k^{r}$; in particular, if $r=1$, then $C V(\boldsymbol{f})$ is finite.
Proof. Fix a basis of $V$ so that $V \simeq k^{m}$. The field generated by the coefficients of the $f_{j}$, say $k_{1} \supset k$, can be embedded in $\mathbb{C}$. It is thus enough to prove Proposition 2.4 over $\mathbb{C}$ itself, where it is just the statement that the fibers of $f$ are generically smooth. Over $\mathbb{C}$ this is essentially Sard's lemma for affine algebraic varieties treated by Mumford [1995].

Analytic Sard's theorem in characteristic $\boldsymbol{p}>\mathbf{0}$. In characteristic $p>0$, the algebraic Sard's lemma fails abysmally [Mumford and Oda 2015, p. 179] over algebraically closed fields. Indeed, let $f$ be a polynomial in two variables $X, Y$ giving rise to a map $K^{2} \rightarrow K$ where $K$ is algebraically closed and of characteristic $p>0$, for example,

$$
f=X^{p+1}+X^{p} Y+Y^{p}
$$

Then the gradient of $f$ vanishes precisely on the $Y$-axis, and $f$ on the $Y$-axis is the map $y \mapsto y^{p}$ which is surjective. So the image of the singular set is all of $K$, and every fiber has a singular point. But if we replace $K$ by a local field, then $y \mapsto y^{p}$ is not surjective, and in fact the image under $f$ of the singular set is $k^{p}$ which is a closed proper subset of $k$ (in the $k$-topology), and is of measure zero in $k$. Thus the generic fiber (in the $k$-topology) is smooth in $k$.

We shall now consider the situation over local fields of characteristic $p>0$. From Sard [1942] we know that when $k=\mathbb{R}$ and the map is of class $C^{(a)}(a>0)$, $f(C)$ has measure zero when $a>m-r$. Now, when $k$ has characteristic $p>0$, the derivatives of $f$ are not enough to determine the coefficients of the power series expansion of $f$ whose order is greater than $p-1$. So there is an analogy with the case of $C^{(p-1)}$ over $\mathbb{R}$, suggesting that over $k$ the condition $p>m-r+1$ would be sufficient to guarantee that $f(C)$ is a null set. This suggestion, which leads to Theorem 2.5, is due to Professor Pierre Deligne (personal communication, 2016), which we gratefully acknowledge.
Theorem 2.5. Let $X, Y$ be analytic manifolds over a local field $k$ of characteristic $p>0$, of dimensions $m, r$ respectively. Let $f: X \rightarrow Y$ be an analytic. Let $C$ be the critical set for $f$. Then $f(C)$ has measure zero in $Y$ if $p>m-n+1$.

Proof. The proof that $f(C)$ has measure zero in $Y$ when $p>m-r+1$ is a minor adaptation of [Guillemin and Pollack 1974], needed because we have an additional restriction on $p$.

The result is local and so we may take $X$ to be a compact open set $U \in k^{m}$. We use induction on $m$. We define the filtration $C=C_{0} \supset C_{1} \supset \cdots \supset C_{p-1}$, where $C_{s}$ $(1 \leq s \leq p-1)$ is the set where all derivatives of the components of $f$ of order $\leq s$ vanish. The sets $C, C_{s}$ are compact while $C_{s} \backslash C_{s+1}$ is locally compact and second
countable, hence a countable union of compact sets. So $f(C), f\left(C_{s}\right)$ are compact, and $f\left(C_{s} \backslash C_{s+1}\right)$ is a countable union of compact sets.

The inductive proof that $f\left(C \backslash C_{1}\right)$ is a null set reduces to the case when $(m, r)$ becomes $(m-1, r-1)$. Since $m-r=(m-1)-(r-1)$, the condition on $p$ remains the same and induction applies.

The inductive proof that $f\left(C_{s} \backslash C_{s+1}\right)$ is a null set reduces to the case when $(m, r)$ becomes $(m-1, r)$. Since $p>m-r+1>(m-1)-r+1$, induction applies again.

It remains to show that $f\left(C_{p-1}\right)$ is a null set when $p>m-r+1$. We shall show actually that $f\left(C_{p-1}\right)$ is a null set when $p>m / r$. This is enough since $m / r \leq m-r+1$. This is a local result and so we may work around a point of $C_{p-1}$ which can be taken to be the origin. We use the max norm on $k^{m}$ and $k^{r}$ so that the norms take values in $q^{\mathbb{Z}}$, where $q>1$ is the cardinality of the residue field of $k$. By scaling, if necessary, we may assume that all components of $f$ are given by power series expansions, absolutely convergent on the ball $B(q):=\left\{x \in k^{m} \mid\|x\| \leq q\right\}$. Note that $B(1)=R^{m}$, where $R$ is the ring of integers of $k$. In order to estimate the growth of these series we need a lemma:
Lemma 2.6. Let $g$ be an analytic function on $B(q)$ given by an absolutely convergent power series expansion about 0 on $B(q)$. Let $D$ be the set in $B(1)$ where $\partial^{\beta} f=0$ for all $\beta$ with $|\beta| \leq p-1$. Then we have

$$
|g(x+h)-g(x)| \leq A\|h\|^{p}
$$

uniformly for $x \in D,\|h\| \leq 1 \leq q-1$, the constant $A>0$ depending only on $g$.
Proof. We use [Serre 2006, pp. 67-75]. We have

$$
g(x)=\sum_{\alpha} c_{\alpha} X^{\alpha}, \quad \sum_{\alpha}\left|c_{\alpha}\right|=A<\infty
$$

For $x \in B(1)$ we have $g(x+h)=\sum_{\beta} \Delta^{\beta} g(x) h^{\beta}$, where

$$
\Delta^{\beta} g(x)=\sum_{\alpha \geq \beta} c_{\alpha}\binom{\alpha}{\beta} x^{\alpha-\beta}, \quad \beta!\Delta^{\beta} g(x)=\partial^{\beta} g(x)
$$

Then $\left|\Delta^{\beta} g(x)\right| \leq A$ on $B(1)$. If $x \in D,\|h\| \leq 1 \leq q-1$, then $x+h \in B(1)$. Moreover, for $|\beta| \leq p-1, \beta!\Delta^{\beta} g(x)=0$ so that $\Delta^{\beta} g(x)=0$. Hence,

But, for $y \in B(1)$,

$$
g(x+h)=g(x)+\sum_{|\beta| \geq p}\left(\Delta^{\beta} g\right)(x) h^{\beta}
$$

So,

$$
\left|\Delta^{\beta} g(y)\right| \leq \sum\left|c_{\alpha}\right|=A
$$

$$
|g(x+h)-g(x)| \leq A\|h\|^{p} \quad(x \in D,\|h\| \leq 1)
$$

proving the lemma.

We now divide $B(1)^{m}$ into very small "cells". Let $P$ be the maximal ideal in $R$. Let $N$ be any integer $\geq 1$. Then $B(1)$ is the disjoint union of $q^{N}$ cosets of $P^{N}$ each of which is a compact open set that has diameter $\leq q^{-N}$ and volume $q^{-N}$. This gives a partition of $B(1)^{m}$ into $q^{m N}$ compact open sets ("cells") of diameter $\leq q^{-N}$ and volume $q^{-N m}$. By the above lemma, if $x, x+h \in D$ and are in one of these cells, say $\gamma$, then

$$
\begin{equation*}
\|f(x+h)-f(x)\| \leq A\|h\|^{p} \leq A q^{-N p} \tag{2-5}
\end{equation*}
$$

where $A$ is a constant independent of $x$. Hence, $f(\gamma)$ is contained in a set of diameter $\leq q^{-N p}$ and hence volume $\leq q^{-N p r}$. Thus $f\left(D \cap C_{p-1}\right)$ is enclosed in a set of volume $\leq q^{m N-N p r}=q^{-N(p r-m)}$. If $p>m / r$ this expression goes to 0 as $N \rightarrow \infty$, and we are done.

Remark 2.7. If $\boldsymbol{f}=\left(f_{1}, \ldots, f_{r}\right)$ is a polynomial map of $k^{m}$ into $k^{r}$ such that $d f_{1} \wedge \cdots \wedge d f_{r} \neq 0$, then $\boldsymbol{f}\left(k^{m}\right)$ is open and Sard's theorem shows that almost every fiber of $f$ is smooth in $k$. So there are always noncritical values. Whether some of them are stable is not known to us.

Remark 2.8. When $r=1$, the above condition reduces to $p>m$. Both this condition and the fact that when $m \geq p+1$ it is possible that the image of the critical set can be all of $k$ were communicated to us by Professor Pierre Deligne (2016). We are grateful for his generosity and for giving us permission to discuss his example.

Example 2.9 (Deligne). We take $m=p+1$ with coordinates $y, x_{1}, \ldots, x_{p}$. The field $k:=\mathbb{F} \llbracket t \rrbracket[1 / t]$, where $\mathbb{F}$ is a finite field of characteristic $p$, is a local field of characteristic $p$. Then $k$ is a vector space of dimension $p$ over $k^{(p)}:=\left\{x^{p} \mid x \in k\right\}$. Let $\left(a_{i}\right)_{1 \leq i \leq p}$ be a basis for $k / k^{(p)}$, for instance $a_{i}=t^{i-1},(1 \leq i \leq p)$. Consider, for an integer $n>1$, prime to $p$,

$$
f=y^{n}+a_{1} x_{1}^{p}+\cdots+a_{p} x_{p}^{p}
$$

Then the critical locus is given by $y=0$. Its image under $f$ is obviously all of $k$. If we do not insist that $d f \not \equiv 0$, we can omit $y$ so that $f$ maps the critical set $k^{p}$ onto $k$.

This example is easily modified for the case $r>1$. We consider $k^{p+r}$ with coordinates $y_{1}, \ldots, y_{r-1}, y, x_{1} \ldots, x_{p}$ and take the map $f: k^{p+r} \rightarrow k^{r}$ defined by

$$
f:\left(y_{1}, \ldots, y_{r-1}, y, x_{1} \ldots, x_{p}\right) \mapsto\left(y_{1}, \ldots, y_{r-1}, y^{n}+\sum_{i=1}^{p} a_{i} x_{i}^{p}\right)
$$

where the notation is as before. The critical set is again given by $y=0$, and the map restricted to this set is

$$
f:\left(y_{1}, \ldots, y_{r-1}, 0, x_{1} \ldots, x_{p}\right) \mapsto\left(y_{1}, \ldots, y_{r-1}, \sum_{i=1}^{p} a_{i} x_{i}^{p}\right),
$$

whose range is $k^{r}$. Exactly as before, if we omit $y$, we get a map where $d f_{1} \wedge \cdots \wedge d f_{r}$ is zero but $f$ maps the critical set $k^{p+r-1}$ onto $k^{r}$.

## 3. Construction of definite forms and their associated norms

As mentioned in Remark 2.3 we begin by discussing the construction of definite forms in an arbitrary number of variables over $k$.

Proposition 3.1. Let $V$ be a finite dimensional vector space over a local field. If $k=\mathbb{R}$, and $\nu(x)$ is a positive definite quadratic form on $V$, then $|\nu(x)|^{1 / 2}$ is a norm on $V$. If $k$ is nonarchimedean, and $r$ is an integer such that $r^{2} \geq m$, then there is a homogeneous polynomial $v: V \rightarrow k$ of degree $r$ such that
(a) $v$ is definite, i.e., for $x \in V, v(x)=0$ if and only if $x=0$;
(b) $|v(x)|^{1 / r}$ is a nonarchimedean norm on $V$.

Proof. We deal only with the case of nonarchimedean $k$. By the theory of the Brauer group of $k$ [Weil 1967, chapter XII, theorem 1] and its corollary we can find a division algebra $D$ over $k$ which is central over $k$ and $\operatorname{dim}_{k}(D)=r^{2}$. Since $V \hookrightarrow D$, it is enough to prove the proposition for $V=D$. The advantage is that we can use the algebraic structure of $D$.

Let $v$ be the reduced norm [Weil 1967, chapter IX, proposition 6] of $D$. Then, $v: D \rightarrow k$ is a homogeneous polynomial function on $D$ of degree $r$, and $v(x)^{r}=$ $\operatorname{det}(\lambda(x))$ where $\lambda(x)$ is the endomorphism $y \mapsto x y$ of $D$. Note that $\operatorname{det}(\lambda)$ is a polynomial function on $D$ with values in $k$, homogeneous of degree $r^{2}$. As $\lambda(x)$ is invertible for any $x \neq 0$ in $D, \operatorname{det}(\lambda(x))$ and hence $v(x)$, is nonzero for $x \neq 0$ in $D$. Hence, $v$ is a definite form of degree $r$ on $D$. It remains to prove that $N(x):=|\nu(x)|^{1 / r}$ is a nonarchimedean norm on $D$. This reduces to showing that $N(1+u) \leq 1$ if $u \in D$ and $N(u) \leq 1$, or equivalently, that $|\lambda(1+u)| \leq 1$ if $u \in D$ and $|\lambda(u)| \leq 1$, which follows from [Weil 1967, chapter I, section 4].

Remark 3.2. Actually, $v(x)^{r}=\operatorname{det} \lambda(x)$ will serve our purposes as well and is obviously a homogeneous polynomial of degree $r^{2}$, Then $|\nu(x)|^{1 / r}=|\operatorname{det} \lambda(x)|^{1 / r^{2}}$. We introduced $v$ because it is of smaller degree and this may be of use in other contexts.

## 4. Hörmander's inequalities over nonarchimedean local fields

Let $V$ be a finite dimensional vector space over a local, nonarchimedean field $k$, with its canonical norm $|\cdot|$. Let $\|\cdot\|$ be a nonarchimedean norm on $V$. We may assume that the norms on $k$ and $V$ take values in the set $\left\{0, q^{ \pm 1}, q^{ \pm 2}, \ldots\right\}$, where $q$ is the cardinality of the residue field of $k$. Also, let $f: V \rightarrow k$ be a polynomial function, and let $Z(f)$ denote its zero locus. For $x \in V$ and nonempty $E \subset V$ let $\operatorname{dist}(x, E):=\inf _{y \in E}\|x-y\|$.

Theorem 4.1 (H1). Let $f: V \rightarrow k$ be a polynomial function on $V$. Suppose that $Z(f) \neq \varnothing$. Then there exist constants $C>0, \alpha \geq 0$ such that

$$
\begin{equation*}
|f(x)| \geq C \cdot \operatorname{dist}(x, Z(f))^{\alpha} \tag{4-1}
\end{equation*}
$$

for all $x \in V$ with $\|x\| \leq 1$.
Theorem 4.2 (H2). Let $f: V \rightarrow k$ be a polynomial function, $Z(f)$ as above. Then
(a) if $Z(f)=\varnothing$, then there exist constants $C>0$ and $\beta \geq 0$ such that

$$
\begin{equation*}
|f(x)| \geq C \cdot \frac{1}{\|x\|^{\beta}} \quad(x \in V,\|x\| \geq 1) \tag{4-2}
\end{equation*}
$$

(b) if $Z(f) \neq \varnothing$, then there exist constants $C>0$ and $\alpha, \beta \geq 0$ such that

$$
\begin{equation*}
\|f(x)\| \geq C \cdot \frac{\operatorname{dist}(x, Z(f))^{\alpha}}{\|x\|^{\beta}} \quad(x \in V,\|x\| \geq 1) \tag{4-3}
\end{equation*}
$$

Remark 4.3. Theorem 4.1 and 4.2 were proved by Hörmander [1958] when $k=\mathbb{R}$. Also, H 1 is a special case of the Łojasiewicz inequality for $f$ a real analytic function [Łojasiewicz 1959].

In proving H1 we may assume that $V=k^{m}$ and $f \in R\left[x_{1}, \ldots, x_{m}\right], R$ being the ring of integers in $k$. Let $P \subset R$ be the maximal ideal of $R$. Suppose that $Z(f) \neq \varnothing$ but $Z(f) \cap R^{m}=\varnothing$. Then there exists a constant $b>0$ such that $|f(x)| \geq b>0$ for $x \in R^{m}$. On the other hand, as $R^{m}$ is compact, there exists $b_{1}>0$ such that $\operatorname{dist}(x, Z(f)) \leq b_{1}$ for all $x \in R^{m}$. Hence $|f(x)| \geq b b_{1}^{-1} b_{1} \geq b b_{1}^{-1} \operatorname{dist}(x, Z(f))$ for all $x \in R^{m}$. Hence we may assume in addition that $Z(f) \cap R^{m} \neq \varnothing$ in the proof of H1.

Proof of H1: $\boldsymbol{k}$ nonarchimedean. We follow [Greenberg 1966], specialized to the case of a single polynomial.

Proof. By theorem 1 there, applied to the single polynomial $f$, we can find integers, $N, c \geq 1$ and $s \geq 0$ such that if $v \geq N$ and $f(x) \equiv 0\left(\bmod P^{v}\right)$, and $x \in R^{m}$, then there exists $y \in R^{m}$ such that $f(y)=0$ and $x_{i}-y_{i} \equiv 0\left(\bmod P^{[\nu / c]-s}\right)$ for all $i$.

Assume $|f(x)|=q^{-(N+\ell)}, \ell \geq 0$. Then there exists $y \in Z(f) \cap R^{m}$ such that

$$
\|x-y\| \leq q^{-[(N+\ell) / c]+s} \leq q^{-[(N+\ell) / c-1]+s} \leq q^{s+1}|f(x)|^{1 / c},
$$

which implies that

$$
\operatorname{dist}\left(x, Z(f) \cap R^{m}\right) \leq q^{s+1}|f(x)|^{1 / c}
$$

so that

$$
|f(x)| \geq \frac{\operatorname{dist}\left(x, Z(f) \cap R^{m}\right)^{c}}{q^{c(s+1)}} \geq \frac{\operatorname{dist}(x, Z(f))^{c}}{q^{c(s+1)}}
$$

Thus, H1 is proved for $x \in R^{m}$ with $|f(x)| \leq q^{-N}$. For $x$ in $R^{m}$ with $|f(x)|>q^{-N}$, we have $q^{-N}<|f(x)| \leq 1$, while $\operatorname{dist}\left(x, Z(f) \cap R^{m}\right) \leq 1$ since $\|x-y\| \leq 1$ for $x, y \in R^{m}$. Hence
$|f(x)| \geq q^{-N} \operatorname{dist}\left(x, Z(f) \cap R^{m}\right) \geq q^{-N} \operatorname{dist}\left(x, Z(f) \cap R^{m}\right)^{c} \geq q^{-N} \operatorname{dist}(x, Z(f))^{c}$.
If $C=\min \left(q^{-N}, q^{-(s+1) c}\right)$, then we have H1 with $\alpha=c$.
Remark 4.4. That the local version of the Łojasiewicz inequality comes out of [Greenberg 1966] has been observed in [Bollaerts 1990]; we give this proof since it includes the case when $k$ has characteristic $>0$. Greenberg's result is applicable here because $R$ is then complete ( $k^{*}=k$ in his notation).

## Proof of H2.

Lemma 4.5. If H 2 is true for a $k$-vector space $V$, then it is also true for any subspace $W$ of $V$. In particular, for a central division algebra, $D_{r}$ over $k$, of dimension $r^{2} \geq \operatorname{dim}_{k} V$, it is enough to prove H 2 for $D_{r}$.
Proof. Let $W \subseteq V$ be a subspace, and $U \subseteq V$ such that $V=W \oplus U \simeq W \times U$. Let $f$ be a polynomial on $W$. Define the polynomial $g$ on $V$ by $g(w+u):=f(w)$. For $w \in W, u \in U$, we take $\|w+u\|=\max (\|u\|,\|w\|)$; because $U$ and $W$ are complementary, this is nonarchimedean. Clearly $Z(g)=Z(f) \times U$.

Suppose $Z(f)=\varnothing$. Then $Z(g)=\varnothing$. Since H 2 is true for $V$ and $W \subset V$, there exist constants $C>0, \beta \geq 0$ such that $|f(w)| \geq C\|w\|^{-\beta}$ for $w \in W,\|w\| \geq 1$. We may therefore assume that $Z(f) \neq \varnothing$, so $Z(g) \neq \varnothing$.

Then, $|g(x)| \geq C \operatorname{dist}(x, Z(g))^{\alpha}\|x\|^{-\beta}$ for $x \in V,\|x\| \geq 1$ where $C>0, \alpha, \beta \geq 0$ are constants. If $x=w \in W$, $\operatorname{dist}(w, Z(g))=\operatorname{dist}(w, Z(f))$.

Now we prove H2 for $D_{r}$. Our proof is inspired by Hörmander's [1958]. It replaces the inversion in his proof by the involution $x \mapsto x^{-1}$ of $D_{r}^{\times}:=D_{r} \backslash\{0\}$.

For a division algebra $D_{r}$ of dimension $r^{2}$, central over $k$, let us recall $v:=$ $v_{r}: D_{r} \rightarrow k$ of Proposition 3.1, and note that it has the following property: if $k^{\prime}$ is any field containing $k$ such that there exists an isomorphism $F: k^{\prime} \otimes_{k} D_{r} \xrightarrow{\sim}$ $M_{r}\left(k^{\prime}\right)=M_{r}$ where $M_{r}$ is the algebra of $r \times r$ matrices over $k^{\prime}$, then $\nu(a)=\operatorname{det} F(a)$ for $a \in D_{r}$ [Weil 1967, Proposition 6, p. 168],

Lemma 4.6. For any polynomial function $f: D \rightarrow k$ of degree $d$, $f$ not necessarily homogeneous, let $f^{*}(x):=f\left(x^{-1}\right) \nu(x)^{d}$ for $x \neq 0$; then $f^{*}(x)$ extends uniquely to a polynomial function $D_{r} \rightarrow k$. Moreover, for nonzero $x, x \in Z(f)$ if and only if $x^{-1} \in Z\left(f^{*}\right)$.
Proof. Uniqueness is obvious. To prove that $f^{*}$ has a polynomial extension it suffices to prove it for $k^{\prime} \otimes_{k} D_{r}$, where $k^{\prime}$ is a separable extension of $k$ such that $k^{\prime} \otimes_{k} D_{r} \simeq M_{r}\left(k^{\prime}\right)$. The required result is compatible with addition and multiplication
of the $f$ so that it is enough to verify it for $f=1$ (obvious) and $f=a_{i j}$, a matrix entry; then $f^{*}=a^{i j} \operatorname{det}=A_{i j}$, the corresponding cofactor. The last statement of the lemma is obvious
Remark 4.7. From now on we use the norm $\|x\|=|\nu(x)|^{1 / r}$ for $D_{r}, r \geq 2$.
Lemma 4.8. If $x, y, x-y$ are all nonzero, then $\|x-y\|=\left\|x^{-1}-y^{-1}\right\|\|x\|\|y\|$
Proof. Use $y-x=x\left(x^{-1}-y^{-1}\right) y$ and the multiplicativity of $\|\cdot\|$.
The next two lemmas are auxiliary before we prove H 2 for $D_{r}$.
Lemma 4.9. If $Z(f)$ is nonempty, there exists a constant $A \geq 1$ such that

$$
\operatorname{dist}(x, Z(f)) \leq A\|x\| \quad \text { for all } x \text { with }\|x\| \geq 1
$$

Proof. Choose $z_{0} \in Z(f)$. Then $\operatorname{dist}(x, Z(f)) \leq\left\|x-z_{0}\right\| \leq \max \left(\|x\|,\left\|z_{0}\right\|\right)$. If $\|x\| \geq\left\|z_{0}\right\|$, then $\operatorname{dist}(x, Z(f)) \leq\|x\|$ and we can take $A=1$. If $\|x\|<\left\|z_{0}\right\|$ then $\left\|x-z_{0}\right\|=\left\|z_{0}\right\| \leq\left\|z_{0}\right\|\|x\|$ for $\|x\| \geq 1$; and as $\left\|z_{0}\right\| \geq 1$, the lemma is proved if we take $A=1+\left\|z_{0}\right\|$.
Lemma 4.10. Suppose $Z(f)$ contains a nonzero element. Then there exists a constant $C>0$ such that

$$
\begin{equation*}
\operatorname{dist}\left(x^{-1}, Z\left(f^{*}\right)\right) \geq C \frac{\operatorname{dist}(x, Z(f))}{\|x\|^{2}} \quad(\|x\| \geq 1) \tag{4-4}
\end{equation*}
$$

Proof. First assume $0 \notin Z\left(f^{*}\right)$. Then $Z\left(f^{*}\right)=Z\left(f^{*}\right) \backslash\{0\} \neq \varnothing$. Then, with $\|x\| \geq 1$,
$\operatorname{dist}\left(x^{-1}, Z\left(f^{*}\right) \backslash\{0\}\right)=\inf _{0 \neq z \in Z\left(f^{*}\right)}\left\|x^{-1}-z\right\|=\inf _{0 \neq y \in Z(f)}\left\|x^{-1}-y^{-1}\right\|=\inf _{0 \neq y \in Z(f)} E$, where $E:=\|x-y\|\|x\|^{-1}\|y\|^{-1}$.

We consider cases: (a) $\|y\|>\|x\|$ and (b) $\|y\| \leq\|x\|$. In case (a) $\|x-y\|=\|y\|$ so that $E=\|x\|^{-1}=\|x\|\|x\|^{-2} \geq A^{-1} \operatorname{dist}(x, Z(f))\|x\|^{-2}$, where $A \geq 1$ is as in Lemma 4.9. In case (b) $E \geq\|x-y\|\|x\|^{-2}$ so that $\inf E \geq \operatorname{dist}(x, Z(f))\|x\|^{-2}$. These give (4-4) with $C=1 / A$.

If $0 \in Z\left(f^{*}\right)$, then $\operatorname{dist}\left(x^{-1}, Z\left(f^{*}\right)\right)=\min \left(\operatorname{dist}\left(x^{-1}, Z\left(f^{*}\right) \backslash\{0\}\right),\left\|x^{-1}\right\|\right)$. Now $\|x\|^{-1}=\|x\|\|x\|^{-2} \geq C\|x\|^{-2} \operatorname{dist}(x, Z(f))$ by Lemma 4.9 where $C=1 / A$, while $\operatorname{dist}\left(x^{-1}, Z\left(f^{*}\right) \backslash\{0\}\right) \geq C\|x\|^{-2} \operatorname{dist}(x, Z(f))$, by above.
Proof of H 2 for $D_{r}$. We consider two cases: (a) $Z(f)=\varnothing$, (b) $Z(f) \neq \varnothing$.
Case (a): Then $Z\left(f^{*}\right)=\varnothing$ or $\{0\}$. If $Z\left(f^{*}\right)=\varnothing$, then there exists a constant $C>0$ such that $\left|f^{*}(x)\right| \geq C>0$ with $\|x\| \leq 1$. So, $\left|f^{*}(y)\right|=\left|f\left(y^{-1}\right)\right|\|y\|^{r d} \geq C>0$ for $0<\|y\| \leq 1$, which becomes $|f(x)| \geq C\|x\|^{r d} \geq C>0$ for $\|x\| \geq 1$.

If $Z\left(f^{*}\right)=\{0\}$, then $\operatorname{dist}\left(z, Z\left(f^{*}\right)\right)=\|z\|$, and $\left|f^{*}(y)\right| \geq C\|y\|^{\beta}$ with $0<\|y\| \leq 1$ for constants $C>0, \beta \geq 0$ by Theorem 4.1. Then $\left|f\left(y^{-1}\right)\right|\|y\|^{r d} \geq C\|y\|^{\beta}$ with $\|y\| \leq 1$ or $|f(x)| \geq C\|x\|^{r d}\|x\|^{-\beta} \geq C\|x\|^{-\beta}$ with $\|x\| \geq 1$.

Case (b): $Z(f)$ is now nonempty, and hence either $Z(f)=\{0\}$ or $Z(f)$ contains a nonzero element. If $Z(f)=\{0\}$, then $Z\left(f^{*}\right)=\varnothing$ or $\{0\}$. This comes under case (a), above, and we have $|f(x)| \geq C\|x\|^{-\beta}$ with $\|x\| \geq 1$ which gives (a).

Suppose $Z(f)$ contains a nonzero element. By H1, there exists constants $C_{1}>0$, $\alpha \geq 0$ such that $\left|f^{*}\left(x^{-1}\right)\right| \geq C_{1} \operatorname{dist}\left(x^{-1}, Z\left(f^{*}\right)\right)^{\alpha}$ with $\|x\| \geq 1$. So by Lemma 4.10, for $C_{2}=C_{1} C^{\alpha},|f(x)| \geq C_{2} \operatorname{dist}(x, Z(f))^{\alpha}\|x\|^{-2 \alpha}$ for $\|x\| \geq 1$, proving (b).

Criterion for a polynomial not to be rapidly decreasing on a set S. In [Igusa 1978] Igusa and Raghavan develop what is essentially a criterion for a polynomial on an real vector space not to be rapidly decreasing on a set of vectors of norm $\geq 1$. In this section we generalize that method to all local fields, introducing several polynomials in the criterion.
Lemma 4.11. Let $f: V \rightarrow k^{r}$ be a polynomial map and $d$ the maximum of the degrees of its components. Then there exists a constant $C>0$ such that for all $x, y \in V$ with $\|x\| \geq 1$,

$$
\|\boldsymbol{f}(x)-\boldsymbol{f}(y)\| \leq C\|x\|^{d-1} \max _{0 \leq r \leq d}\left(\|x-y\|^{r}\right)
$$

Proof. It is enough to prove this for $r=1, f=f$. The estimate is compatible with addition in $f$ and so we may assume $f$ to be a monomial of degree $d$ in some coordinate system on $V$. Assume the result for all monomials of degree $d-1$. Then $f=x_{i} g$, where $g$ is a monomial of degree $d-1$. We have

$$
x_{i} g(x)-y_{i} g(y)=x_{i}(g(x)-g(y))+\left(x_{i}-y_{i}\right)(g(y)-g(x))+\left(x_{i}-y_{i}\right) g(x)
$$

and the estimate is obvious for each of the three terms.
Proposition 4.12. Let $S \subseteq V$ be a set with $\|x\| \geq 1$ for all $x \in S$. Let $g$ be polynomial on $V$. If $Z(g)=\varnothing$, we have

$$
|g(x)| \geq \frac{C}{\|x\|^{\gamma}} \quad(\|x\| \geq 1)
$$

for some $C>0, \gamma \geq 0$. Suppose $Z(g) \neq \varnothing$ and suppose that there exist polynomials $f_{i}: V \rightarrow k, i=1, \ldots r$, and a constant $b>0$ such that $\max \left|f_{i}(x)-f_{i}(y)\right| \geq b>0$ for all $x \in S, y \in Z(g)$. Then there exist constants $C>0$ and $\gamma \geq 0$ such that

$$
\begin{equation*}
|g(x)| \geq \frac{C}{\|x\|^{\gamma}} \quad(x \in S) \tag{4-5}
\end{equation*}
$$

Proof. The first statement is (a) of H2. We now assume $Z(g) \neq \varnothing$. We identify $V \simeq k^{m}$ and work in coordinates. Set $d:=\max _{i}\left(\operatorname{deg}\left(f_{i}\right)\right)$. In what follows, $C_{1}, C_{2}, \ldots$, are constants $>0$.

For all $x \in S$ and $y \in Z(g)$, by Lemma 4.11 for some constant $C>0$, we have $0<b \leq \max _{1 \leq i \leq r}\left|f_{i}(x)-f_{i}(y)\right| \leq C\|x\|^{d-1} \max _{1 \leq r \leq r}\|x-y\|^{r}$ for all $x \in S$, $y \in Z(g)$.

Choose $y \in Z(g)$ such that $\|x-y\|=\operatorname{dist}(x, Z(g))$. Then for all $x \in S$, we have

$$
0<b \leq C_{1}\|x\|^{d-1} \max _{1 \leq r \leq d}\left(\operatorname{dist}(x, Z(g))^{r}\right)
$$

We consider the two cases (a) $\operatorname{dist}(x, Z(g)) \leq 1$, so the maximum above is $\operatorname{dist}(x, Z(g))$, and (b) $\operatorname{dist}(x, Z(g))>1$, so the maximum is $\operatorname{dist}(x, Z(g))^{d}$.

By H2, there exist constants $C_{2}>0, \alpha, \beta \geq 0$ such that

$$
|g(x)| \geq C_{2} \operatorname{dist}(x, Z(g))^{\alpha}\|x\|^{-\beta}
$$

so dist $(x, Z(g)) \leq C_{3}|g(x)|^{1 / \alpha}\|x\|^{\beta / \alpha}$. In case (a), $0<b \leq C_{3}|g(x)|^{1 / \alpha}\|x\|^{\beta / \alpha+(d-1)}$, and in case (b), $0<b \leq C_{4}|g(x)|^{d / \alpha}\|x\|^{d \beta / \alpha+(d-1)}$. So in both cases, with $\delta=$ $d \beta / \alpha+(d-1)$, one has

$$
0<b \leq C_{5}\|x\|^{\delta} \max \left(|g(x)|^{1 / \alpha},|g(x)|^{d / \alpha}\right) .
$$

Hence, $\max \left(|g(x)|,|g(x)|^{d}\right) \geq C_{6}\|x\|^{-\delta \alpha}$, giving in all cases $|g(x)| \geq C_{7}\|x\|^{-\delta \alpha}$ with $x \in S$.

Lower bounds of $\left\|\nabla_{r} \boldsymbol{f}\right\|$ on stably noncritical level sets. Let $V$ and $\boldsymbol{f}=: V \rightarrow k^{r}$ $\left(\boldsymbol{f}=\left(f_{1}, \ldots f_{r}\right), r \leq m=\operatorname{dim}_{k} V\right)$ be as usual. Let $C(\boldsymbol{f})$ be the critical set of $\boldsymbol{f}$, and $C V(f)=f(C(f))$ have their usual meanings. Write $W=C V(\boldsymbol{f})$. We assume that the closure $\bar{W}$, in the $k$-topology of $k^{r}$, of $W$ is a proper subset of $k^{r}$. Our assumption is equivalent to assuming that stably noncritical values of $\boldsymbol{f}$ exist, which is true in characteristic zero (see page 232). Let $L_{\boldsymbol{c}}, \nabla_{r} f$, and $\partial_{J} f$ be defined as in Section 2.

If $\omega \subset k^{r} \backslash \bar{W}$ is a compact set, then there exists $b>0$ such that $\|u-v\| \geq b>0$ for $u \in \omega, v \in \bar{W}$. This means $\max _{i}\left|f_{i}(x)-f_{i}(y)\right| \geq b>0$, with $\boldsymbol{c} \in \omega, x \in L_{\boldsymbol{c}}, y \in$ $C(f)$.

Proposition 4.13. Let $\omega \subset k^{r}$ be an open set whose closure consists entirely of noncritical values of $\boldsymbol{f}=\left(f_{1}, \ldots, f_{r}\right)$. For $\boldsymbol{c} \in \omega$, let $L_{\boldsymbol{c}}$ be defined as above. Then there exist constants, $C, \gamma>0$ such that

$$
\begin{equation*}
\left\|\nabla_{r} \boldsymbol{f}(x)\right\| \geq \frac{C}{\|x\|^{\gamma}} \quad\left(x \in L_{c}, \boldsymbol{c} \in \omega,\|x\| \geq 1\right) \tag{4-6}
\end{equation*}
$$

Proof. We write $\left(y_{J}\right)$ for the coordinates on $k\binom{m}{r}$ and select a definite homogeneous form $v$, which is positive definite of degree 2 if $k$ archimedean, and of degree $R$ on $k\binom{m}{r}$, where $R$ is any integer $\geq 2$ such that $R^{2} \geq\binom{ m}{r}$, with the property that $|\nu(y)|^{1 / R}$ is a norm on $k^{\left({ }_{r}^{m}\right)}$, if $k$ is nonarchimedean. Then $v\left(\nabla_{r} \boldsymbol{f}(x)\right)=0$ if and only if $\nabla_{r} \boldsymbol{f}(x)=0$, i.e., if and only if $x$ is a critical point of $\boldsymbol{f}$. Let $g(x)=v\left(\nabla_{r} \boldsymbol{f}(x)\right)$. Then $Z(g)$ is the set of critical points of $f$. Suppose first that $Z(g) \neq \varnothing$. Now
there exists $b>0$ such that

$$
\|u-v\|=\max _{1 \leq i \leq r}\left|u_{i}-v_{i}\right| \geq b>0 \quad(u \in \omega, v \in \bar{W})
$$

Hence, as $\boldsymbol{f}(x) \in \omega$ for $x \in L_{\boldsymbol{c}}(c \in \omega)$ and $\boldsymbol{f}(y) \in \bar{W}$ for $y \in Z(g),\|\boldsymbol{f}(x)-\boldsymbol{f}(y)\| \geq$ $b>0$. So by Proposition 4.12 there exist constants $C>0, \delta \geq 0$ such that

$$
\left|v\left(\nabla_{r} f(x)\right)\right|=|g(x)| \geq \frac{C}{\|x\|^{\delta}} \quad\left(x \in L_{c}, \boldsymbol{c} \in \omega,\|x\| \geq 1\right)
$$

But $v$ is homogeneous of degree $d(d=2$ for archimedean and $R$ for nonarchimedean $k$ ) and definite. So there exist constants $C_{1}, C_{2}>0$ such that

$$
C_{1}\left\|\nabla_{r} \boldsymbol{f}(x)\right\|^{d} \leq\left|\nu\left(\nabla_{r} \boldsymbol{f}(x)\right)\right|=|g(x)| \leq C_{2}\left\|\nabla_{r} \boldsymbol{f}(x)\right\|^{d}
$$

So for suitable $C>0, \gamma \geq 0$, we have $\left\|\nabla_{r} \boldsymbol{f}(x)\right\| \geq C\|x\|^{-\gamma}$. The case $Z(g)=\varnothing$ is taken care of by the first statement of Proposition 4.12.

Remark 4.14. We cannot make $\gamma=0$ in all cases. For instance, let char $k=0$ and $r=1, f(x, y, z)=x^{2} z^{2}+y^{3} z$ and $c=-1$. Consider $x_{n}=n, z_{n}=1 / n$, $y_{n}=-(2 n)^{1 / 3}$. Then $F\left(x_{n}, y_{n}, z_{n}\right)=1-2=-1, \partial F / \partial X\left(x_{n}, y_{n}, z_{n}\right)=2 x_{n} z_{n}^{2} \rightarrow 0$, and $\partial F / \partial Y\left(x_{n}, y_{n}, z_{n}\right)=3 y_{n}^{2} z_{n} \rightarrow 0, \partial F / \partial Z\left(x_{n}, y_{n}, z_{n}\right)=2 x_{n}^{2} z_{n}+y_{n}^{3}=2 n-2 n=0$. But $\left\|\left(x_{n}, y_{n}, z_{n}\right)\right\|=n,\left\|\nabla f\left(x_{n}, y_{n}, z_{m}\right)\right\| \sim$ Const $\cdot 1 / n^{1 / 3}$. So $\gamma \geq 1 / 3$. We do not know the minimal value of $\gamma$.

## 5. Proof of temperedness of canonical measures on stably noncritical level sets

Consequences of Krasner's lemma. The well-known lemma of Krasner [Artin 1967] has an important consequence (Corollary 5.3). Let $k$ be a local field of arbitrary characteristic and $K$ its algebraic closure. The following lemma must be well known, but we prove it in this form.

Lemma 5.1. We can find a countable family $\left\{k_{n}\right\}$ of finite extensions of $k$ with the property that any finite extension of $k$ is contained in one of the $k_{n}$. In particular $K=\bigcup_{n} k_{n}$.
Proof. We first work with separable extensions of fixed degree $n$ over $k$. Let $S_{n}$ be the set of monic, irreducible and separable elements of $k[X]$ of degree $n$. Then it follows from Krasner's lemma that if $f \in S_{n}$, there is an $\varepsilon=\varepsilon(f)>0$ with the following property: if $g$ is monic and $\|f-g\|<\varepsilon$, then $g \in S_{n}$ and $K(f)=K(g)$ in $K$, where $K(h)$ denotes the splitting field of $h$. Since $S_{n}$ is a separable metric space, it follows that there are at most a countable number of these splitting fields, and any separable extension of degree $n$ over $k$ is contained in one of these. Let us enumerate these splitting fields as $\left\{k_{n j}\right\}(j=1,2, \ldots)$. If $k$ has characteristic 0
we are already finished. Suppose $k$ has characteristic $p>0$. Let $F\left(x \mapsto x^{p}\right)$ be the Frobenius automorphism of $K$. Define the extension $k_{n j r}=F^{-r}\left(k_{n j}\right)$ for $r=1,2, \ldots$, which are clearly finite over $k$. Clearly, any finite extension of $k$ of finite degree is contained in one of the $k_{n j r}$.
Remark 5.2. If $k$ has characteristic 0 , then there are only a finite number of extensions of fixed degree $n$. But in prime characteristic this is not true: the field $k=F_{2} \llbracket X \rrbracket\left[X^{-1}\right]$ of Laurent series in $X$ with $F_{2}$ a finite field of characteristic 2 has a countably infinite number of separable quadratic extensions. Indeed, the extensions defined by $T^{2}-T-c=0$ are distinct for infinitely many values of $c$.

Corollary 5.3. If $M$ is an affine subvariety of some $A_{K}^{n}$ and $M\left(k^{\prime}\right)$ is countable for all finite extensions $k^{\prime}$ of $k$, then $M$ is finite.
Proof. By Lemma 5.1, $M(K)=\bigcup_{k^{\prime}} M\left(k^{\prime}\right)$ is countable, hence finite.
A consequence of the refined Bézout's theorem. The refinement of Bézout's theorem due to Fulton [1998, Example 8.4.7, p. 148, and Section 12.3] (see also [Vogel 1984, Corollary 2.26, p. 85]), is the statement that if $Z_{i}(1 \leq i \leq r)$ are $r(r \geq 2)$ pure dimensional varieties in $\mathbb{P}_{K}^{m}$, then the number of irreducible components of $\bigcap_{i} Z_{i}$ is bounded by the Bézout number $\prod_{i} \operatorname{deg}\left(Z_{i}\right)$. It has the following simple consequence.
Lemma 5.4. Let $U$ be a nonempty Zariski open subset of $\mathbb{A}_{K}^{r}$ so that $U \subset \mathbb{A}_{K}^{r} \subset \mathbb{P}_{K}^{r}$. Let $h_{i}(i=1,2, \ldots, r)$ be polynomials on $A_{K}^{r}$ with $\operatorname{deg} h_{i}=: d_{i}$, and let $Z_{i}$ be the zero locus of $h_{i}$. Let $Z_{i}^{\times}=Z_{i} \cap U$ and $\bar{Z}_{i}$ the closure of $Z_{i}$ in $\mathbb{P}_{K}^{r}$. If $\bigcap_{i} Z_{i}^{\times}=F$ is nonempty and finite, then $F$ has at most $D:=\prod_{i} d_{i}$ elements.
Proof. Since $\mathbb{A}_{K}^{r}$ is Zariski dense in $\mathbb{P}_{K}$ we have $\bar{Z}_{i} \cap A_{K}^{r}=Z_{i}$; moreover, $\bar{Z}_{i}$ is of pure degree $d_{i}$. Let $W_{0}$ be an irreducible component of $W:=\bigcap \bar{Z}_{i}$ that meets $U$. Since $W_{0}$ is irreducible and $W_{0} \cap U$ is nonempty open in $W_{0}$, it is dense in $W_{0}$. Let $w \in W_{0} \cap U$. Then $w$ is in each of the $\bar{Z}_{i} \cap U$ and so $w \in F$. So $W_{0} \cap U$ is finite and contained in $F$. Since $W_{0} \cap U$ is dense in $W_{0}$, it follows that $W_{0} \cap U$ must consist of a single element of $F$ and $W_{0}$ itself consists of that point. Moreover all points of $F$ are accounted for in this manner as $F$ is contained in the union of irreducible components of $W$ which meet $U$. Hence the cardinality of $F$ is at most the number of irreducible components of $W$, which is at most $D$.

The maps $\pi_{J}$ and a universal bound for the cardinality of their fibers. Let $V \simeq k^{m}$ so that $\boldsymbol{f}=\left(f_{1}, \ldots, f_{r}\right)$ with $f_{j} \in k\left[x_{1}, \ldots, x_{m}\right]$. Assume that $\boldsymbol{c}$ is a noncritical value of $\boldsymbol{f}$ so that $L_{\boldsymbol{c}}$ has no singularities. Fix $J \subset \underline{m}:=\{1, \ldots, m\}$, and let $\pi_{J}: k^{m} \rightarrow k^{m-r}$ map $\left(x_{1}, \ldots, x_{m}\right)$ to $\left(y_{1}, \ldots, y_{m-r}\right)$, where $\left\{y_{j}\right\}_{j=1}^{m-r}=\left\{x_{i} \mid i \in \underline{m} \backslash J\right\}$. We wish to prove that the map $\pi_{J}$ restricted to $L_{c}$ has fibers of cardinality $\leq D:=d_{1} \cdots d_{r}$, where $d_{i}:=\operatorname{deg}\left(f_{i}\right)$. Without loss of generality assume $J=\{1, \ldots, r\}$, so that
$\pi_{J}:\left(x_{1}, \ldots, x_{m}\right) \mapsto\left(x_{r+1}, \ldots, x_{m}\right)$. Write $x=\left(x_{1}, \ldots, x_{m}\right)$ and $y=\left(x_{r+1}, \ldots, x_{m}\right)$. Define $z$ so that $x=(z, y)$.

We regard $L_{\boldsymbol{c}}$ as an affine variety and $L_{c, J}$ as an affine open subvariety. For any $k^{\prime}$ with $k \subset k^{\prime} \subset K$ we have the respective sets of $k^{\prime}$-points, $L_{\boldsymbol{c}}\left(k^{\prime}\right)$ and $L_{\boldsymbol{c}, J}\left(k^{\prime}\right)$. Denote the restriction of $\pi_{J}$ to $L_{c, J}$ by $\bar{\pi}_{J}$.
Proposition 5.5. Let $D=\prod_{1 \leq i \leq r} d_{i}$. Then the fibers of $\bar{\pi}_{J}$ are all of cardinality $\leq D$.
Proof. Note that $d \bar{\pi}_{J}$ is an isomorphism on $L_{\boldsymbol{c}, J}(k)$. Hence $U_{J}(k):=\bar{\pi}_{J}\left(L_{\boldsymbol{c}, J}(k)\right)$ is open in $k^{m-r}$ and $\bar{\pi}_{J}$ is a local analytic isomorphism of $L_{\boldsymbol{c}, J}(k)$ onto $U_{J}(k)$. For any field $k^{\prime}$ between $k$ and $K$, we write again $\bar{\pi}_{J}$ for the map $L_{c, J}\left(k^{\prime}\right) \rightarrow k^{\prime m-r}$, and $U_{J}\left(k^{\prime}\right)$ for its image. If $k^{\prime}$ is a finite extension of $k$, then $k^{\prime}$ is again a local field; exactly as for $k$, we have $d \bar{\pi}_{J}: L_{c, J}\left(k^{\prime}\right) \rightarrow U_{J}\left(k^{\prime}\right)$ is an analytic isomorphism. For any $k^{\prime}, k \subset k^{\prime} \subset K$ with $k^{\prime} / k$ finite, $U_{J}\left(k^{\prime}\right)$ is open in $k^{\prime m-r}$ and the fibers of $\bar{\pi}_{J}$ on $L_{c, j}\left(k^{\prime}\right)$ are discrete and at most countable. If we then fix $y \in U_{J}(k)$, and write $W_{y}$ for the affine variety $\bar{\pi}_{J}^{-1}(y)$, then $W_{y}\left(k^{\prime}\right)$ is at most countable for all finite extensions $k^{\prime} / k$. Hence, by Corollary 5.3, $W_{y}(K)$ is finite. Let $F:=W_{y}(K)$.

On the other hand, $\pi_{J}^{-1}(y)(K)=K^{r} \times\{y\} \simeq K^{r}$. Let $h_{i}(z):=f_{i}(z, y)-c_{i}$. Then $h_{i}$ is a polynomial on $K^{r}$ of degree $\leq d_{i}$. Moreover, since $\bar{\pi}_{J}^{-1}(y)(k)$ is nonempty, $\partial\left(h_{1}, \ldots, h_{r}\right) / \partial\left(x_{1}, \ldots, x_{r}\right)=\partial_{J}(z, y)$ is not identically zero on $K^{r}$. Thus, $\left\{z \mid \partial_{J}(z, y) \neq 0\right\}$ is a nonempty affine open $U_{1}$ in $K^{r}$. Moreover, $F=\bigcap_{1 \leq i \leq r} Z\left(h_{i}\right)^{\times}$where $Z\left(h_{i}\right)^{\times}:=Z\left(h_{i}\right) \cap U_{1}$. So Lemma 5.4 applies and proves that $\# F \leq D$.
Lemma 5.6. Let $\partial_{J}$ be as on page 232. Then if $\omega_{m-r}$ is the exterior form corresponding to the Haar measure on $k^{m-r}$, the exterior form

$$
\rho_{\boldsymbol{c}}:=\frac{1}{\partial_{J}(x)} \bar{\pi}_{J}^{*}\left(\omega_{m-r}\right)
$$

on $L_{\boldsymbol{c}, J}$ has the property that $\left|\rho_{\boldsymbol{c}}\right|$ generates the measure $\mu_{\boldsymbol{c}}:=\mu_{\boldsymbol{f}, \boldsymbol{c}}$. In particular, if $\lambda$ is the Haar measure on $k^{m-r}$ and $v$ is the measure generated by $\left|\bar{\pi}_{J}^{*}\left(\omega_{m-r}\right)\right|$, then $\bar{\pi}_{J}$ takes $v$ to $\lambda$ in small open neighborhoods of each point of $L_{c, J}(k)$, and $d \mu_{c}=\left|\partial_{J}(x)\right|^{-1} d \nu$.
Proof. This is clear from (2-4).
Proof of Theorem 1.2. This follows from three things: the lower bounds on $\left\|\nabla_{r}\right\|$ when $\boldsymbol{c}$ is a stably noncritical value of $\boldsymbol{f}$, the relationship between $\lambda, \nu, \mu_{\boldsymbol{f}, \boldsymbol{c}}$, and the temperedness of $\lambda$. The simple measure-theoretic lemma below explains this. Let $R, S$ be locally compact metric spaces which are second countable, with Borel measures $r, s$ respectively on them, and $\pi: R \rightarrow S$ a continuous surjective map which is a local homeomorphism, and takes $r$ to $s$ in a small neighborhood of each point of $R$ : this means that for each $x \in R$ there are open sets $M_{x}, N_{\pi(x)}$ containing $x$ and $\pi(x)$ respectively, such that $\pi$ is a homeomorphism of $M_{x}$ with $N_{\pi(x)}$ and takes $r$ to $s$.

Lemma 5.7. If there is a natural number $d$ such that all fibers of $\pi$ have cardinality at most $d$, then for each Borel set $E \subset R, \pi(E)$ is a Borel set in $S$, and we have

$$
r(E) \leq d \cdot s(\pi(E))
$$

Moreover if $f \geq 0$ is a continuous function on $R$ and $t$ is the Borel measure on $R$ defined by $d t=f d r$, then for any Borel set $E \subset R$ we have

$$
t(E) \leq \sup _{E}|f| \cdot d \cdot s(\pi(E))
$$

Proof. The second inequality follows trivially from the first, so that we need only prove the first. We use induction on $d$. For $d=1, \pi$ is a continuous bijection of $R$ with $S$; being a local homeomorphism, it is then a global homeomorphism. It is easy to see that it takes $r$ to $s$ globally, and so the results are trivial. Let $d>1$, assume the results for $d-1$, and suppose that there are points of $S$ the fibers over which have cardinality exactly $d$. Let $S_{d}$ be the set of such points in $S$. Now, if the fiber above a point has $e$ elements, the fibers of neighboring points have cardinality $\geq e$, and so $S_{d}$ is open in $S$. Let $R_{d}=\pi^{-1}\left(S_{d}\right)$. Then $\pi: R_{d} \rightarrow S_{d}$ is a $d$-sheeted covering map. If $x \in R_{d}$, we can find an open set $M$ containing $\pi(x)$ such that $N:=\pi^{-1}(M)=\bigsqcup_{1 \leq j \leq d} N_{j}$ where $\pi: N_{j} \rightarrow M$ is a homeomorphism taking $r$ to $s$. If $E \subset N$ is a Borel set, then $E=\bigsqcup_{j} E \cap N_{j}$, so that $\pi(E)=\bigcup_{j} \pi\left(E \cap N_{j}\right)$ is Borel as $\pi$ is a homeomorphism on each $N_{j}$. Moreover,

$$
r(E)=\sum_{j} r\left(E \cap N_{j}\right)=\sum_{j} s\left(\pi\left(E \cap N_{j}\right) \leq d \cdot s(\pi(E))\right.
$$

These two properties are true with any Borel $M^{\prime} \subset M$ and $N^{\prime}=\pi^{\prime}\left(M^{\prime}\right)$ replacing $M, N$ respectively. Write now $S_{d}=\bigcup_{n} M_{n}$ where the $M_{n}$ are open and have the properties described above for $M$. Then $S_{d}=\bigsqcup_{n} M_{n}^{\prime}$ where $M_{n}^{\prime} \subset M_{n}$, so that $R_{d}=\bigsqcup_{n} \pi^{-1}\left(M_{n}^{\prime}\right)$. The two properties above are valid for any Borel set contained in any $\pi^{-1}\left(M_{n}^{\prime}\right)$, hence they follow for any Borel set $E \subset R_{d}$. Write $S^{\prime}=S \backslash S_{d}, R^{\prime}=\pi^{-1}\left(S^{\prime}\right)=R \backslash R_{d}$. Then ( $\left.R^{\prime}, S^{\prime}, \pi\right)$ inherit the properties of $(R, S, \pi)$ with $d-1$ instead of $d$. The result is valid for $\left(R^{\prime}, S^{\prime}, \pi\right)$ and hence for ( $R, S, \pi$ ), as is easily seen.

We are now ready to prove Theorem 1.2. Assume that $\boldsymbol{c}$ is a stably noncritical value of $\boldsymbol{f}$. For simplicity of notation we will suppress mentioning $\boldsymbol{c}$, because all of our estimates are locally uniform in $\boldsymbol{c}$. On $L_{\boldsymbol{c}}=L$ we have the estimate

$$
\left\|\nabla_{r}(x)\right\|=\max _{J}\left|\partial_{J}(x)\right|>\frac{C}{\|x\|^{\gamma}} \quad(\|x\| \geq 1)
$$

where $C>0, \gamma \geq 0$ are constants that remain the same when $\boldsymbol{c}$ is varied in a small neighborhood of $\boldsymbol{c}$. Let us write $L^{+}$for the subset of $L$ where $\|x\|>1$. Now, at
each point $x \in L^{+}$some $\left|\partial_{J}(x)\right|$ equals $\left\|\nabla_{r}(x)\right\|$. Hence if we write

$$
M_{J}=\left\{x \in L^{+}| | \partial_{J}(x) \mid>C\|x\|^{-\gamma}\right\}
$$

then

$$
L^{+}=\bigcup_{J} M_{J}
$$

The map $\bar{\pi}_{J}$ is open on $M_{J}$ onto its image $W_{J}$ and is a local analytic isomorphism. Moreover, if $\lambda, v, \mu=\mu_{\boldsymbol{c}}$ have the same meaning as before, we have, on $M_{J}$,

$$
d \mu=\left|\partial_{J}(x)\right|^{-1} d \nu
$$

and hence, for any Borel set $E \subset M_{J}$, with $D$ as in Lemma 5.4,

$$
\mu(E) \leq D \cdot \sup _{E}\left|\partial_{J}(x)^{-1}\right| \cdot \lambda\left(\pi_{J}(E)\right)
$$

Remembering that $\left|\partial_{J}(x)\right|^{-1}<C^{-1}\|x\|^{\gamma}$, we get from this that

$$
\mu(E) \leq D C^{-1} \cdot \sup _{E}\|x\|^{\gamma} \cdot \lambda\left(\pi_{J}(E)\right) .
$$

If we take $E=B_{R} \cap M_{J}$ where $B_{R}=\left\{x \in k^{m} \mid\|x\|<R\right\}$, we see that $\pi_{J}(E)$ is a subset of the open ball of $k^{m-r}$ of radius $R$, and hence $\lambda\left(\pi_{J}(E)\right) \leq A R^{m-r}$ where $A$ is a universal constant. Hence

$$
\mu\left(B_{R} \cap M_{J}\right) \leq A D C^{-1} \cdot R^{m-r+\gamma} .
$$

Since this is true for all $J$, the temperedness of $\mu$ together with the growth estimate is proved, as well as the assertion that the last estimate remains unchanged if $\boldsymbol{c}$ varies in a small neighborhood of its original value. This finishes the proof of Theorem 1.2.

## 6. Invariant measures on regular adjoint orbits of a semisimple Lie algebra

As an application of our Theorem 1.2 we shall prove that the invariant measures on regular semisimple orbits of a semisimple Lie algebra $\mathfrak{g}:=\mathfrak{g}_{K}$ over a local field $k$ of characteristic 0 are tempered.

The restriction to regular orbits is a consequence of the methods we use; the result is expected to be true without any condition on the orbit of the adjoint action.

For the moment let $k$ be any field of characteristic 0 and $K$ the algebraic closure of $k$. We write $\mathfrak{g}_{K}=K \otimes_{k} \mathfrak{g}_{k}$. Let $P(K)$ be the $K$-algebra of polynomial functions on $\mathfrak{g}_{K}$ with values in $K$. Since such a polynomial is determined by its restriction to $\mathfrak{g}_{k}$, the restriction to $k$ defines an isomorphism of $P(K)$ with the $K$-algebra $P_{k}(K)$ of $K$-valued polynomial functions on $\mathfrak{g}_{k}$.

Let $G$ be the connected adjoint group of $\mathfrak{g}_{k}$. It is a linear algebraic group defined over $k$ and we write $G\left(k^{\prime}\right)$ for the group of its points over $k^{\prime}, k \subset k^{\prime} \subset K$. We
regard $G\left(k^{\prime}\right)$ as a subset of $G=G(K)$. From [Borel 1991] we know that $G(k)$ is Zariski-dense in $G(K)$. Now $G(K)$ acts on $P(K)$ and we denote by $J(K)$ the $K$-algebra of invariants of this action, which is a graded algebra in the obvious way. By a theorem of Chevalley, $J(K)$ is freely generated by homogeneous elements $p_{1}, \ldots, p_{r}$ of degrees $d_{1}, \ldots, d_{r}$ respectively, where $r$ is the rank of $\mathfrak{g}_{k}$. In view of our remarks above, $J(K)$ is isomorphic to the graded $K$-algebra of invariants of $G(k)$ in $P_{k}(K)$. The action by $G(k)$ leaves $P_{k}(k)$ invariant, and we write $J(k)$ for the graded $k$-subalgebra of $G(k)$ - invariants in $P_{k}(k)$. It is clear that

$$
J(k) \simeq J(K)^{\mathrm{Gal}(K / k)}
$$

as graded $k$-algebras.
The following lemma is surely known but we include it for the sake of completeness.

Lemma 6.1. The graded $k$-algebra $J(k)$ is freely generated by homogeneous elements $q_{1}, \ldots, q_{r}$ of degrees $d_{1}, \ldots, d_{r}$ respectively.

Proof. There is a finite extension $k^{\prime}$ of $k$ with $k \subset k^{\prime} \subset K$ such that the free homogeneous generators $p_{i}$ of $J(K)$ have their coefficients in $k^{\prime}$. Hence we may come down from $K$ to $k^{\prime}$. Let $\left(e_{\alpha}\right)$ be a $k$-basis for $k^{\prime}$. Then we can write each $p_{i}$ as

$$
p_{i}=\sum_{\alpha} p_{i, \alpha} e_{\alpha} \quad\left(p_{i, \alpha} \in P_{k}(k)\right)
$$

Since the $p_{i, \alpha}$ are $k$-valued, the $G(k)$-invariance of the $p_{i}$ implies that the $p_{i, \alpha}$ are in $J(k)$. Now the $p_{i}$ are algebraically independent, and so, $\omega:=d p_{1} \wedge \cdots \wedge d p_{r} \not \equiv 0$. Let

$$
\omega_{\alpha_{1}, \ldots, \alpha_{r}}=d p_{1, \alpha_{1}} \wedge \cdots \wedge d p_{r, \alpha_{r}} .
$$

Then

$$
\omega=\sum_{\alpha_{1}, \ldots, \alpha_{r}} \omega_{\alpha_{1}, \ldots, \alpha_{r}} e_{\alpha_{1}} \wedge e_{\alpha_{2}} \cdots \wedge e_{\alpha_{r}} \not \equiv 0
$$

Hence we can choose $\alpha_{1}, \ldots, \alpha_{r}$ such that $\omega_{\alpha_{1}, \ldots, \alpha_{r}} \not \equiv 0$. With this choice, let

$$
q_{i}=p_{i, \alpha_{i}} \quad(1 \leq i \leq r)
$$

Then the $q_{i}$ are homogeneous elements of $J(k)$ and $\operatorname{deg}\left(q_{i}\right)=d_{i}$, and they are algebraically independent.

Now $J\left(k^{\prime}\right)$ is freely generated by the $p_{i}$ of degree $d_{i}$. Hence, its Poincaré series is $\prod_{i}\left(1-T^{d_{i}}\right)^{-1}$. For any integer $m \geq 1$ let $D_{m}$ be the dimension of $J\left(k^{\prime}\right)_{m}$, the subspace of degree $m$ in $J\left(k^{\prime}\right)$. So $\operatorname{dim}\left(J(k)_{m}\right) \leq D_{m}$. On the other hand, let $J_{1}(k)$ be the subalgebra of $J(k)$ generated by the $q_{i}$. Since the $q_{i}$ are homogeneous, this
is a graded subalgebra of $J(k)$, and it has the same Poincaré series as $J\left(k^{\prime}\right)$. Now $J_{1}(k)_{m} \subset J(k)_{m}$ for all $m$, and so

$$
D_{m}=\operatorname{dim} J_{1}(k)_{m} \leq \operatorname{dim} J(k)_{m} \leq \operatorname{dim} J\left(k^{\prime}\right)_{m}=D_{m}
$$

This proves that $J_{1}(k)_{m}=J(k) m$ for all $m$, so that $J_{1}(k)=J(k)$. This finishes the proof of the lemma.

Let $r=\operatorname{rank}(\mathfrak{g})$. Then by assumption we can choose $g_{1}, \ldots, g_{r} \in J(k)$ freely generating $J(k)$, hence also $J(K)$ (over $K$ ). An element $H \in \mathfrak{g}_{K}$ is semisimple (resp. nilpotent) if ad $X$ is semisimple (resp. nilpotent). A semisimple element $H$ is called regular if its centralizer is a Cartan subalgebra (CSA) of $\mathfrak{g}_{K}$. There is an invariant polynomial $D \in J(k)$, called the discriminant of $\mathfrak{g}$, such that if $X \in \mathfrak{g}_{k}, X$ is semisimple and regular if and only if $D(X) \neq 0$. If $Y \in \mathfrak{g}$ is any element, we can write $Y=H+X$ where $H$ is semisimple and $X$ is a nilpotent in the derived algebra of the centralizer of $H$ in $\mathfrak{g}_{K}$ (which is semisimple). It is known [Kostant 1963] that the orbit of $H+X$ has $H$ in its closure, and so, for any $g \in J(K)$, we have $g(H)=g(H+X)$. If $\mathfrak{h}_{K}$ is a CSA of $\mathfrak{g}_{K}$, it is further known that the restriction map from $\mathfrak{g}_{K}$ to $\mathfrak{h}_{K}$ is an isomorphism of $J(K)$ with the algebra $J\left(\mathfrak{h}_{K}\right)$ of polynomials on $\mathfrak{h}_{K}$ invariant under the Weyl group $W_{K}$ of $\mathfrak{h}_{K}$. It is known that the differentials $d g_{1}, \ldots, d g_{r}$ are linearly independent at an element $Y$ of $\mathfrak{g}_{K}$ if and only if $Y$ lies in an adjoint orbit of maximal dimension, which is $\operatorname{dim}\left(\mathfrak{g}_{K}\right)-\operatorname{rank}\left(\mathfrak{g}_{K}\right)=n-r$, where $n=\operatorname{dim}\left(\mathfrak{g}_{K}\right)$ [Kostant 1963]. If $Y$ is semisimple, this happens if and only if $Y$ is regular. Let $\mathfrak{g}_{K}^{\prime}$ be the invariant open set of regular semisimple elements. We write

$$
\boldsymbol{F}=\left(g_{1}, \ldots, g_{r}\right): \mathfrak{g}_{K} \mapsto K^{r}
$$

and view it as a polynomial map of $\mathfrak{g}_{K}$ into $K^{r}$ commuting with the action of the adjoint group. Before we apply Theorem 1.2 to this set up, we need some preliminary discussion. Let $\mathcal{R}=\boldsymbol{F}\left(\mathfrak{g}_{K}^{\prime}\right)$. The next lemma deals with the situation over $K$.
Lemma 6.2. We have $\mathfrak{g}_{K}^{\prime}=\boldsymbol{F}^{-1}(\mathcal{R})$. Moreover $\mathcal{R}$ is Zariski open in $K^{r}$, and is precisely the set of noncritical values of $\boldsymbol{F}$, so that all the noncritical values are also stably noncritical. Moreover, for any $\boldsymbol{c} \in \mathcal{R}$, the preimage $\boldsymbol{F}^{-1}(\boldsymbol{c})$ is an orbit under the adjoint group, consisting entirely of regular semisimple elements, hence smooth.

Proof. Since $d g_{1} \wedge \cdots \wedge d g_{r} \neq 0$ everywhere on $\mathfrak{g}_{K}^{\prime}$, the map $\boldsymbol{F}$ is smooth on $\mathfrak{g}_{K}^{\prime}$. Hence it is an open map [Görtz and Wedhorn 2010, Corollary 14.34], showing that $\boldsymbol{F}\left(\mathfrak{g}_{K}^{\prime}\right)=\mathcal{R}$ is open in $K^{r}$.

We shall prove that if $Y \in \mathfrak{g}_{K}$ and $X \in \mathfrak{g}_{K}^{\prime}$ are such that $\boldsymbol{F}(Y)=\boldsymbol{F}(X)$, then $Y$ is regular semisimple, and is conjugate to $X$ under the adjoint group. Suppose $Y$ is not regular semisimple. Write $Y=Z+N$, where $Z$ is semisimple and $N$ is a nilpotent in the derived algebra of the centralizer of $Z$. The $\boldsymbol{F}(Y)=\boldsymbol{F}(Z)=\boldsymbol{F}(X)$.

Using the action of the adjoint group separately on $X$ and $Z$ we may assume that $X, Z \in \mathfrak{h}_{K}$ where $\mathfrak{h}_{K}$ is a CSA, and $\boldsymbol{F}(X)=\boldsymbol{F}(Z)$. Then all Weyl group invariant polynomials take the same value at $Z$ and $X$ and so $Z$ and $X$ are conjugate under the Weyl group. But as $X$ is regular, so is $Z$, hence $N=0$ or $Y$ itself is regular semisimple. So, $\mathfrak{g}_{K}^{\prime}=\boldsymbol{F}^{-1}(\mathcal{R})$. But then the above argument already shows that $Y$ and $X$ are conjugate under the adjoint group. Since the fibers of $\boldsymbol{F}$ above points of $\mathcal{R}$ are smooth, all points of $\mathcal{R}$ are stably noncritical. It remains to show that there are no other noncritical values. Suppose $Y \in \mathfrak{g}_{K}$ is such that $\boldsymbol{d}=\boldsymbol{F}(Y)$ is a noncritical value where $\boldsymbol{d} \notin \mathcal{R}$. Then $Y \notin \mathfrak{g}_{K}^{\prime}$. Now $Y=Z+N$ as before, where $Z$ is no longer regular (it is semisimple still). Then $\boldsymbol{F}(Z)=\boldsymbol{F}(Y)$ and so $Z \in \boldsymbol{F}^{-1}(\boldsymbol{d})$. But as $Z$ is semisimple but not regular, $d g_{1} \wedge \cdots \wedge d g_{r}$ is zero at $Z$ [Kostant 1963]. Thus $Z$ is a singular point of $\boldsymbol{F}^{-1}(\boldsymbol{d})$, contradicting the fact that $\boldsymbol{d}$ is noncritical. The lemma is thus completely proved.

We now come to the case where the ground field is $k$, a local field of characteristic 0 . We assume that the $g_{i}$ have coefficients in $k$. Fix a regular semisimple element $H_{0}$ in $\mathfrak{g}_{k}$. Let

$$
W(k):=W_{H_{0}}(k)=\left\{X \in \mathfrak{g}(k) \mid g_{i}(X)=g_{i}\left(H_{0}\right)(1 \leq i \leq r)\right\} .
$$

Theorem 6.3. Then the canonical measure on $W(k)$ is tempered, and the growth estimate G (see Section 1) is uniform when $H$ varies in a neighborhood of $H_{0}$.
Proof. For the map $\boldsymbol{F}$ on $\mathfrak{g}_{k}$ we know that $\left(g_{1}\left(H_{0}\right), \ldots, g_{r}\left(H_{0}\right)\right)$ is a stably noncritical value and so the theorem follows at once from Theorem 1.2.

Although $W(K)$ is a single orbit under $G(K)$, this may no longer true over $k$. $W(k)$ is a $k$-analytic manifold of dimension $n-r$. On the other hand, the stabilizer in $G(k)$ of any point of $W(k)$ has dimension $r$ and so its orbit under $G(k)$ is an open submanifold of $W(k)$. If we do this at every point of $W(k)$ we obtain a decomposition of $W(k)$ into a disjoint union of $G(k)$-orbits which are open submanifolds of dimension $n-r$ and so all these submanifolds are closed also. Thus the orbit $G(k) \cdot H_{0}$ is an open and closed submanifold of $W(k)$ of dimension $n-r$. Now the canonical measure on $W(k)$ is invariant under $G(k)$ and so on the orbit $G(k) . H_{0}$ it is a multiple of the invariant measure on the orbit. Note that the orbit being closed, the invariant measure on it is a Borel measure on $\mathfrak{g}_{k}$. Since the canonical measure is tempered on $W(k)$ by Theorem 6.3, it is immediate that the invariant measure on the orbit $G(k) . H_{0}$ is also tempered. Hence we have proved the following theorem:

Theorem 6.4. The orbits of regular semisimple elements of $\mathfrak{g}_{k}$ are closed, and the invariant measures on them are tempered.

For temperedness of invariant measures on semisimple symmetric spaces at the Lie algebra level over $\mathbb{R}$; see [Heckman 1982].

Remark 6.5. Ranga Rao [1972] and Deligne have independently shown that for any $X \in \mathfrak{g}_{k}$, there is an invariant measure on the adjoint orbit of $X$, and this measure extends to a Borel measure on the $k$-closure of the adjoint orbit of $X$. It is natural to ask if these are tempered in our sense when $k$ is nonarchimedean. We shall consider this question in another paper since it does not follow from the results proved here.

## 7. Examples

In this section we give some examples. We consider only single polynomials ( $r=1$ ) of degree $d \geq 3$, defined over a local field $k$ of characteristic 0 . Let $f \in k\left[x_{1}, \ldots, x_{m}\right]$.

Elementary methods when $\boldsymbol{r}=\mathbf{1}$ and $\boldsymbol{f}$ is homogeneous. For $f$ homogeneous we have Euler's theorem on homogeneous functions, which asserts that $\sum_{i} x_{i} \partial f / \partial x_{i}=$ $d \cdot f$. Let $L_{c}=\left\{x \in k^{m} \mid f(x)=c\right\}$ for $c \in k$. Then, for any critical point $x$ of $f$, we have $f(x)=0$, i.e., $L_{0}$ contains all the critical points. So every $c \in k \backslash\{0\}$ is a noncritical value and so is also stably noncritical. Moreover, Euler's identity for $x \in L_{c}, c \neq 0$, gives $\sum_{i} x_{i} \partial f / \partial x_{i}=d c$, so that we have

$$
\left|d \left\|c \left|=\left|\sum_{i} x_{i} \frac{\partial f}{\partial x_{i}}\right| \leq C\|x\|\|\nabla f(x)\| \quad(C>0)\right.\right.\right.
$$

giving the estimate, with $A$ a constant $>0$,

$$
\|\nabla f(x)\| \geq A\|x\|^{-1}, \quad\|x\| \geq 1, x \in L_{c}
$$

Moreover the projection $\left(x_{1}, \ldots, x_{m}\right) \mapsto\left(x_{1}, \ldots, \hat{x}_{i}, \ldots x_{m}\right)$ has the property that all fibers have cardinality $\leq d$. We thus have Theorem 1.2 with

$$
\mu_{f, c}=O\left(R^{m}\right) \quad(R \rightarrow \infty)
$$

where $O$ is uniform locally around $c$. We can actually say more.
Proposition 7.1. Suppose $\mathbf{0}$ is the only singularity in $L_{0}$, i.e., the projective locus of $L_{0}$ is smooth. Then for any compact set $W \subset k \backslash\{\mathbf{0}\}$, we have

$$
\begin{equation*}
\inf _{c \in W, x \in L_{c},\|x\| \geq 1}\|\nabla f(x)\|>0 \tag{7-1}
\end{equation*}
$$

Moreover, the measure $\mu_{f, 0}$ defined on $L_{0} \backslash\{0\}$ is finite in open neighborhoods of $\mathbf{0}$ if $m>d$, so that it extends to a Borel measure on $L_{0}$. Finally, for all $c \in k$,

$$
\mu_{f, c}\left(B_{r}\right)=O\left(R^{m-1}\right)
$$

If $m \leq d$, there are examples where $\mu_{f, 0}$ is not finite in neighborhoods of $\mathbf{0}$.

Proof. To prove (7-1) assume (7-1) is not true. Then we can find sequences $c_{n} \in W$, $x_{n} \in L_{c_{n}}$ such that $c_{n} \rightarrow c \in W, \nabla f\left(x_{n}\right) \rightarrow 0$ as $n \rightarrow \infty$. By passing to a subsequence and permuting the coordinates we may assume that $x_{n}=\left(x_{n 1}, \ldots, x_{n m}\right)$ where $\left|x_{n 1}\right| \geq\left|x_{n j}\right|(j \geq 2)$ and $\left|x_{n 1}\right| \rightarrow \infty$. Now,

$$
f\left(x_{n 1}, \ldots, x_{n m}\right)=x_{n 1}^{d} f\left(1, x_{n 1}^{-1} x_{n 2}, \ldots, x_{n 1}^{-1} x_{n m}\right)=c_{n} \rightarrow c
$$

and

$$
(\nabla f)\left(x_{n 1}, \ldots, x_{n m}\right)=x_{n 1}^{d-1}(\nabla f)\left(1, x_{n 1}^{-1} x_{n 2}, \ldots, x_{n 1}^{-1} x_{n m}\right) \rightarrow 0
$$

Now $\left|x_{n 1}^{-1} x_{n j}\right| \leq 1$ for $2 \leq j \leq m$ and so, passing to a subsequence, we may assume that $x_{n 1}^{-1} x_{n j} \rightarrow v_{j}$ for $j \geq 2$. Hence,

$$
f\left(1, v_{2}, \ldots, v_{m}\right)=0 \quad \text { and } \quad(\nabla f)\left(1, v_{2}, \ldots, v_{m}\right)=0
$$

showing that $\left(1, v_{2}, \ldots, v_{m}\right) \neq(0, \ldots, 0)$ is a singularity of $L_{0}$. Then (7-1) leads to the conclusion

$$
\mu_{f, c}\left(B_{R}\right)=O\left(R^{m-1}\right) \quad(R \rightarrow \infty)
$$

locally uniformly at each $c \neq 0$.
For $\mu_{f, 0}$ defined on $L_{0} \backslash\{0\}$, one must first show that it is finite on small neighborhoods of 0 , i.e., it extends to a Borel measure on $L_{0}$, if $m>d$. Let $S=\left\{u \in L_{0} \mid\|u\|=1\right\}$. Then there exist constants $a, b>0$ such that $a \leq\|\nabla f(x)\| \leq b$ for all $x \in S$. Hence, by homogeneity,

$$
a\|x\|^{d-1} \leq\|\nabla f(x)\| \leq b\|x\|^{d-1} \quad\left(x \in L_{0} \backslash\{0\}\right)
$$

Hence

$$
\|\nabla f(x)\| \geq a>0 \quad\left(x \in L_{0},\|x\| \geq 1\right)
$$

As before, this leads to $\mu_{f, 0}\left(B_{R} \backslash B_{1}\right)=O\left(R^{m-1}\right)$ as $R \rightarrow \infty$. Around $\mathbf{0}$ we obtain the finiteness of $\mu_{f, 0}$ from the estimate $b^{-1}\|x\|^{-(d-1)} \leq\|\nabla f(x)\|^{-1} \leq a^{-1}\|x\|^{-(d-1)}$ and the fact that

$$
\int_{x \in k^{m-1}, 0<\|x\|<1}\|x\|^{-(d-1)} d^{m-1} x<\infty
$$

if $m>d$ for both $k=\mathbb{R}$ and $k$ nonarchimedean. We shall now suppose that $f=X^{4}+Y^{4}-Z^{4}$. Then $\mathbf{0}$ is the only critical point. The map $(x, y, z) \mapsto(x, y)$ on $L_{0} \cap\{(x, y, z) \mid x>0\}$ is a diffeomorphism and the measure $\mu_{f, 0}$ is

$$
\frac{1}{|\partial f / \partial z|} d x d y=\frac{1}{4} \frac{d x d y}{\left(x^{4}+y^{4}\right)^{3 / 4}}
$$

and it is easy to verify that

$$
\iint_{N} \frac{d x d y}{\left(x^{4}+y^{4}\right)^{3 / 4}}=\infty
$$

for any neighborhood $N$ of $(0,0)$.

Remark 7.2. It follows from Proposition 7.1 that to have

$$
\begin{equation*}
\inf _{x \in L_{c},\|x\| \geq 1}\|\nabla f(x)\|=0 \quad(c \neq 0) \tag{7-2}
\end{equation*}
$$

we must look for $f$ such that $L_{0}$ has singular points $\neq \mathbf{0}$. In the next section we describe some of these examples.
Some hypersurfaces in $\mathbb{P}_{\boldsymbol{k}}^{\boldsymbol{m}-1}$ with $[1: 0: \ldots: 0]$ as an isolated singularity. We do not try to give a "normal form" for such hypersurfaces; nevertheless large families of these can be described. We work in $k^{m}, k$ a local field of characteristic 0 . Since the first coordinate axis in $k^{m}$ is chosen to be an isolated critical line (ICL), the first variable will be distinguished in what follows. Let us write $X, Y_{1}, \ldots, Y_{m-1}$ as the variables. Write $\boldsymbol{Y}=\left(Y_{1}, \ldots, Y_{m-1}\right)$. Let $C(\varepsilon)=\{(X, \boldsymbol{Y})|\|\boldsymbol{Y}\| \leq \varepsilon| X \mid\}$
Lemma 7.3. Suppose $\left(X_{n}, \boldsymbol{Y}_{n}\right)$ is a sequence of points in $L_{c}(c \neq 0)$ such that they are in $C(\varepsilon)$ for some $\varepsilon<1$. Let $F\left(X_{n}, \boldsymbol{Y}_{n}\right)=c \neq 0$ and $\nabla F\left(X_{n}, \boldsymbol{Y}_{n}\right) \rightarrow 0$. Then if the $X$-axis is an (ICL) for $F$, we must have $X_{n} \rightarrow \infty, 1 / X_{n} \boldsymbol{Y}_{n} \rightarrow \mathbf{0}$ as $n \rightarrow \infty$.

Proof. By Euler's theorem, there is no singularity on $L_{c}(c \neq 0)$. Hence $\|\nabla F\|$ is bounded away from 0 on each compact subset of $L_{c}$. Hence, item 2 above implies $\left\|\left(X_{n}, \boldsymbol{Y}_{n}\right)\right\|=\left|X_{n}\right| \rightarrow \infty$. Then $\left\|X_{n}{ }^{-1} \boldsymbol{Y}_{n}\right\| \leq 1$ and has a limit point $\boldsymbol{\eta}$. Passing to a subsequence, if necessary, we have $X_{n}{ }^{-1} \boldsymbol{Y}_{n} \rightarrow \boldsymbol{\eta}$ as $n \rightarrow \infty$. If $d=\operatorname{deg}(F)$ we have $X_{n}^{d} F\left(1, X_{n}{ }^{-1} \boldsymbol{Y}_{n}\right)=c, X_{n}^{d-1} \partial_{X} F\left(1, X_{n}{ }^{-1} \boldsymbol{Y}_{n}\right) \rightarrow 0$, and $X_{n}^{d-1} \partial_{Y_{i}} F\left(1, X_{n}{ }^{-1} \boldsymbol{Y}_{n}\right) \rightarrow 0$. So $F(1, \boldsymbol{\eta})=0$ and $\nabla F(1, \eta)=0$, while $\eta \in C(\varepsilon)$. Hence $\boldsymbol{\eta}=\mathbf{0}$ since $\varepsilon$ can be arbitrarily small.

Lemma 7.4. If $(1, \mathbf{0})$ is a critical point of $F$, then $F$ has the form

$$
F=X^{d-2} p_{2}+X^{d-3} p_{3}+\cdots+p_{d}
$$

where $p_{r}$ is a homogeneous polynomial in $\boldsymbol{Y}$ of degree $r$.
Proof. Write $F=X^{d-2} p_{2}+X^{d-3} p_{3}+\cdots+p_{d}$. Then $p_{0}$ is a constant, and $F(1, \mathbf{0})=\mathbf{0}$ gives $p_{0}=0$. Then, $\partial F / \partial Y_{i}(1, \mathbf{0})=0$ gives $p_{1}=0$.

From now on we let $d \geq 3$ and write

$$
F=X^{d-2} p_{2}+\cdots+p_{d}, \quad G=p_{2}+\cdots+p_{d}
$$

Note that $G$ is a polynomial in $\boldsymbol{Y}$, but not necessarily homogeneous.
Lemma 7.5. If $\mathbf{0}$ is an isolated critical point (ICP) of $G$, then the $X$-axis is an ICL of $F$. In particular, this is so if the quadratic form $p_{2}$ is nondegenerate.
Proof. We must prove that if $\left(1, \boldsymbol{Y}_{n}\right)$ is a CP for $F$ with $\boldsymbol{Y}_{n} \rightarrow \mathbf{0}$, then $\boldsymbol{Y}_{n}=0$ for $n \geq 1$. The conditions for $\left(1, \boldsymbol{Y}_{n}\right)$ to be a CP of $F$ are

$$
F\left(1, \boldsymbol{Y}_{n}\right)=0, \quad \frac{\partial}{\partial X} F\left(1, \boldsymbol{Y}_{n}\right)=0, \quad \frac{\partial}{\partial \boldsymbol{Y}_{i}} F\left(1, \boldsymbol{Y}_{n}\right)=0 \text { for all } i .
$$

Consequently $G\left(\boldsymbol{Y}_{n}\right)=0$ and $\partial G / \partial \boldsymbol{Y}_{i}\left(\boldsymbol{Y}_{n}\right)=0$ for all $i$. Since $\boldsymbol{Y}_{n} \rightarrow \mathbf{0}$ and $\mathbf{0}$ is an ICP for $G, \boldsymbol{Y}_{n}=0$ for all $n \gg 1$.

For the second statement, suppose $p_{2}$ is nondegenerate. By Morse's lemma [Duistermaat 1973] for local fields $k, \mathrm{ch} .=0$, there is a local diffeomorphism of $k^{m-1}$ fixing $\mathbf{0}$ taking $G$ to $p_{2}$. But $\mathbf{0}$ is an isolated CP for $p_{2}$, which makes it isolated for $G$.

We remark that Duistermaat's proof [1973] of Morse's lemma is over $\mathbb{R}$, but its proof applies to the nonarchimedean case without any change, so we omit it.
Lemma 7.6. The converse to the first statement of Lemma 7.5 is true if

$$
F=X^{d-r} p_{r}+p_{d} \quad(r \geq 2)
$$

Proof. We must show that $G=p_{r}+p_{d}$ has $\mathbf{0}$ as an ICP if $(1, \mathbf{0})$ is an ICP for $F$. Suppose $\boldsymbol{w}_{n}$ are CPs for $G=p_{r}+p_{d}$ with $\boldsymbol{w}_{n} \rightarrow \mathbf{0}$. Then $G\left(w_{n}\right)=F\left(1, \boldsymbol{w}_{n}\right)=0$ for all $n$, and $G_{i}\left(\boldsymbol{w}_{n}\right)=\partial F / \partial Y_{i}\left(1, \boldsymbol{w}_{n}\right)=0$ for all $n$. Hence, $p_{r, i}\left(\boldsymbol{w}_{n}\right)+p_{d, i}\left(\boldsymbol{w}_{n}\right)=0$ for all $n$. By Euler's theorem, $r p_{r}\left(\boldsymbol{w}_{n}\right)+d p_{d}\left(\boldsymbol{w}_{n}\right)=0$ for all $n$. But, $p_{r}\left(\boldsymbol{w}_{n}\right)+p_{d}\left(\boldsymbol{w}_{n}\right)=0$ for all $n$ as well. So, $p_{r}\left(\boldsymbol{w}_{n}\right)=p_{d}\left(\boldsymbol{w}_{n}\right)=0$ for all $n$. Hence, $\partial F / \partial X\left(1, \boldsymbol{w}_{n}\right)=$ $(d-r) p_{r}\left(\boldsymbol{w}_{n}\right)=0$ for all $n$. So $\left(1, \boldsymbol{w}_{n}\right)$ is a CP of $F$ for all $n$. As $(1, \boldsymbol{0})$ is assumed to be an ICP for $F, \boldsymbol{w}_{n}=\mathbf{0}$ for $n \gg 1$. So $\mathbf{0}$ is an ICP for $F$.

Study of condition (7-2) for $F=X^{d-2} p_{2}+p_{d}$ where $G=p_{2}+p_{d}$ has 0 as an ICP. Let us consider $F=X^{2}+P_{4}(Y)$ where $P_{4}$ is a homogeneous quartic polynomial in $Y, Z$. For this to have $(t, 0,0)$ as and ICL we must have $(0,0)$ as an ICP for $G=Z^{2}+P_{4}(Y, Z)$.
Lemma 7.7. $G=Z^{2}+P_{4}(Y, Z)$ has $\mathbf{0}$ as an ICP if and only if $Z^{2} \nmid P_{4}(Y, Z)$, i.e.,

$$
P_{4}(Y, Z)=a_{0} Y^{4}+a_{1} Y^{3} Z+a_{2} Y^{2} Z^{2}+a_{3} Y Z^{3}+a_{4} Z^{4}
$$

where at least one of $a_{0}, a_{1}$ is nonzero. In this case $\mathbf{0}$ is its only CP.
Proof. The equations which determine whether $(y, z)$ is a CP of $G$ are

$$
z^{2}+P_{4}(y, z)=0, \quad \frac{\partial P_{4}}{\partial Y}(y, z)=0 \quad \text { and } \quad 2 Z+\frac{\partial P_{4}}{\partial Z}(y, z)=0
$$

From the second and third equations just defined, using Euler's theorem, we have $2 z^{2}+4 P_{4}(y, z)=0$, which implies $z^{2}=0$ and $P_{4}(y z)=0$.

So the only critical points are of the form $(y, 0)$. Then $(0,0)$ is certainly a CP. If $(y, 0)$ is a critical point for some $y \neq 0$, then $4 a_{0} y^{3}=0, a_{1} y^{3}=0$ which implies $a_{0}, a_{1}$ both vanish. The entire $Y$-axis consists of critical points, and so for $(0,0)$ to be an ICP, at least one of $a_{0}, a_{1} \neq 0$. in which case $(0,0)$ is the only CP.

We consider the cases (I) $a_{0} \neq 0$ and (II) $a_{0}=0, a_{1} \neq 0$. We consider case (I). We shall now verify that $\inf _{\|\boldsymbol{u}\|>1}\|\nabla F(\boldsymbol{u})\|>0$ if $\boldsymbol{u} \in L_{c},\|\boldsymbol{u}\| \geq 1$. Assume $F=X^{2} Z^{2}+P_{4}(Y, Z)$, and in view of Lemma 7.3, choose a sequence $\left(x_{n}, y_{n}, z_{n}\right)$ such that $x_{n} \rightarrow \infty, y_{n} / x_{n} \rightarrow 0, z_{n} / x_{n} \rightarrow 0$ and:
(i) $x_{n}^{2} y_{n}^{2}+P_{4}\left(y_{n}, z_{n}\right)=c$,
(ii) $\partial F / \partial X=2 x_{n} z_{n}^{2} \rightarrow 0$,
(iii) $\partial P_{4} / \partial Y\left(y_{n}, z_{n}\right) \rightarrow 0$,
(iv) $2 x_{n}^{2} z_{n}+\partial P_{4} / \partial Z\left(y_{n}, z_{n}\right) \rightarrow 0$.

From (ii) we get $z_{n} \rightarrow 0$. Assuming we are in case (I), $y_{n}$ is bounded. Otherwise, by passing to a subsequence we may assume $y_{n} \rightarrow \infty$ giving $\partial P_{4} / \partial Y\left(y_{n}, z_{n}\right)=$ $4 a_{0} y_{n}^{3}+3 a_{1} y_{n}^{2} z_{n}+\cdots \rightarrow 0$. If $a_{0} \neq 0$, then $\partial P_{4} / \partial Y\left(y_{n}, z_{n}\right)=4 a_{0} y_{n}^{3}\left(1+o\left(z_{n} / y_{n}\right)\right) \rightarrow$ $\infty$, which is a contradiction. But if $\eta \neq 0$ is a limit point of $y_{n}$, then

$$
\frac{\partial P_{4}}{\partial Y}\left(z_{n}, y_{n}\right) \rightarrow 4 a_{0} \eta^{3} \neq 0
$$

which is a contradiction. So, $y_{n} \rightarrow 0$ necessarily. Then, $P_{4}\left(y_{n}, z_{n}\right) \rightarrow 0$ and $\partial P_{4} / \partial Z\left(y_{n}, z_{n}\right) \rightarrow 0$. Hence by (iv), $x_{n}^{2} z_{n} \rightarrow 0$, by (i) $x_{n}^{2} z_{n}^{2} \rightarrow c \neq 0$, a contradiction. This finishes case (I).

Assuming we are in case (II), $a_{0}=0, a_{1} \neq 0$, we claim $y_{n} \rightarrow \infty$. Otherwise, by passing to a subsequence, we may assume $y_{n} \rightarrow \eta$. Then $P_{4}\left(y_{n}, z_{n}\right)=a_{1} y_{n}^{3} z_{n}+\cdots$ so that $P_{4}\left(y_{n}, z_{n}\right) \rightarrow 0$. Hence, $x_{n}^{2} z_{n}^{2} \rightarrow c$. But $\partial P_{4} / \partial Z\left(y_{n}, z_{n}\right)=a_{1} y_{n}^{3}+\cdots \rightarrow a_{1} \eta^{3}$. Hence, by (iv), $x_{n}^{2} z_{n}=o(1)$. So, as $z_{n} \rightarrow 0$, we have $x_{n}^{2} y_{n}^{2} \rightarrow 0$. Hence, $c=0$ is a contradiction.

We are left with the case $x_{n} \rightarrow \infty, y_{n} \rightarrow \infty, z_{n} \rightarrow 0,\left(y_{n} / x_{n}\right)\left(z_{n} / x_{n}\right) \rightarrow 0$, and $P_{4}(Y, Z)=a_{1} Y^{3} Z+\cdots$, for $a_{1} \neq 0$. But $\partial P_{4} / \partial Y\left(y_{n}, z_{n}\right)=3 a_{1} y_{n}^{2} z_{n}\left(1+o\left(z_{n} / y_{n}\right)\right) \rightarrow$ 0 if and only if $y_{n}^{2} z_{n} \rightarrow 0$. In this case may we have a counterexample to statement (7-2). Remark 4.14 gives an example of this kind. Note that case (I) is generic among the families we consider.

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