# Orbits on Homogeneous Spaces of Arithmetic Origin and Approximations 

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#### Abstract

. We prove an $S$-arithmetic version, in the context of algebraic groups defined over number fields, of Ratner's theorem for closures of orbits of subgroups generated by unipotent elements. We apply this result in order to obtain a generalization of results of Margulis and of Borel-Prasad about values of irrational quadratic forms at integral points to the general setting of hermitian forms over division algebras with involutions of first or second kind. As a byproduct of our considerations we obtain another proof of the strong approximation theorem for algebraic groups defined over number fields.


## Introduction

Many problems from number theory and, in particular, in Diophantine approximations can be reformulated in terms of dynamics of actions of subgroups on homogeneous spaces. In this way, ideas and methods from dynamical systems can be successfully used in number theory and, vice versa, problems from number theory stimulate the study of certain kinds of dynamical systems. One of the most impressing example is provided by the following conjecture formulated by Oppenheim in 1929 [Op 1,2]: If $f$ is a real nondegenerate indefinite quadratic form of $n \geq 5$ variables and $f$ is not a multiple of a quadratic form with rational coefficients, then for any real $\varepsilon>0$ there exists a nonzero vector $z \in \mathbb{Z}^{n}$ such that $|f(z)|<\varepsilon$. Let us briefly recall some of the conjectures and the subsequent results connected with Oppenheim's conjecture. (For detail, we refer the reader to the exhaustive review paper [M6], as well as to the earlier review papers [B2],[D4],[M5],[Rat4] and [S2].) In the

[^0]mid-seventies, M.S.Raghunathan formulated a conjecture that the closure of an orbit of a unipotent subgroup on a homogeneous space $G / \Gamma$, where $G$ is a Lie group and $\Gamma$ is a lattice in $G$, is an orbit of a larger subgroup, and also noted that this statement implies the Oppenheim conjecture. In 1981, S.G.Dani [D2] formulated the metric version of the Raghunathan conjecture, often called "measure rigidity" : any $U$ invariant $U$-ergodic Borel probability measure on $G / \Gamma$, where $U$ is a unipotent subgroup of $G$, coincides with the Haar measure on a closed orbit of a connected subgroup containing $U$. In [M2,3], using the homogeneous space approach, G.A.Margulis proved that $f\left(\mathbb{Z}^{n}\right)$ is dense in $\mathbb{R}$ which, in particular, confirms Oppenheim's conjecture . The conjectures of Dani and Raghunathan have been proved in full generality by M.Ratner in [Rat1] and [Rat2], respectively. (We will refer to the first result as to "the measure classification theorem" and to the second one as to "the theorem for orbit closures.") The main part in [Rat2] is dedicated to the proof of a theorem about uniform distribution of unipotent flows on the homogeneous space $G / \Gamma$. Later Dani and Margulis [DM2] applied other methods in order to prove a refined version of Ratner's uniform distribution theorem.

In [ BPr$]$ Borel and Prasad obtained the following generalization of the Margulis theorem. Let $S$ be a finite set of normalized valuations of a number field $K$ containing the set $S_{\infty}$ of archimedean ones, $K_{S}$ the direct sum of the completions $K_{v}$ of $K$ at $v \in S$, and $\mathcal{O}$ the ring of $S$-integers of $K$. Let $f_{S}$ be a nondegenerate quadratic form on $K_{S}^{n}$; equivalently, let $f_{S}$ be a collection $f_{v}, v \in S$, where $f_{v}$ is a nondegenerate quadratic form on $K_{v}^{n}$. Assume that $f_{S}$ is $K$-irrational (i.e. $f_{S}$ is not a multiple of a quadratic form on $K^{n}$ ), $f_{v}$ is isotropic for all $v \in S$ and $n \geq 3$. Then $f_{S}\left(\mathcal{O}^{n}\right)$ is dense in $K_{S}$. (Note that $K$ (respectively, $\mathcal{O}^{n}$ ) is diagonally embedded in $K_{S}$ (respectively, in $K_{S}^{n}$ ). Under these embeddings, $K_{S}$ is a $K$-algebra and $\mathcal{O}^{n}$ is a lattice in $K_{S}^{n}$.) The BorelPrasad result can be regarded as an analog for the irrational quadratic forms of the local-global principle for quadratic forms over number fields (that is, of the Hasse-Minkowski theorem [Se, ch. 4, Theorem 8]). In fact, Landherr, Kneser and Springer proved the local-global principle for $K$-rational hermitian forms over finite-dimensional division algebras with involutions of first or second kind(cf. [L], [K1] and [Sch,ch. 10]). This suggests the problem about the extension of the results [M2,3] and [ BPr ] to the general framework of the hermitian forms over division algebras. In [BPr] and [Pr2] Borel and Prasad raised the problem of generalizing Ratner's theorems to the case when $G$ is a finite direct product of real and $p$-adic Lie groups. The generalizations of these theorems were obtained by Ratner herself [Rat3] and, independently,
the generalization of the measure classification theorem was obtained by Margulis and the author [MTo1,2].

Note that for the arithmetic applications of the theorem for orbit closures we need only the $S$-arithmetic version of this theorem when $\mathbf{G}$ is a $K$-algebraic group, $G=\mathbf{G}\left(K_{S}\right)$ and $\Gamma$ is an $S$-arithmetic subgroup of $G$ in the sense that $\Gamma$ and $\mathbf{G}(\mathcal{O})$ are commensurable subgroups of $\mathbf{G}(K)$. (Here and further on $S$ and $K_{S}$ are as above and $G$ is identified with the direct product $\prod_{v \in S} \mathbf{G}\left(K_{v}\right)$.) The proof of the theorem in the $S$-arithmetic case allows to avoid some technical complications although the main ideas from the general case remain involved. In the present paper we give a proof of the theorem for orbit closures in the $S$-arithmetic case using the approach and methods from Dani-Margulis paper [DM2] and, subsequently, we apply this theorem in order to obtain the analogs for the hermitian forms over division algebras of the result of Borel and Prasad [ BPr ]. Our versions in the $S$-arithmetic case of both the theorem for orbit closures and the measure classification theorem (see Theorem 1 and Theorem 2 below) are somewhat more precise than in the general case of direct products of real and $p$-adic Lie groups (cf.[Rat3], [MTo1,2]). Some of our arguments are applied in order to give a short proof of the strong approximation theorem for the simply connected algebraic groups defined over number fields (see [Pl] and also [Pr1] where the strong approximation theorem is proved for any global field $K$ ).

Let us fix the following notions. We say that a connected $K$-algebraic subgroup $\mathbf{P}$ of $\mathbf{G}$ is a subgroup of class $\mathcal{F}$ (relatively to $S$ ) if for each proper normal $K$-algebraic subgroup $\mathbf{Q}$ of $\mathbf{P}$ there exists $v \in S$ such that $(\mathbf{P} / \mathbf{Q})\left(K_{v}\right)$ contains a unipotent element different from the identity. Recall that according to [B1], an $S$-arithmetic subgroup $\Gamma$ of $G$ is a lattice (i.e. $\Gamma$ has finite covolume in $G$ ) if and only if the connected component of $\mathbf{G}$ does not admit nontrivial $K$-rational characters. In the latter case, $\Gamma$ is called $S$-arithmetic lattice. Note that if $\mathbf{P}$ is a subgroup of class $\mathcal{F}$ in $\mathbf{G}$ then $P^{\prime} \cap \Gamma$ is an $S$-arithmetic lattice in $P^{\prime}$ for any subgroup of finite index $P^{\prime}$ in $\mathbf{P}\left(K_{S}\right)$ and any $S$-arithmetic subgroup $\Gamma$ in $G$. Given a subgroup $H \subset G$, we will denote by $H_{u}$ the subgroup generated by all 1-parameter unipotent subgroups of $H$ (see 1.5).

Theorem 1. Let $\mathbf{G}$ be a $K$-algebraic group, $\Gamma$ an $S$-arithmetic lattice in $G, H$ a subgroup of $G$ such that $H=H_{u}$ and $x=g \Gamma$ a point in $G / \Gamma$. Then there exists a subgroup $\mathbf{P} \subset \mathbf{G}$ of class $\mathcal{F}$ and a subgroup of finite index $P^{\prime}$ in $\mathbf{P}\left(K_{S}\right)$ such that $g P^{\prime} g^{-1}$ contains $H$ and the closure of $H x$ in $G / \Gamma$ coincides with $g P^{\prime} g^{-1} x$.

In sections 3 and 4 we will give a direct proof of Theorem 1. (It can also be deduced from [Rat3, Theorem 2] and Theorem 3 below.) The proof of Theorem 1 essentially uses the following measure classification theorem, very important by itself.

Theorem 2. Let $\mathbf{G}$ be a $K$-algebraic group, $\Gamma$ an $S$-arithmetic subgroup of $G, H$ a subgroup of $\mathbf{G}$ such that $H=H_{u}$ and $\mu$ an $H$ invariant $H$-ergodic Borel probability measure on $G / \Gamma$. Then there exist a subgroup $\mathbf{P} \subset \mathbf{G}$ of class $\mathcal{F}$, a subgroup of finite index $P^{\prime}$ in $\mathbf{P}\left(K_{S}\right)$ and a point $x=g \Gamma$ in $G / \Gamma$ such that $g P^{\prime} g^{-1}$ contains $H, g P^{\prime} g^{-1} x$ is closed in $G / \Gamma$ and the measure $\mu$ is $g P^{\prime} g^{-1}$-invariant and concentrated on $g P^{\prime} g^{-1} x$.

Recall that in the usual formulations of the above theorems (see [Rat3] and [MTo1]) $P^{\prime}$ is a closed subgroup of $G$ such that $P^{\prime} \cap \Gamma$ has finite covolume in $P^{\prime}$ (without the additional specification that $P^{\prime}$ has finite index in $\mathbf{P}\left(K_{S}\right)$, where $\mathbf{P}$ is a subgroup of class $\left.\mathcal{F}\right)$. Theorem 2 follows from [MTo1,Theorem 2] or [Rat3. Theorem 1] and from the next theorem.

Theorem 3. Let $\mathbf{G}$ be a $K$-algebraic group, $\Gamma$ an $S$-arithmetic subgroup of $G, M$ a closed subgroup of $G, x=g \Gamma$ a point in $G / \Gamma$ and $M_{x}=\{a \in G \mid a x=x\}$. Assume that $M x$ is closed, $M_{u} x$ is dense in $M x$ and $M x$ admits $M$-invariant Borel probability measure $\mu$. Let $P=\mathbf{P}\left(K_{S}\right)$ where $\mathbf{P}$ is the connected component of the Zariski closure of $g^{-1} M_{x} g$ in $\mathbf{G}$. Then $P^{\prime}=\left\{a \in P \mid a g^{-1} \mu=g^{-1} \mu\right\}$ is a subgroup of finite index in $P, g P^{\prime} g^{-1}=M x$, and $g^{-1} M g \cap P$ is an open subgroup in $P^{\prime}$ containing $g^{-1} M_{u} g$. Furthermore, $\mathbf{P}$ is an algebraic subgroup of class $\mathcal{F}$ in $\mathbf{G}$ and it is uniquely defined by $\mu$.

Theorem 3 will be proved in section 2. Also in section 2, we easily derive from Theorem 3 the following strong approximation theorem.

Theorem 4. (cf. [Pl]) Let $\mathbf{G}$ be a connected, simply connected, algebraic group defined over a number field $K$. Let $\mathcal{V}$ be the adele ring of $K$ and $T$ be a finite set of normalized valuations of $K$. Assume that for any proper $K$-algebraic subgroup $\mathbf{N}$ of $\mathbf{G}$ there exists a valuation $v \in T$ such that $(\mathbf{G} / \mathbf{N})\left(K_{v}\right)$ is not compact. Then $\mathbf{G}\left(K_{T}\right) \mathbf{G}(K)$ is dense in $\mathbf{G}(\mathcal{V})$.

Note that the group $\mathbf{G}$ in the formulation of Theorem 4 is actually a group of class $\mathcal{F}$. Indeed, since $\mathbf{G}$ is a simply connected algebraic group, the solvable radical of $\mathbf{G}$ coincides with its unipotent radical. This
implies that $\mathbf{G}$ is of class $\mathcal{F}$ relatively to any finite set $S$ of valuations of $K$ containing $T$ and $S_{\infty}$.

In order to formulate the results about the hermitian forms, we need to fix some standard algebraic notions and concepts (cf. [Sch.], ch. 8,10$]$ ). We denote by $D$ a central division algebra over a number field $L$ and of finite degree $r$ i.e. $\operatorname{dim}_{L} D=r^{2}$. We fix a subfield $K$ of $L$ such that either $L=K$ or $L$ is a quadratic extension of $K$. We let $S, S_{\infty}, \mathcal{O}, K_{v}$ and $K_{S}$ be the same as before. In addition, we denote by $\Lambda$ an $\mathcal{O}$-order in $D$ (i.e. $\Lambda$ is an $\mathcal{O}$-algebra of finite type such that $\left.D=\Lambda \bigotimes_{\mathcal{O}} K\right)$. Tensoring with $K_{v}$ gives the topological $K_{v}$-algebras $D_{v}=D \bigotimes_{K} K_{v}$ and $L_{v}=L \bigotimes_{K} K_{v}$. Let $T$ be a subset of $S$. The direct sums $D_{T}=\bigoplus_{v \in T} D_{v}, L_{T}=\bigoplus_{v \in T} L_{v}$ and $K_{T}=\bigoplus_{v \in T} K_{v}$ are endowed with the product topology. We will identify $D, L$ and $K$ with their diagonal embeddings in $D_{T}, L_{T}$ and $K_{T}$, respectively. It is well known that each of these embeddings is dense. Furthermore, $\Lambda$ is a discrete cocompact abelian subgroup of $D_{S}$. Note that $L_{S}$ (resp. $L_{v}$ ) coincides with the center of $D_{S}\left(\right.$ resp. $\left.D_{v}\right)$.Let $v \in S$ and $\tau_{v}$ be a $L_{v} / K_{v}$-involution on $D_{v}$, i.e. $\tau_{v}$ is an antiautomorphism on $D_{v}$ such that $\tau_{v}^{2}=i d$ and $K_{v}=\left\{\left.x \in L_{v}\right|^{\tau_{v}} x=x\right\}$. Clearly, $\tau_{T}=\bigoplus_{v \in T} \tau_{v}$ is a $L_{T} / K_{T}$-involution on $D_{T}$ for any $T \in S$ and, conversely, every $L_{T} / K_{T^{-}}$ involution on $D_{T}$ is a direct sum of $L_{v} / K_{v}$-involution on $D_{v}, v \in T$. Note that any $L / K$-involution $\tau$ on $D$ extends in a unique way to a $L_{v} / K_{v}$-involution on $D_{v}$ (resp. $L_{T} / K_{T}$-involution on $D_{T}$ ). If $L=K$ (resp. $L$ is a quadratic extension of $K$ ), then $\tau_{v}, \tau_{T}$ and $\tau$ are involutions of first (resp. second) kind (cf.[Sch]).

Let $\lambda_{T}=\left(\lambda_{v}\right)_{v \in T} \in L_{T}$ and $h_{T}$ be a nondegenerate $\lambda_{T}$-hermitian form on $D_{T}^{n}$ (see 5.1). Then $h_{T}$ can be equivalently viewed as a collection $h_{v}, v \in T$, where $h_{v}$ is a $\lambda_{v}$-hermitian form (with respect to $\tau_{v}$ ) on $D_{v}^{n}$. (Note that if $\tau_{T}=\mathrm{id}$ and $\lambda_{T}=1$ (resp. $\lambda_{T}=-1$ ) then $D_{T}=K_{T}$ and each $h_{v}, v \in T$, is a bilinear symmetric (resp. symplectic) form. In the first case $h_{v}(x, x)$ is a quadratic form which is the object under investigation in [ $\mathrm{M} 2,3$ ] and $[\mathrm{BPr}]$.)

The hermitian form $h_{T}$ is called $K$-rational if there exists an invertible element $a$ in $D_{T}$ and a $\lambda$-hermitian form $h$ on $D^{n}$ such that $h_{T}=a h$ and $h_{T}$ is called $K$-irrational in the opposite case. If $h_{S}$ is a nondegenerate hermitian form on $D_{S}^{n}$, we denote by $S_{0}$ the set of all $v \in S$ such that $h_{v}$ is isotropic (i.e. the $K_{v}$-algebraic group $\mathbf{S U}\left(h_{v}\right)$ corresponding to $h_{v}$ is $K_{v}$-isotropic, see 5.1). The form $h_{S}$ is called isotropic if $S_{0} \neq \emptyset$. We will denote by $h_{S_{0}}$ the hermitian form on $D_{S_{0}}^{n}$ given by all $h_{v}, v \in S_{0}$. The hermitian forms $h_{S}$ and $h_{S}^{\prime}$ are properly equivalent if there exists $g \in \mathrm{SL}_{n}\left(D_{S}\right)$ such that $h_{S}^{\prime}=h_{S}^{g}\left(\right.$ where $\mathrm{SL}_{n}\left(D_{S}\right)=\prod_{v \in S} \mathrm{SL}_{n}\left(D_{v}\right)$ acts in the usual way on the hermitian forms). We will say that $h_{S}$ and $h_{S}^{\prime}$
are almost $S$-integer equivalent if for any $\varepsilon>0$ and any $\mathcal{O}$-order $\Lambda$ in $D$ there exists $g \in \mathrm{SL}_{n}(\Lambda)$ such that

$$
\left\|h_{S}^{g}-h_{S}^{\prime}\right\|<\varepsilon
$$

where $\left\|\|\right.$ is a norm on the space of all hermitian forms on $D_{S}^{n}$ comparable with the topology on this space induced by the topology on $K_{S}$.

It is easy to prove that the almost $S$-integer equivalence between two hermitian forms implies their proper equivalence (see 5.6). The next theorem shows in particular that the converse is true for all $K$-irrational isotropic hermitian forms of dimension $n \geq 3$.

Theorem 5. With the above notations, let $h_{S}$ be a nondegenerate isotropic hermitian form on $D_{S}^{n}$. Assume that (a) rn $\geq 3$ if $\tau_{S}$ is of first kind and $\tau_{S} \neq i d$, (b) rn $\geq 2$ if $\tau_{S}$ is of second kind, and (c) $n \geq 3$ if $\tau_{S}=i d$. Then the following conditions are equivalent:
(i) $h_{S_{0}}$ is $K$-irrational hermitian form;
(ii) $h_{S}$ is properly equivalent to an hermitian form $h_{S}^{\prime}$ if and only if $h_{S}$ is almost $S$-integer equivalent to $h_{S}^{\prime}$.

Theorem 5 and the next corollary will be proved in section 5. (Also in $\S 5$ we give examples which show that the restrictions in the formulation of the theorem can not be weekened.)

Corollary 1. Let $n \geq 2, h_{S}$ be as in Theorem 1 and $h_{S_{0}}$ be $K$-irrational. Then for any $\varepsilon>0$, any $\mathcal{O}$-order $\Lambda$ in $D$ and any $x_{1}, x_{2}, \ldots, x_{n-1} \in D_{S}^{n}$ there exist $z_{1}, \ldots, z_{n-1} \in \Lambda^{n}$ such that

$$
\begin{equation*}
\left|h_{S}\left(x_{i}, x_{j}\right)-h_{S}\left(z_{i}, z_{j}\right)\right|<\varepsilon, \tag{1}
\end{equation*}
$$

for all $i, j$ (here $\left|\mid\right.$ stands for a norm on $D_{S}$ comparable with the topology on $\left.D_{S}\right)$. In particular, the closure of $\left\{h_{S}(z, z) \mid z \in \Lambda^{n}\right\}$ in $D_{S}$ coincides with $\left\{h_{S}(x, x) \mid x \in D_{S}^{n}\right\}$.

In the case of quadratic forms (i.e. $\tau_{S}=\mathrm{id}$ and $\lambda_{S}=1$ ), the above corollary was proved by Borel and Prasad (see [BPr], [B2]) for primitive vectors $z_{1}, \ldots z_{n-1} \in \mathcal{O}^{n}$ assuming that $S=S_{0}$. (For 3-dimensional real quadratic forms the result had been earlier proved by Dani and Margulis [DM1].) The observation in the case of quadratic forms that $S_{0} \subset S$ may be taken arbitrary nonempty belongs to Margulis.

In contrast to the case of quadratic forms, it is generally possible that all nondegenerate hermitian forms on $D_{S}^{n}$ are $K$-irrational. This occurs when $D$ is not abelian and $D$ does not admit $L / K$-involutions.

Recall that $D$ admits a nontrivial involution of first kind if and only if $D$ is a quaternion division algebra. The conditions for the existence on $D$ of an $L / K$-involution of second kind is given by relations between the local invariants of $D$ (cf. [Sch., 10.2.4]). Our theorem implies immediately the following criterion.

Corollary 2. Let $D$ be non-abelian and $L=K$ (resp. $L$ is a quadratic extension of $K$ ). Then the following assertions are equivalent :
(i) $D$ does not admit an involution of first (resp. an $L / K$-involution of second) kind;
(ii) If $r n \geq 3$ (resp. $r n \geq 2$ ) and two nondegenerate isotropic hermitian forms on $D_{S}^{n}$ are properly equivalent, then they are almost S-integer equivalent.

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## §1. Notation and terminology

1.1 As usual, $\mathbb{C}, \mathbb{R}, \mathbb{Q}, \mathbb{Q}_{p}, \mathbb{Z}, \mathbf{N}$ denote the complex, real, rational, rational p-adic, integer and natural numbers. Furthermore, $\mathbb{R}^{+}$will be the set of all strictly positive real numbers. $K$ will denote a number field, i.e.a finite extension of $\mathbb{Q}$. All valuations of $K$ under consideration will be normalized valuations (see [CaF, ch.1]). In particular, if $v$ is a valuation of $K$ and $K_{v}$ is the completion of $K$ with respect to $v$ then $\left|\left.\right|_{v}\right.$ and $\theta_{v}$ denote the normalized norm and the normalized Haar measure on $K_{v}$, respectively. The field $K_{v}$ contains a unique completion $\mathbb{Q}_{p_{v}}$ of Q, where $p_{v}$ is a prime number if $v$ is nonarchimedean and $p_{v}=\infty$ and $Q_{\infty}=\mathbb{R}$ if $v$ is archimedean.

If $v$ is a nonarchimedean valuation we denote by $\mathcal{O}_{v}$ the ring of integers of $K_{v}$. We will denote by $\mathcal{V}$ the adele ring of $K$ (i.e. $\mathcal{V}$ is the restricted topological product of all completions $K_{v}$ of $K$ relative to $\mathcal{O}_{v}$, $\left.v \notin S_{\infty}[\mathrm{CaF}, \operatorname{ch} .2]\right)$.
1.2 Throughout the paper, we fix a finite set $S$ of valuations of $K$ containing the set $S_{\infty}$. We denote by $\mathcal{O}$ the ring of $S$-integers of $K$ and by $\mathcal{O}^{*}$ the group of units in $\mathcal{O}$.

Given a subset $T$ of $S$, we denote by $\theta_{T}=\prod_{v \in T} \theta_{v}$ the Haar measure on $K_{T}$ and by $\left\|\|_{T}=\sup _{v \in T}| |_{v}\right.$, the norm on $K_{T}$.

Let $r: T \rightarrow \mathbb{R}^{+}$and $a=\left(a_{v}\right)$ be a point in $K_{T}$. By an interval with radius $r$ and center $a$ in $K_{T}$, we mean the set

$$
I(r, a)=\left\{\left(x_{v}\right) \in K_{T}| | x_{v}-\left.a_{v}\right|_{v} \leq r(v) \text { for all } v \in T\right\}
$$

If $a=0$, we write $I(r)$ instead of $I(r, 0)$.
1.3 By a $K_{T^{-}}$algebraic variety $\mathbf{M}$, we mean a (formal) direct product $\prod_{v \in T} \mathbf{M}_{v}$ of $K_{v}$-algebraic varieties $\mathbf{M}_{v}$. A map $f: \mathbf{M} \rightarrow \mathbf{M}^{\prime}$, where $\mathbf{M}$ and $\mathbf{M}^{\prime}$ are $K_{T}$-algebraic varieties, is called $K_{S}$-rational (resp. $K_{S}$-regular) if $f$ is a product of $K_{v}$-rational (resp. $K_{v}$-regular) maps $f_{v}: \mathbf{M}_{v} \rightarrow$ $\mathbf{M}_{v}^{\prime}, v \in S$. If in the above definition of $\mathbf{M}$ all $\mathbf{M}_{v}, v \in S$, are $K_{v^{-}}$ algebraic groups we say that $\mathbf{M}$ is a $K_{S^{-}}$algebraic group. Analogously, we define the notions of linear $K_{S}$-space, $K_{S}$-rational representation of a $K_{S^{-}}$algebraic group etc. The set $\prod_{v \in S} \mathbf{M}_{v}\left(K_{v}\right)$ will be denoted by $\mathbf{M}\left(K_{S}\right)$ or simply $M$ and called set of $K_{S^{-}}$-rational points of $\mathbf{M}$ (respectively, group of $K_{S}$-rational points, in the case of a $K_{S}$-algebraic group M). Naturally if $\mathbf{V}$ is a $K$-algebraic variety, we associate to it a $K_{S}$-algebraic variety, also denoted by $\mathbf{V}$, such that $\mathbf{V}_{v}=\mathbf{V}$ for all $v \in S$.
1.4 By Zariski topology on a $K_{S}$-algebraic variety $\mathbf{M}$, we mean the product of the Zariski topologies on $\mathbf{M}_{v}, v \in S$. $\mathbf{M}$ is called connected if every $\mathbf{M}_{v}, v \in S$, is connected for the Zariski topology.

On $\mathbf{M}\left(K_{S}\right)$ we have two topologies : one induced by the Zariski topology on $\mathbf{M}$ and another which is a product of the locally compact Hausdorff topologies on $\mathbf{M}\left(K_{v}\right), v \in S$. In order to distinguish the two topologies, all topological notions connected with the first one will be used with the prefix "Zariski". (We will say : Zariski closure, Zariski closed, etc.) Given a subset $X$ in $\mathbf{M}\left(K_{S}\right)$, we will denote by $\bar{X}$ the Zariski closure of $X$ in $\mathbf{M}\left(K_{S}\right)$. By a $K_{S}$-algebraic subvariety of $\mathbf{M}\left(K_{S}\right)$, we mean a Zariski closed subset of $\mathbf{M}\left(K_{S}\right)$.
1.5 Let $\mathbf{G}$ be an algebraic group defined over $K$. Every $\mathbf{G}\left(K_{v}\right)$ is naturally embedded in $G$. By a 1 -parameter unipotent $K_{v}$-subgroup $U_{v}=\left\{u_{v}(t)\right\}$ of $G$, we mean a nontrivial $K_{v}$-rational homomorphism $u_{v}: K_{v} \rightarrow \mathbf{G}\left(K_{v}\right)$. Given a subgroup $M$ of $G$, we will denote by $M_{u}$ the subgroup of $M$ generated by all 1-parameter unipotent $K_{v}$-subgroups for all $v \in S$.

Let $T \subset S$ and for each $v \in T$ let $U_{v}=\left\{u_{v}\left(t_{v}\right) \mid t_{v} \in K_{v}\right\}$ be a 1-parameter unipotent $K_{v}$-subgroup. Then the homomorphism $u_{T}$ :
$K_{T} \rightarrow G,\left(t_{v}\right)_{v \in T} \rightarrow\left(u_{v}\left(t_{v}\right)\right)_{v \in T}$, defines a 1-parameter unipotent $K_{T^{-}}$ subgroup $U_{T}=\prod_{v \in T} U_{v}$ of $G$.

Finally, given a 1-parameter unipotent $K_{v}$-subgroup $U_{v}=\left\{u_{v}(t)\right\}$ of $G$ and a 1-dimensional over $\mathbb{Q}_{p_{v}}$ linear subspace $l$ in $K_{v}$ the restriction of $u_{v}$ on $l$ is called 1-parameter unipotent $\Phi_{p_{v}}$-subgroup of $G$.

## §2. Proof of Theorem 3 and of the strong approximation theorem.

Theorem 3 will be proved in 2.1-2.4. We preserve the notations from its announcement.
2.1 Denote by $\mu$ the $M$-invariant Borel probability measure on $M x$. Note that $\mu$ is $M_{u}$-ergodic and $M^{\prime} x=M x$ for each open subgroup $M^{\prime} \subset M$ containing $M_{u}$. Replacing, if necessary, $M$ by its subgroup of finite index, we will assume that $M$ satisfies the following conditions:
$\left(^{*}\right)$ Every open subgroup of $M$ is Zariski dense in $\bar{M}$ and $M=$ $\prod_{v \in S} M_{v}$ where $M_{v} \subset \mathbf{G}\left(K_{v}\right)$.

We need the following version of the Borel density theorem.
Lemma $M_{x}$ is Zariski dense in $M$.
Proof. Since $M x$ is closed $M x$ is homeomorphic to $M / M_{x}$ and, therefore, $\mu$ can be considered as a $M$-invariant $M_{u}$-ergodic Borel probability measure on $M / M_{x}$. Let $\mu_{1}$ be the image of the measure $\mu$ on $M / M \cap \bar{M}_{x}$ under the natural map $M / M_{x} \rightarrow M / M \cap \bar{M}_{x}$. Clearly the measure $\mu_{1}$ is $M_{u}$-invariant and $M_{u}$-ergodic. On the other hand, $M / M \cap \bar{M}_{x}$ can be regarded as a Borel subset of a $K_{S^{-}}$algebraic variety on which $M_{u}$ acts rationally. This in view of [MTo1, 3.1] implies that $\mu_{1}$ is concentrated in a point. The lemma is proved.
2.2 We will denote by $R(\bar{M})$ the solvable radical of $\bar{M}$ (i.e. $R(\bar{M})$ is the maximal connected in the Zariski topology solvable normal subgroup of $\bar{M})$. Clearly $\bar{M}=\prod_{v \in S} \bar{M}_{v}$ and $R(\bar{M})=\prod_{v \in S} R\left(\bar{M}_{v}\right)$ where $R\left(\bar{M}_{v}\right)$ is the solvable radical of $\bar{M}_{v}$.

Lemma Assume that $R(\bar{M})=R(\bar{M})_{u}$. Then $M$ is open in the Hausdorff topology of $\bar{M}$.

Proof. It is enough to prove the lemma when $S=\{v\}$. Let $\mathcal{G}$ be the Lie algebra of $M_{v}$. It is well known that the commutator of $\mathcal{G}$ coincides with the Lie algebra of the commutator $\mathcal{D}^{1}\left(\bar{M}_{v}\right)$ of $\bar{M}_{v}$. (cf. [C, ch.2, theorem 13]). Therefore the commutator $\mathcal{D}^{1}\left(M_{v}\right)$ of $M_{v}$ is open
in $\mathcal{D}^{1}\left(\bar{M}_{v}\right)$. This allows to reduce the proof of the lemma to the case when $M_{v}$ is abelian. In this case $\bar{M}_{v}$ is isomorphic to a vector space $K_{v}^{n}$. Now the lemma follows from the assumption $\left(^{*}\right)$ and from the fact that if $v$ is archimedean (respectively, nonarchimedean) then the connected component of $M_{v}$ is Zariski closed (respectively, $\{1\}$ is the only discrete subgroup of $M_{v}$ ).
2.3 Lemma Let $R(\bar{M})=R(\bar{M})_{u}$ and also let $\bar{M} x$ be closed and admit $\bar{M}$-invariant Borel probability measure $\bar{\mu}$. Then $M^{\prime}=\{a \in \bar{M} \mid$ $a M x=M x\}$ is a subgroup of finite index in $\bar{M}$.

Proof. Since $M_{u}$ is a normal subgroup of $\bar{M}$ and $M_{u} x$ is dense in $M x$ we get that for any $a \in \bar{M}$ either $a M x=M x$ or $a M x \cap M x=\emptyset$. In view of $2.2 \bar{\mu}(a M x)=\bar{\mu}(M x)>0$ for all $a \in \bar{M}$. Therefore there exists a finite subset $\left\{a_{1}, a_{2}, \ldots, a_{r}\right\} \subset M$ such that $\bar{M} x=a_{1} M x \cup a_{2} M x \cup$ $\ldots \cup a_{r} M x$ and the multiplication from the left by an element from $\bar{M}$ permutes the subsets $a_{i} M x, i=1,2, \ldots, r$. This implies the lemma.
2.4 Proof of Theorem 3. Replacing $\mu$ by $g^{-1} \mu$ and $M$ by $g^{-1} M g$ we may (and will) assume that $g=e$. Let $\mathbf{P}_{1}$ be the largest normal subgroup of class $\mathcal{F}$ in $\mathbf{P}$ and $P_{1}=\mathbf{P}_{1}\left(K_{S}\right)$. In view of Lemma $2.1 M$ is contained in $P$ and, therefore, $M_{u} \subset P_{1}$. Since $P_{1} \cap \Gamma$ is a lattice in $P_{1}$ the orbit $P_{1} x$ is closed and contains $M x$. In particular, $P_{1}$ contains an open subgroup of $M$. In view of the assumption $\left(^{*}\right)$, Lemma 2.1 and the definition of $\mathbf{P}$ we get that $\mathbf{P}=\mathbf{P}_{1}$ and $P=\bar{M}$. It follows from 2.2 and 2.3 that $M$ is open in $P$ and $P^{\prime}=\{a \in P \mid a \mu=\mu\}$ has finite index in $P$. The uniqueness of $\mathbf{P}$ follows from the fact that the Lie algebras of $\mathbf{P}\left(K_{v}\right)$ and $M_{v}, v \in S$, coincide. The theorem is proved.
2.5 Proof of Theorem 4. Let us first consider the case when $\mathbf{G}$ is semisimple. In view of Weil's restriction of scalars functor [W2, ch.1] and the result of Borel and Tits [BTi2, 6.21(ii)] we may (and will) assume that $\mathbf{G}$ is absolutely almost simple. Also, it is easy to see that it is enough to prove the theorem for $T=\left\{v_{o}\right\}$. Assuming all this, denote by $S_{o}$ the (finite) set of all nonarchimedean valuations $v$ of $K$ such that $\mathbf{G}$ is $K_{v}$-anisotropic. Let $S_{1}$ be any finite set of valuations such that $S_{1} \cap\left(S_{o} \cup\left\{v_{o}\right\}\right)=\emptyset$ and $S_{1} \cup\left\{v_{o}\right\}$ contains all archimedean valuations of $K$. For each $v \in S_{o}$, fix an open subgroup $R_{v}$ in $\mathbf{G}\left(K_{v}\right)$. Define an open subgroup $A$ in $\mathbf{G}(\mathcal{V})$ as follows:

$$
A=\prod_{v \in S_{1} \cup\left\{v_{o}\right\}} \mathbf{G}\left(K_{v}\right) \times \prod_{v \in S_{o}} R_{v} \times \prod_{v \notin S_{o} \cup S_{1} \cup\left\{v_{o}\right\}} \mathbf{G}\left(\mathcal{O}_{v}\right) .
$$

The group $\Gamma=A \cap \mathbf{G}(K)$ is an $S$-arithmetic subgroup, where $S=$ $S_{o} \cup S_{1} \cup\left\{v_{o}\right\}$. Since $\mathbf{G}$ is simply connected, $\mathbf{G}\left(K_{v_{o}}\right)=\mathbf{G}\left(K_{v_{o}}\right)_{u}$ (cf. $\left.[\mathrm{Pl}]\right)$.

Let $G^{\prime}$ be the closure of $\mathbf{G}\left(K_{v_{o}}\right) \Gamma$ in $\mathbf{G}\left(K_{S}\right)$. In view of Theorem $3, G^{\prime}$ is a subgroup of finite index in $\mathbf{G}\left(K_{S}\right)$. Note that $\mathbf{G}\left(K_{v}\right)$ does not contain a subgroup of finite index if $v$ is archimedean (Cartan) and, also, if $v$ is nonarchimedean and $\mathbf{G}$ is $K_{v}$-isotropic (cf. [Ti1],[Pl]). Therefore $G^{\prime} \supset \mathbf{G}\left(K_{S_{1} \cup\left\{v_{o}\right\}}\right)$. Since the subgroups $R_{v}, v \in S_{o}$, can be chosen arbitrary small and the finite set $S_{1}$ arbitrary large, we obtain that the closure of $\mathbf{G}\left(K_{v_{o}}\right) \mathbf{G}(K)$ in $\mathbf{G}(\mathcal{V})$ contains all adeles $x=\left(x_{v}\right)$ with $x_{v}=1$ for each $v \in S_{o}$. If $S_{o} \neq \emptyset$ then $\mathbf{G}$ is a group of type $A_{n}$ and, for every $v \in S_{o}, \mathbf{G}\left(K_{v}\right)$ is the group of elements with reduced norm 1 in a central division algebra over $K_{v}$. By a result of Kneser [K2] the diagonal embedding of $\mathbf{G}(K)$ in $\prod_{v \in S_{o}} \mathbf{G}\left(K_{v}\right)$ is dense. Therefore $\mathbf{G}\left(K_{v_{o}}\right) \mathbf{G}(K)$ is dense in $\mathbf{G}(\mathcal{V})$.

For a non-semisimple $K$-algebraic group $\mathbf{G}$, we apply the following standard argument. The group $\mathbf{G}$ is a semidirect product over $K$ of its unipotent radical $\mathbf{U}$ and its semisimple $K$-algebraic subgroup $\mathbf{L}$. Therefore $\mathbf{G}(\mathcal{V}) \mathbf{L}(\mathcal{V}) \mathbf{U}(\mathcal{V})$ and $\mathbf{G}(K)=\mathbf{L}(K) \mathbf{U}(K)$. Now, the theorem follows from the validity of the strong approximation for both $\mathbf{L}$ and $\mathbf{U}$. This completes the proof.

## §3. Ratner's uniform distribution theorem

3.1 Similarly to [Rat2,3], we will deduce Theorem 1 from its stronger version for 1-parameter unipotent $K_{T}$-subgroups, the so-called uniform distribution theorem.

Theorem. Let $G, \Gamma$ and $x=g \Gamma$ be as in Theorem 1. Furthermore, let $T$ be a nonempty subset of $S$ and $U=\left\{u(t) \mid t \in K_{T}\right\}$ be a 1parameter unipotent $K_{T}$-subgroup of $G$ and let $\left\{I\left(r_{i}\right)\right\}$ be an increasing sequence of intervals in $K_{T}$ such that $\cup_{i} I\left(r_{i}\right)=K_{T}$. Then there exists a subgroup $\mathbf{P} \subset \mathbf{G}$ of class $\mathcal{F}$ and a subgroup of finite index $P^{\prime}$ in $P=\mathbf{P}\left(K_{S}\right)$ such that the closure of the orbit $U x$ coincides with $g P^{\prime} g^{-1} x$ and

$$
\lim _{i \rightarrow \infty} \frac{1}{\theta_{T}\left(I\left(r_{i}\right)\right)} \int_{I\left(r_{i}\right)} f(u(t) x) d \theta_{T}(t)=\int_{g P^{\prime} g^{-1} x} f(y) d \mu(y)
$$

for any bounded continuous function $f$ on $G / \Gamma$, where $\mu$ is the Haar measure on $g P^{\prime} g^{-1} x$.

For 1-parameter real and p-adic subgroups $U=\{u(t)\}$ the uniform distribution theorem was proved in the general context of direct products of real and $p$-adic Lie groups in [Rat3, Theorem 3]. (In fact in [Rat3] and, earlier, in [DM2] for real Lie groups, a stronger version of this theorem
is proved : the point $x$ is replaced by a sequence of points converging to a generic point $x$ (cf [Rat3, Theorem 4] and [DM2, Theorem 2]. In the present paper we do not treat this more technical case.) Our proof of Theorem 3.1 uses methods, with some modifications, of the proof in [DM2].
3.2 Deduction of Theorem 1 from Theorem 3.1. It is enough to prove Theorem 1 for $x=\Gamma$. Let $H=H_{u}$ as in the formulation of Theorem 1. Denote by $M$ the subset of all unipotent elements in $H$. It is clear that $M$ is a $K_{S}$-algebraic subvariety of $G$. Since each element in $M$ is contained in a maximal unipotent subgroup of $H$ and any two maximal unipotent subgroups of $H$ are conjugated [B3,15.9], $M$ is Zariski connected. On the other hand, every element $a \in M$ is contained in a 1-parameter unipotent $K_{T}$-subgroup $U(a)$ of $H$. In view of Theorem 3.1, there exists a $K$-algebraic subgroup $\mathbf{P}(a) \subset \mathbf{G}$ of class $\mathcal{F}$ and a subgroup of finite index $P(a)^{\prime}$ in the group of $K_{S}$-rational points of $\mathbf{P}(a)$ such that the closure of $U(a) x$ coincides with $P(a)^{\prime} x$. Let $A$ be an open subset in $M$ homeomorphic to a neighbourhood of 0 in some linear $K_{S}$-space. Then since $\left\{P(a)^{\prime} \mid a \in M\right\}$ is a countable set, there exists $a_{o} \in M$ such that $P\left(a_{o}\right)^{\prime} \cap A$ has positive Lebesque measure. As $M$ is Zariski connected, we get that $P\left(a_{o}\right)^{\prime} \cap A$ is Zariski dense in $M$. Therefore, $H \subset P\left(a_{o}\right)^{\prime}$, which implies Theorem 1.
3.3 The following result is important for the proof of Theorem 3.1.

Theorem. With $G, \Gamma$ and $U$ as in Theorem 3.1, let $\varepsilon>0$ and $\mathcal{K} \subset G / \Gamma$ be a compact. Then there exists a compact $\mathcal{K}_{1}$ in $G / \Gamma$ such that for any $x \in \mathcal{K}_{2}$ and any interval $I(r)$ in $K_{T}$,

$$
\frac{1}{\theta_{T}(I(r))} \theta_{T}\left\{t \in I(r) \mid u(t) x \in \mathcal{K}_{1}\right\} \geq 1-\varepsilon
$$

In the case when $T=\{v\}$, the above theorem was announced in [MT1,11.4] with indications about the proof. (The details will appear elsewhere.) The general case follows from this one by a simple application of the Fubini theorem. In the real case the theorem is proved in [DM2,6.1] using earlier results [M1], [D1,3,4] and the arithmeticity theorem [M4,ch.9].
3.4 Singular and generic points. Given a subgroup $U$ of $G$ and a proper subgroup $\mathbf{P} \subset \mathbf{G}$ of class $\mathcal{F}$, we put $X(P, U)=\{g \in G \mid U g \subset$ $g P\}$. It is clear that $X(P, U)$ is a $K_{S}$-algebraic subvariety of $G$. We denote $\mathcal{S}(U)=\cup_{\mathbf{P} \in \mathcal{F}, \mathbf{P} \neq \mathbf{G}} X(P, U) \Gamma / \Gamma$ and $\mathcal{G}(U)=G / \Gamma-\mathcal{S}(U)$. As in [DM2], the points from $\mathcal{S}(U)$ (resp. $\mathcal{G}(U)$ ) are called singular (resp. generic) points with respect to $U$.

Theorem 3.1 will be derived from the following

Proposition. Let $G, \Gamma, U, \varepsilon$ and $\mathcal{K}$ be as in Theorem 3.3. Also let $\mathbf{P}$ be a proper subgroup of $\mathbf{G}$ of class $\mathcal{F}$ and $C=\prod_{v \in S} C_{v}$ be a compact subset of $X(P, U)$, where $C_{v} \subset \mathbf{G}\left(K_{v}\right)$ for each $v \in S$. Then there exists a compact $D=\prod_{v \in S} D_{v}$ in $X(P, U)$ such that $D_{v} \supset C_{v}$ for all $v \in S$, $D_{v}=C_{v}$ for $v \notin T$ and the following holds: For any neighbourhood $\Phi_{o}$ of $D$ in $G$ there exists a neighbourhood $\Phi$ of $C$ in $G$, such that for any $x \in \mathcal{K}-\Phi_{o} \Gamma / \Gamma$ and any interval $I(r)$ in $K_{T}$,

$$
\begin{equation*}
\frac{1}{\theta_{T}(I(r))} \theta_{T}\{t \in I(r) \mid u(t) x \in \Phi \Gamma / \Gamma\}<\varepsilon \tag{2}
\end{equation*}
$$

3.5 Deduction of the theorem for uniform distribution from Proposition 3.4. Let $\mathbf{P}$ be the smallest subgroup of class $\mathcal{F}$ in $\mathbf{G}$ such that $U g \subset g P$ (if $x$ is singular then $\mathbf{P} \neq \mathbf{G}$ ). Put $U_{1}=g^{-1} U g$ and $\Delta=\Gamma \cap P$. Then $\Delta$ is an $S$ - arithmetic lattice in $P, y=\Delta$ is a generic point in $P / \Delta$ with respect to $U_{1} \subset P$, and $P \Gamma / \Gamma$ is closed in $G / \Gamma$ and homeomorphic to $P / \Delta$. This reduces the proof of the theorem to the case when $x$ is generic, which we will assume from now on.

For any interval $I(r)$ define a probability measure $\mu_{r}$ on $G / \Gamma$ by the formula

$$
\int_{G / \Gamma} f(y) d \mu_{r}(y)=\frac{1}{\theta_{T}(I(r))} \int_{I(r)} f(u(t) x) d \theta_{T}(t)
$$

where $f$ is a bounded continuous function on $G / \Gamma$. We denote by $\widetilde{G / \Gamma}$ the one-point compactification of $G / \Gamma$ if $G / \Gamma$ is not compact, and $G / \Gamma$ itself if $G / \Gamma$ is compact. It is well-known that, given a compact metrizable topological space $Y$, the space $\mathcal{P}(Y)$ of all Borel probability measures on $Y$ is compact with respect to the weak * topology. Let $\left\{I\left(r_{i}\right)\right\}$ be a sequence of intervals as in the formulation of Theorem 3.1. Put $\mu_{i}=\mu_{r_{i}}$ for every i. The sequence of measures $\left\{\mu_{i}\right\}$ is naturally embedded in $\mathcal{P}(\widetilde{G / \Gamma})$. Let $\lambda$ be a limit point of $\left\{\mu_{i}\right\}$ in $\mathcal{P}(\widetilde{G / \Gamma})$. It follows from 3.3 that $\lambda$ is concentrated on $G / \Gamma$. Using the fact that $f$ is bounded and performing the linear substitution $s=t_{o}+t$ in the integrals

$$
\int_{I\left(r_{i}\right)} f\left(u\left(t_{o}+t\right) x\right) d \theta_{T}(t)
$$

a simple argument shows that

$$
\lambda\left(u\left(t_{o}\right) f\right)=\left(u\left(-t_{o}\right) \lambda\right)(f)=\lambda(f)
$$

for any bounded continuous function $f$ on $G / \Gamma$. This means that $\lambda$ is $U$-invariant. On the other hand, since $X(P, U)$ is second countable
topological space and $\mathcal{F}$ is countable it follows from Proposition 3.4 applied for $\mathcal{K}=\{x\}$ that $\lambda(\mathcal{S}(G / \Gamma))=0$. In view of Theorem 2 , every $U$ invariant $U$-ergodic measure on $G / \Gamma$ which is not supported by an orbit of an open subgroup of $G$ is supported by $\mathcal{S}(G / \Gamma)$. This, in view of the decomposition of $\lambda$ into a continuous sum of its $U$-ergodic components implies that $\lambda$ is $U$-ergodic and coincides with the Haar measure on $G^{\prime} x$ where $G^{\prime}$ is an open subgroup of $G$. So, assuming the validity of Proposition 3.4, the proof of Theorem 3.1 is complete.

## §4. Polynomial-like behaviour of the unipotent orbits and proof of Proposition 3.4

4.0 Let us make the following simple observation : Given a field $F$ and a 1-parameter unipotent subgroup $U=\{u(t)\}$ in $G L_{n}(F)$, the $\operatorname{map} t \rightarrow u(t) x$ is polynomial of degree less that or equal to $n$ for each $x \in F^{n}$. This fact is in the root of the phenomenon that the dynamics of the actions of subgroups on homogeneous spaces are much easier to be understood when the subgroups are generated by unipotent elements than, say, when they are generated by split semisimple elements. (In more geometrical terms, the first type of actions corresponds to a horosperical flow and the second one to a geodesic flow on a Riemannian manifold with constant negative curvature.) The above observation will be used in the proof of the key Proposition 4.2. For the proof of 4.2, we need a property of polynomial maps given by the following lemma. (Note that Lemma 4.1 plays also a crucial role in the proof of Theorem 3.3.)
4.1 Lemma. Let $v$ be either a real or a nonarchimedean valuation of $K$ and $\epsilon$ and $\alpha$ be positive reals. Also, let $f=\left(f_{1}, \ldots, f_{s}\right)$ be a polynomial map $K_{v} \rightarrow K_{v}^{s}$ of degree not greater than $n \in \mathbb{N}$ (i.e. degf $f_{i} \leq n$ for all i). Put $\delta=\frac{\varepsilon^{n}}{(n+1)^{n+1}}$ and denote by $\left\|\|_{v}\right.$ the norm sup on $K_{v}^{s}$. Then for any interval $I \subset K_{v}$ which contains a number $t_{o}$ with $\left\|f\left(t_{o}\right)\right\|_{v} \geq \alpha$,

$$
\theta_{v}\left\{t \in I \mid \| f\left(t_{o} \|_{v} \leq \delta \alpha\right\} \leq \varepsilon \theta_{v}\left\{t \in I \mid \| f\left(t_{o} \| \leq \alpha\right\}\right.\right.
$$

Proof. We will omit the case when $K_{v}=\mathbb{R}$ which is considered in detail in [DM2,4.1]. Without loss of generality we may assume that $\alpha$ is a value of the norm $\left|\left.\right|_{v}\right.$. It is easy to see that, given two intervals in $K_{v}$ with nonempty intersection, one of them contains the other. So, for any $x \in I$ there exists a unique interval $J(x)$ in $I$ which contains $x$ and is maximal with the property : $\|f(y)\|_{v} \leq \alpha$ for all $y \in J(x)$. Therefore, it
is enough to consider the case $I=J(x)$. Clearly, $I$ contains an element $y_{o}$ with $\left\|f\left(y_{o}\right)\right\|_{v}=\alpha$. Since there exists an $i$ such that $\left|f_{i}\left(y_{o}\right)\right|_{v}=\alpha$, the proof is reduced to the case when $s=1$. Furthermore, replacing $f$ by $a f$, where $a \in K_{v}$ and $|a|_{v}=\alpha^{-1}$, and doing a linear substitution $t \rightarrow t+t_{o}$, we reduce the proof of the lemma to the following case : $I=\left\{\left.t \in K_{v}| | t\right|_{v \leq 1}\right\}, t_{o}=0,|f(0)|_{v}=\alpha=1, \varepsilon \leq 1$, and $|f(t)|_{v} \leq 1$ for all $t \in I$. Assume that the statement of the lemma is false, that is,

$$
\theta_{v}(A)>\varepsilon,
$$

where $A=\left\{\left.t \in I| | f(t)\right|_{v} \leq \delta\right\}$. Denote $I_{1}=\left\{\left.t \in K_{v}| | t\right|_{v} \leq \frac{\varepsilon}{n+1}\right\}$. Since the norm $\left|\left.\right|_{v}\right.$ and the measure $\theta_{v}$ on $K_{v}$ are normalized, we obtain

$$
\theta_{v}\left(I_{1}\right) \leq \frac{\varepsilon}{n+1}
$$

cf. [CaF, ch.1]. Therefore, there exist points $\alpha_{1}, \ldots, \alpha_{n+1}$ in $A$ such that $\left\{\alpha_{i}+I_{1}\right\} \cap\left\{\alpha_{j}+I_{1}\right\}=\emptyset$ for all $i \neq j$. (Note that $I_{1}$ is an ideal in the ring $I$.) In particular,

$$
\begin{equation*}
\left|\alpha_{i}-\alpha_{j}\right|_{v}>\frac{\varepsilon}{n+1} \tag{3}
\end{equation*}
$$

for all $i \neq j$. Let us write the Lagrange interpolation formula for $f$ at the points $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{n+1}$ :

$$
f(t)=\sum_{i=1}^{n+1} f\left(\alpha_{i}\right) \frac{\left(x-\alpha_{1}\right) \ldots\left(x-\alpha_{i-1}\right)\left(x-\alpha_{i+1}\right) \ldots\left(x-\alpha_{n+1}\right)}{\left(\alpha_{i}-\alpha_{1}\right) \ldots\left(\alpha_{i}-\alpha_{i-1}\right)\left(\alpha_{i}-\alpha_{i+1}\right) \ldots\left(\alpha_{i}-\alpha_{n+1}\right)}
$$

The substitution $t=0$, the inequality (5), and the ultrametric inequality in $K_{v}$ imply that

$$
|f(0)|_{v}<\frac{\delta(n+1)^{n}}{\varepsilon^{n}}<1
$$

Contradiction. The lemma is proved.
4.2 Proposition. Let $M$ be a Zariski closed subset in $K_{v}^{m}$. Then for any compact subset $A$ of $M$ and any $\varepsilon>0$ there exists a compact $B$ in $M$ containing $A$ such that the following holds: given a compact neighbourhood $W_{o}$ of $B$ in $K_{v}^{m}$, there exists a neighbourhood $W$ of $A$ in $K_{v}^{m}$ such that for any 1-parameter unipotent subgroup $\{u(t)\}$ in $G L_{m}\left(K_{v}\right)$, any $a \in K_{v}^{m}-W_{o}$ and any interval $I$ in $K_{v}$ containing 0 , we have

$$
\begin{equation*}
\theta_{v}\{t \in I \mid u(t) a \in W\} \leq \varepsilon\left\{t \in I \mid u(t) a \in W_{o}\right\} \tag{4}
\end{equation*}
$$

Proof. Note that the case $K_{v}=\mathbb{C}$ can be easily reduced to the real one by embedding $M$ in $\mathbb{R}^{2 m}$ via Weil's restriction of scalars and by using Lemma 4.4 below. The real case itself is considered in [DM2, 4.2]. So, we will assume that $v$ is nonarchimedean. First note that if the proposition is true for a compact subset of $M$ containing $A$, then it is also true for $A$. Therefore, it is enough to consider the case $A=\left\{x \in M \mid\|x\| \leq R_{o}\right\}$ where $R_{o}>0$ is a constant and $\left\|\|\right.$ is the norm sup on $K_{v}^{m}$. Let $f_{1}, f_{2}, \ldots, f_{r} \in K_{v}\left[x_{1}, . ., x_{m}\right]$ be such that $M=\left\{x \in K_{v}^{m} \mid f_{i}(x)=\right.$ 0 for all $i=1,2, \ldots, r\}$. Let $n \in \mathbf{N}$ be such that the degree of each polynomial $f_{i}$ is $\leq n$. Put $\delta=\frac{\varepsilon^{m n}}{(m n+1)^{m n+1}}$. Let $R$ be a real number such that $R \geq R_{o} \delta^{-1}$ and $R$ be a value of the norm $\|\|$. Denote,

$$
B=\{x \in M \mid\|x\| \leq R\}
$$

Let $W_{o}$ be a neighbourhood of $B, a \notin W_{o}$ and $u(t)$ be a 1-parameter unipotent subgroup of $G L_{m}\left(K_{v}\right)$. Denote $g(t)=u(t) a, t \in K_{v}$. Then $g(t)=\left(g_{1}(t), \ldots, g_{m}(t)\right)$ where $g_{i}(t)$ are polynomials of degree $\leq m$. Also denote $F(t)=\left(F_{1}(t), \ldots, F_{r}(t)\right)$ where $F_{i}(t)=f_{i}\left(g_{1}(t), \ldots, g_{m}(t)\right)$ for all $i$. It is easy to see (for example, by an argument from the contrary and using the compactness of $B$ ) that there exists an $\alpha>0$ which is a value of the norm $\left|\left.\right|_{v}\right.$ and such that

$$
W_{1}=\left\{x \in K_{v}^{m} \mid\|x\| \leq R \text { and }\left|f_{i}(x)\right|_{v} \leq \alpha \text { for all } i\right\}
$$

is a neighbourhood of $B$ contained in $W_{o}$. We will prove that

$$
W=\left\{x \in K_{v}^{m} \mid\|x\| \leq R_{o} \text { and }\left|f_{i}(x)\right|_{v} \leq \alpha \delta \text { for all } i\right\}
$$

is the required neighbourhood of $A$. Let $I \subset K_{v}$ be an interval containing 0 . Put $J=\left\{t \in I \mid u(t) a \in W_{1}\right\}$. For each $t \in J$, we denote by $J(t)$ the maximal (closed) subinterval of $I$ such that $u(J(t)) a \subset W_{1}$. Since $a \notin W_{1}$ and $0 \in I$, for every $t \in J$ there exists $t^{\prime} \in J(t)$ such that either $\left\|u\left(t^{\prime}\right) a\right\|=R$ or $\left\|F\left(t^{\prime}\right)\right\|=\alpha$. Using 4.1 we get

$$
\varepsilon \theta_{v}(J(t)) \geq \theta_{v}\left\{y \in J(t) \mid\|u(y) a\| \leq R_{o} \text { and }\|F(y)\| \leq \alpha \delta\right\}
$$

Since the intervals $J(t)$ form a partition of $J$, the above formula implies (4).
4.3 Lemma. Assume that Proposition 3.4 is valid if $T$ is a singleton. Then it is valid for any $T$.

Proof. The proof is by induction on the cardinality of $T$. Let $T=T_{1} \cup T_{2}$, where both $T_{1}$ and $T_{2}$ are nonempty and $U=U_{1} \times U_{2}$, where $U_{i}=\left\{u_{i}(t) \mid t \in K_{T_{i}}\right\}, i=1,2$, are 1-parameter $K_{T_{i}}$-subgroups
of $G$. Assume that Proposition 3.4 is true for both $U_{1}$ and $U_{2}$. Let $\mathcal{K}, C$ and $\varepsilon$ be as in 3.4. We have to prove that a compact $D$, as in the formulation of Proposition 3.4, exists. Because of 3.3, there exists a compact subset $\mathcal{K}_{1}$ in $G / \Gamma$ such that for any $x \in \mathcal{K}$ and any interval $I_{1}^{\prime}$ in $K_{T_{1}}$ containing 0 ,

$$
\begin{equation*}
\frac{1}{\theta_{1}\left(I_{1}^{\prime}\right)} \theta_{1}\left\{t \in I_{1}^{\prime} \mid u_{1}(t) x \in \mathcal{K}_{1}\right\} \geq 1-\frac{\varepsilon}{3} \tag{5}
\end{equation*}
$$

where $\theta_{1}=\theta_{T_{1}}$ (see 1.2). Note that $X(P, U)=X\left(P, U_{1}\right) \cap X\left(P, U_{2}\right)$. Applying Proposition 3.4 for $U_{2}$, we get a compact $D^{\prime}=\prod_{v \in S} D_{v}^{\prime}$ in $X(P, U)$ with $D_{v}^{\prime}=C_{v}$ for all $v \notin T_{2}$ and such that if $\Psi$ is a neighbourhood of $D^{\prime}$ then there exists a neighbourhood $\Phi$ of $C$ such that for any $y \in \mathcal{K}_{1}-\Psi \Gamma / \Gamma$ and any interval $I_{2}$ in $K_{T_{2}}$ containing 0 , we have

$$
\begin{equation*}
\frac{1}{\theta_{2}\left(I_{2}\right)} \theta_{2}\left\{t \in I_{2} \mid u_{2}(t) y \notin \Phi \Gamma / \Gamma\right\} \geq 1-\frac{\varepsilon}{3} \tag{6}
\end{equation*}
$$

where $\theta_{2}=\theta_{T_{2}}$. Applying again Proposition 3.4 (this time for $U_{1}$ ), we get a compact $D=\prod_{v \in S} D_{v}$ in $X(P, U)$ with $D_{v}=D_{v}^{\prime}$ for all $v \notin$ $T_{1}$ and such that if $\Phi_{o}$ is a neighbourhood of $D$ then there exists a neighbourhood $\Psi$ of $D^{\prime}$ with

$$
\begin{equation*}
\frac{1}{\theta_{1}\left(I_{1}\right)} \theta_{1}\left\{t \in I_{1} \mid u_{1}(t) x \notin \Psi \Gamma / \Gamma\right\} \geq 1-\frac{\varepsilon}{3} \tag{7}
\end{equation*}
$$

for any interval $I_{1}$ in $K_{T_{1}}$ containing 0 and any $x \in \mathcal{K}-\Phi_{o} \Gamma / \Gamma$.
Let $I=I_{1} \times I_{2}$ be an interval in $K_{T}$ containing $0, \Phi_{o}$ a neighbourhood of $D, \Psi$ a neighbourhood of $D^{\prime}$ as given by (7) and $\Phi$ a neighbourhood of $C$ as given by (6). Now it follows from (5)-(7) and the Fubini theorem that

$$
\begin{gathered}
\frac{1}{\theta_{T}(I)} \theta_{T}\{t \in I \mid u(t) x \notin \Phi \Gamma / \Gamma\} \geq \\
\frac{1}{\theta_{T}(I)} \theta_{T}\left\{t=\left(t_{1}, t_{2}\right) \in I \mid u_{1}\left(t_{1}\right) x \in \mathcal{K}_{1}-\Psi \Gamma / \Gamma \text { and } u_{1}\left(t_{1}\right) u_{2}\left(t_{2}\right) x \notin \Phi \Gamma / \Gamma\right\} \geq \\
\left(1-\frac{\varepsilon}{3}\right)^{3} \geq 1-\varepsilon
\end{gathered}
$$

for any $x \in \mathcal{K}-\Phi_{o} \Gamma / \Gamma$. The lemma is proved.
4.4 Let $T=\{v\}, K_{v}=\mathbb{C}$ and $U=\{u(t)\}$ a 1-parameter unipotent $\mathbb{C}$-subgroup of $G$. In this case the proof of Proposition 3.4 can be reduced to the case of actions of 1-parameter unipotent real subgroup of $G$ as follows. For every 1-dimensional real subspace $l \subset \mathbb{C}$ we denote by $u_{l}$ the restriction of $u$ on $l$. Put $U_{l}=\left\{u_{l}(t)\right\}$. In order to prove that (2) holds it is enough to show that for all $l \subset \mathbb{C}$ and all intervals $I(r)=[-r, r]$ in $\mathbb{R}$ we have,

$$
\begin{equation*}
\frac{1}{\theta_{o}(I(r))} \theta_{o}\left\{t \in I(r) \mid u_{l}(t) x \in \Phi \Gamma / \Gamma\right\}<\varepsilon \tag{8}
\end{equation*}
$$

where $\theta_{o}$ is the Lebesques measure on $\mathbb{R}$. The fact that the fulfilment of (8) for all $l$ implies (2) follows from the elementary lemma below.

Lemma.Let $I=\{t \in \mathbb{C}| | t \mid \leq 1\}, \varepsilon>0$ and $A$ be a measurable subset of $I$ such that for any $x \in I$ we have,

$$
\varepsilon \theta_{o}\{a \in \mathbb{R} \mid a x \in I\} \geq \theta_{o}\{a \in \mathbb{R} \mid a x \in I \cap A\} .
$$

Then $\theta_{v}(A)<\varepsilon \pi$.
4.5 Up to the end of section 4 , we preserve the notations from 3.4 and suppose that $T=\{v\}$. Denote by $\mathcal{U}=\{u(t)\}$ a 1-parameter unipotent $\mathbb{R}$-subgroup if $v$ is archimedean and put $\mathcal{U}=U$ if $v$ is nonarchimedean. We let $F=\mathbb{R}$ in the former and $F=K_{v}$ in the latter case. Note that $X(P, U)=X(P, \mathcal{U})$.

Let us fix a $K$-rational representation $\varrho: \mathbf{G} \rightarrow \mathbf{G L}(\mathbf{V})$ such that the normalizer $\mathcal{N}_{\mathbf{G}}(\mathbf{P})$ of $\mathbf{P}$ in $\mathbf{G}$ coincides with the stabilizer in $\mathbf{G}$ of a 1-dimensional subspace of $\mathbf{V}$ spanned by a vector $m \in \mathbf{V}(K)$. (The existence of such a representation follows from the Chevalley theorem [B3,5.1].) Let $\chi$ be the $K$-rational character of $\mathcal{N}_{\mathbf{G}}(\mathbf{P})$ given by $g m=\chi(g) m, g \in \mathcal{N}_{\mathbf{G}}(\mathbf{P})$. We denote $\mathbf{N}=\{g \in \mathbf{G} \mid g m=m\}$, $N=\mathbf{N}\left(K_{S}\right), \Gamma_{N}=\Gamma \cap N$ and $\Gamma_{P}=\Gamma \cap \mathcal{N}_{\mathbf{G}}(\mathbf{P})$. Let $\eta: \mathbf{G} \rightarrow \mathbf{G} m$, $g \rightarrow g m$. (We will denote in the same way the map $G \rightarrow V, g \rightarrow g v$, where $V=\mathbf{V}\left(K_{S}\right)$.) Note that $\mathbf{G} m$ is open in its Zariski closure, $\mathbf{G} m$ is isomorphic to $\mathbf{G} / \mathbf{N}$ and $\eta$ is a quotient map [B3, 6.7]. Let $\mathbf{X}=\{g \in \mathbf{G} \mid U g \subset g \mathbf{P}\}$. Clealy, $\mathbf{X}$ is a $K_{S}$-algebraic variety and $\mathbf{X}\left(K_{S}\right)=X(P, U)$. We will denote $X_{v}=\mathbf{X}_{v}\left(K_{v}\right)$. Since $\mathbf{X} \mathcal{N}_{\mathbf{G}}(\mathbf{P})=\mathbf{X}$, $\mathbf{N} \subset \mathcal{N}_{\mathbf{G}}(\mathbf{P})$ and $\eta$ is a Zariski open map, we get that $\eta(\mathbf{X})$ is Zariski closed in $\mathbf{G} m$. This implies that

$$
\begin{equation*}
\eta^{-1} \overline{(\eta(X(P, U))}=X(P, U) \tag{9}
\end{equation*}
$$

(In the above formula the Zariski closure is taken in $V$.)
4.6 In the following remarks we use some standard facts from algebraic number theory (cf.[CaF],[W1]).
(a) Since $\chi$ is a $K$-character, $\chi(\Gamma) \cap \mathcal{O}^{*}$ has finite index in $\chi(\Gamma)$. In view of the facts that $\mathcal{O}$ is integrally closed and $\chi(\Gamma) \subset K$, we get that $\chi(\Gamma) \subset \mathcal{O}^{*}$. Also since $\mathbf{V}(\mathcal{O})$ is discrete in $V$ and $\eta$ is $K$-rational, $\Gamma m$ is discrete in $V$.
(b) Denote by $K_{S}^{1}$ the set of all $x=\left\{x_{w}\right\}_{w \in S} \in K_{S}$ such that

$$
\prod_{w \in S}\left|x_{w}\right|_{w}=1
$$

The group $\mathcal{O}^{*}$ is diagonally embedded in $K_{S}^{1}$. For each $w \in S$ let $\lambda_{w}: K_{w}^{*} \rightarrow \mathbb{R}$ be the $\operatorname{map} x \rightarrow \log |x|_{w}$. Put $R_{w}=\operatorname{Im}\left(\lambda_{w}\right)$. So, $R_{w}$ coincides with $\mathbb{R}$ if $v$ is archimedean and $R_{w}$ is a cyclic subgroup of $\mathbb{R}$ if $v$ is nonarchimedean. Let $\lambda: K_{S}^{1} \rightarrow \prod_{w \in S} R_{w}$ be the direct sum of all $\lambda_{w}$ and let $R=\operatorname{Im}(\lambda)$. Then $R$ is a locally compact abelian group, $\lambda\left(\mathcal{O}^{*}\right)$ is a lattice in $R$ and $\operatorname{Ker}(\lambda) \cap \mathcal{O}^{*}$ is the group of roots of unity in $K$ $[\mathrm{CaF}, \mathrm{ch} .2,18.1]$. Therefore there exists $\delta>0$ such that if $\xi \in \mathcal{O}^{*}$ and $|1-\xi|_{w}<\delta$ for all $w \neq v$ then $\xi$ is a root of unity in $K$.
(c) Let $A=\prod_{w \in S} A_{w}$ be a subset of $G$. We will say that $A$ is $S(v)$ small if for every $w \neq v$ in $S$ the following holds : if $c \in K_{w}^{*}$ is such that

$$
c\left(A_{w} m\right) \cap A_{w} m \neq \emptyset
$$

then $|c-1|_{w}<\delta$. In particular, if $c \in \mathcal{O}^{*}$ then, in view of (b), $c$ is a root of unity.
(d) Clearly, every element $g \in G$ is contained in a $S(v)$-suffitiently small neighbourhood.

The consideration of $S(v)$-sufficiently small subsets for $\mathbb{Q}$-algebraic varieties was suggested by G.A.Margulis.
4.7 Proposition. Let $\phi: G / \Gamma_{N} \rightarrow G / \Gamma \times V, \phi\left(g \Gamma_{N}\right)=(g \Gamma, g m)$. Then $\phi$ is a proper map.

Proof. Let $\left\{g_{i} \Gamma_{N}\right\}$ be a sequence in $G / \Gamma_{N}$ such that $\phi\left(g_{i} \Gamma_{N}\right)$ converges to $(c \Gamma, q) \in G / \Gamma \times V$. Fix $c_{i} \in G$ and $\gamma_{i} \in \Gamma$ such that $g_{i}=c_{i} \gamma_{i}$ for all $i$ and $\lim _{i} c_{i}=c$. As $\left\{\gamma_{i} m\right\} \subset V$ is discrete (cf. 4.6(a)), there exists $i_{o}$ such that $\gamma_{i} m=\gamma_{i_{o}} m$ for all $i \geq i_{o}$. So, $\gamma_{i} \Gamma_{N}=\gamma_{i_{o}} \Gamma_{N}$ for all $i \geq i_{o}$. Therefore $\left\{g_{i} \Gamma_{N}\right\}$ is bounded in $G / \Gamma_{N}$ which proves that $\phi$ is a proper map.
4.8 The above proposition implies the following

Corollary. Let $D_{o}$ and $L$ be compact subsets in $G$. Then there exists a compact $D$ in $G$ such that $D_{o} \subset D \subset D_{o} N$ and

Furthermore, if $\Omega$ is a neighbourhood of $D$ then $\Omega$ contains a neighbourhood $\Psi$ of $D_{o}$ such that

$$
\Psi N \cap L \Gamma \subset \Psi \Gamma_{N} .
$$

According to $4.7, \phi^{-1}\left(L \Gamma / \Gamma, D_{o} m\right)$ is a compact subset of $D_{o} N / \Gamma_{N}$. Now, the existence of $D$ satisfying (10) follows by a simple continuity argument. The second part can be proved in a similar way.
4.9 Let $A$ be a subset of $G$. Following [DM2], a point $x \in A$ will be called a point of $(P, \Gamma)$-self-intersection in $A$ if there exists $\gamma \in \Gamma-\Gamma_{P}$ such that $x \gamma \in A$. The next proposition corresponds to Corollary 3.5 in [DM2].

Proposition. Let $D_{o}$ and $L$ be compact subsets of $G$ and $Y$ be the (closed) subset of all points of $(P, \Gamma)$-self-intersections in $D_{o}$. Assume that $D_{o} N \cap L \Gamma \subset D_{o} \Gamma_{N}$. Then for every relatively compact neighbourhood $\Psi$ of $Y$ there exists an open neighbourhood $\Omega$ of $D_{o}$ such that

$$
\left(\Omega-\Psi \Gamma_{P}\right) \cap L \Gamma
$$

does not contain points of $(P, \Gamma)$-self-intersections.
Proof. Assume the contrary, that is, there exists a sequence of neighbourhoods $\left\{\Omega_{i}\right\}$ of $D_{o}$ such that $\Omega_{i} \supset \Omega_{i+1}, \bigcap_{i}=D_{o}$ and there exist $g_{i}, g_{i}^{\prime} \in\left(\Omega_{i}-\Psi \Gamma_{P}\right) \cap L \Gamma$ with $g_{i}=c_{i} \gamma_{i}, g_{i}^{\prime}=c_{i} \gamma_{i}^{\prime}$, where $c_{i} \in L$, $\gamma_{i}$ and $\gamma_{i}^{\prime} \in \Gamma$, and $\gamma_{i}^{-1} \gamma_{i}^{\prime} \notin \Gamma_{P}$. Passing to subsequences, we may (and will) assume that each of the sequences $\left\{c_{i}\right\},\left\{g_{i} m\right\}$ and $\left\{g_{i}^{\prime} m\right\}$ converges. Since $\Gamma m$ is discrete, there exists $i_{o}$ such that $\gamma_{i} m=\gamma_{i_{o}} m$ and $\gamma_{i}^{\prime} m=\gamma_{i_{o}}^{\prime} m$ for all $i \geq i_{o}$. Put $c=\lim _{i} c_{i}$. Then $c \gamma_{i_{o}} m$ and $c \gamma_{i_{o}}^{\prime} m \in$ $D_{o} m$. Therefore $c \gamma_{i_{o}}, c \gamma_{i_{o}}^{\prime} \in D_{o} N \cap L \Gamma \subset D_{o} \Gamma_{N}$. As $\gamma_{i_{o}}^{-1} \gamma_{i_{o}}^{\prime} \notin \Gamma_{P}$, we get that $c \gamma_{i_{o}} \in Y \Gamma_{P}$. The latter contradicts the fact that $c \gamma_{i_{o}} \notin \Psi \Gamma_{P}$. The proposition is proved.
4.10 Proof of Proposition 3.4 Let $\mathcal{K} \subset G / \Gamma, \varepsilon>0$ and $C$ be as in the formulation of Proposition 3.4. According to 4.3 and 4.4, it is enough to prove (2) for $\mathcal{U}$ as defined in 4.5. Also, in view 4.6(d) we can (as we will) suppose that $C$ is $S(v)$-small subset of $X(P, U)$.

The proposition will be proved by induction on $\operatorname{dim} \mathbf{P}$. (The proof is trivial for $\operatorname{dim} \mathbf{P}=0$.)

Using 4.2, we can find a compact $B=\prod_{w \in S} B_{w}$ in $\overline{X(P, U) m} \subset V$ such that $B_{w}=C_{w} m$ for all $w \neq v, B_{v} \supset C_{v} m$, and for any neighbourhood $\mathcal{A}_{o}$ of $B_{v}$ in $\mathbf{V}\left(K_{v}\right)$ there exists a neighbourhood $\mathcal{A}$ of $C_{v} m$ in $\mathbf{V}\left(K_{v}\right)$ such that

$$
\begin{equation*}
\theta\{t \in I(r) \mid u(t) a \in \mathcal{A}\} \leq \frac{\varepsilon}{2 k} \theta\left\{t \in I(r) \mid u(t) a \in \mathcal{A}_{o}\right\} \tag{11}
\end{equation*}
$$

for all $a \in X_{v}-\mathcal{A}_{o}$, where $k$ is the order of the group of roots of unity in $K$. (Here and later on $\theta$ is the Haar measure on $F$ and $I(r)$ is an interval in $F$ with radius $r$ centered at 0 .)

Applying 3.3 , we fix a compact $L \subset G$ such that $\mathcal{K} \subset L \Gamma / \Gamma$ and

$$
\begin{equation*}
\frac{1}{\theta(I(r))} \theta\{t \in I(r) \mid u(t) x \in L \Gamma / \Gamma\} \geq 1-\frac{\varepsilon}{4} \tag{12}
\end{equation*}
$$

for all $x \in \mathcal{K}$.
In view of $4.7,4.8$ and (9), there exists a compact $D_{o} \subset X(P, U)$ which satisfies

$$
\phi^{-1}(L \Gamma \times B) \subset D_{o} \Gamma_{N} / \Gamma_{N}
$$

and

$$
D_{o} N \cap L \Gamma \subset D_{o} \Gamma_{N}
$$

Denote by $Y$ the subset of all points of $(P, \Gamma)$-self- intersection in $D_{o}$. If $y \in Y$ there exists $\gamma \in \Gamma-\Gamma_{P}$ such that $y \gamma \in Y$. This implies that $\mathcal{U} y \subset Q y$ where $Q=\mathbf{Q}\left(K_{S}\right)$ and $\mathbf{Q}$ is a group from the class $\mathcal{F}$ contained in $\mathbf{P} \cap \gamma \mathbf{P} \gamma^{-1}$, in particular, $\operatorname{dim} \mathbf{Q}<\operatorname{dim} \mathbf{P}$. Since $D_{o}$ is compact and $\Gamma$ is discrete in $G$, there are finitely many proper algebraic subgroups $\mathbf{P}_{1}, \ldots, \mathbf{P}_{s}$ of $\mathbf{P}$ such that $\mathbf{P}_{i} \in \mathcal{F}$ for all $i$ and $\bigcup_{i \geq 1} X\left(P_{i}, \mathcal{U}\right) \supset$ $Y$. Denote $C_{i}=X\left(P_{i}, \mathcal{U}\right) \cap Y, i=1,2, \ldots, s$. By the induction hypothesis there exists for every $i$ a compact $D_{i} \subset X\left(P_{i}, \mathcal{U}\right)$ so that if $\Phi_{o}$ is any neighbourhood of $\bigcup_{i \geq 1} D_{i}$ then we have an open neighbourhood $\Psi$ of $\bigcup_{i \geq 1} C_{i}$ such that

$$
\begin{equation*}
\frac{1}{\theta(I(r))} \theta\{t \in I(r) \mid u(t) x \in \Psi \Gamma / \Gamma\} \leq \frac{\varepsilon}{4}, \tag{13}
\end{equation*}
$$

for all $x \in \mathcal{K}-\Phi_{o} \Gamma / \Gamma$ and any interval $I(r)$.

Put $D=\bigcup_{i \geq 0} D_{i}$. Now, let us fix a neighbourhood $\Phi_{o}$ of $D$. We will prove that there exists a neighbourhood $\Phi$ of $C$ which satisfies (2). Let $\Psi$ be a neighbourhood of $\bigcup_{i \geq 1} C_{i}$ which satisfies (13) for the last choice of $\Phi_{o}$. Using 4.9, 4.8 and the definition of $D_{o}$, one can find a neighbourhood $\Omega$ of $D_{o}$ such that $\Omega \subset \Phi_{o}$, the set $\left(\Omega-\Psi \Gamma_{P}\right) \cap L \Gamma$ is without points of ( $P, \Gamma$ )-self-intersections, and

$$
\Omega N \cap L \Gamma \subset \Omega \Gamma_{N}
$$

This, together with (9) and the fact that $B$ is $S(v)$-small, implies that there exists a compact $S(v)$-small neighbourhood $W_{o}$ of $B$ in $V$ such that

$$
\begin{equation*}
\phi^{-1}\left(W_{o} \times L \Gamma / \Gamma\right) \subset \Omega \Gamma_{N} / \Gamma_{N} \tag{14}
\end{equation*}
$$

Using the property (11) of $B$, as well as the fact that $\mathcal{U}$ acts trivially on $\mathbf{V}\left(K_{w}\right)$ for all $w \neq v$, we fix a neighbourhood $W$ of $C$ in $V$ such that if $a \in X(P, \mathcal{U}) m-W_{o}$ and $I$ is a maximal subinterval of $I(r)$ with $u(I) a \subset W_{o}$ then

$$
\begin{equation*}
\theta\{t \in I \mid u(t) a \in W\} \leq \frac{\varepsilon}{2 k} \theta(I) \tag{15}
\end{equation*}
$$

We will prove that $\Phi=\eta^{-1}(W)$ is the neighbourhood of $C$ which we need. Let $x=g \Gamma \in \mathcal{K}-\Phi_{o} \Gamma / \Gamma$. Denote

$$
J^{(1)}=\{t \in I(r) \mid u(t) x \notin L \Gamma / \Gamma \text { or } u(t) x \in \Psi \Gamma / \Gamma\}
$$

and

$$
J^{(2)}=\{t \in I(r) \mid u(t) x \in(\Phi \Gamma / \Gamma \cap L \Gamma / \Gamma)-\Psi \Gamma / \Gamma\}
$$

It is clear that

$$
\begin{equation*}
J^{(1)} \cup J^{(2)} \supset\{t \in I(r) \mid u(t) x \in \Phi \Gamma / \Gamma\} \tag{16}
\end{equation*}
$$

In view of (12) and (13)

$$
\begin{equation*}
\theta\left(J^{(1)}\right) \leq \frac{\varepsilon}{2} \theta(I(r)) \tag{17}
\end{equation*}
$$

Assume that there exists $\gamma \in \Gamma$ with $g \gamma m \in W_{o}$. Then since $x \in$ $L \Gamma / \Gamma$, it follows from (14) that $g \gamma \in \Omega \Gamma_{N}$ which, in view of the inclusion $\Omega \subset \Phi_{o}$, implies that $x \in \Phi_{o} \Gamma / \Gamma$. Contradiction. Therefore,

$$
\begin{equation*}
g \gamma m \notin W_{o}, \tag{18}
\end{equation*}
$$

for all $\gamma \in \Gamma$.
Next, for every $q \in g \Gamma m$, we define a subset $J_{q}$ in $I(r)$ in the following way: (i) if $v$ is nonarchimedean then $t \in J_{q}$ iff $t$ is contained by a subinterval $I$ of $I(r)$ such that $u(I) g \gamma m \subset W_{o}$ and $u\left(t^{\prime}\right) x \in L \Gamma / \Gamma-\Psi \Gamma / \Gamma$ for some $t^{\prime} \in I$, and (ii) if $v$ is archimedean then $t \in J_{q}$ iff $t$ is contained by a subinterval $[\alpha, \beta]$ in $I(r)$ such that $u([\alpha, \beta]) g \gamma m \subset W_{o}$ and $u(\beta) x \in L \Gamma / \Gamma-\Psi \Gamma / \Gamma$. Let $t \in J_{q} \cap J_{q^{\prime}}$ where $q=g \gamma m$ and $q^{\prime}=g \gamma^{\prime} m$. Denote by $J_{q}(t)\left(\right.$ resp. $\left.J_{q^{\prime}}(t)\right)$ the maximal interval in $J_{q}$ (resp. $J_{q^{\prime}}$ ) containing $t$. It follows from the definition of $J_{q}$ and $J_{q^{\prime}}$ (and from the fact that in the nonarchimedean case if two intervals have nonempty intersection then one of them contains the other) that there exists $t_{o} \in J_{q}(t) \cap J_{q^{\prime}}(t)$ such that $u\left(t_{o}\right) x \in L \Gamma / \Gamma-\Psi \Gamma / \Gamma$. It follows from (14) that $u\left(t_{o}\right) g \gamma$ and $u\left(t_{o}\right) g \gamma^{\prime}$ belong to $\Omega \Gamma_{N}$. Since $\left(\Omega \Gamma_{N}-\Psi \Gamma_{P}\right) \cap L \Gamma$ is a set without $(P, \Gamma)$-self-intersections, we obtain that $\gamma^{\prime}=\gamma \delta$, where $\delta \in \Gamma_{P}$. Therefore $u\left(t_{o}\right) g \gamma^{\prime} m=\chi(\delta) u\left(t_{o}\right) g \gamma m$.

Since $\chi(\delta) \in \mathcal{O}^{*}, u\left(t_{o}\right) g \gamma m$ and $u\left(t_{o}\right) g \gamma^{\prime} m$ belong to $W_{o}$, and $W_{o}$ is $S(v)$-small set, using 4.6(c) we obtain that

$$
\begin{equation*}
q^{\prime}=\xi q \tag{19}
\end{equation*}
$$

where $\xi$ is a root of unity in $K$.
Applying (15) and (18), we get

$$
\frac{\varepsilon}{2 k} \theta\left(J_{q}(t)\right) \geq \theta\left(J_{q}(t) \cap J^{(2)}\right)
$$

Since for any $t$ and $t^{\prime}$ in $J_{q}$ either $J_{q}(t)=J_{q}\left(t^{\prime}\right)$ or $J_{q}(t) \cap J_{q}\left(t^{\prime}\right)=\emptyset$, we obtain

$$
\begin{equation*}
\frac{\varepsilon}{2 k} \theta\left(J_{q}\right) \geq \theta\left(J_{q} \cap J^{(2)}\right) \tag{20}
\end{equation*}
$$

Now since $\bigcup_{q} J_{q} \supset J^{(2)}$, it follows from (19) and (20) that

$$
\frac{\varepsilon}{2} \theta(I(r)) \geq \frac{\varepsilon}{2} \theta\left(\cup_{q} J_{q}\right) \geq \frac{\varepsilon}{2 k} \sum_{q} \theta\left(J_{q}\right) \geq \sum_{q} \theta\left(J_{q} \cap J^{(2)}\right) \geq \theta\left(J^{(2)}\right)
$$

This, in view of (16) and (17), completes the proof.

## §5. Applications to the Hermitian forms.

5.0 In this section we prove Theorem 5 and its corollaries, after first developping the necessary algebraic background for the irrational hermitian forms. We conclude the section by giving some examples and making some remarks about possible generalizations and strenghtenings of Theorem 5 .
5.1 Let $R$ be a ring with center $Z$ and an involution $\sigma$ (i.e. $\sigma$ is an antiautomorphism of $R$ of order two). Also let $\lambda \in Z$ be such that ${ }^{\sigma} \lambda \lambda=1$. A $\lambda$-hermitian form (relatively to the involution $\sigma$ ) on the right free $R$-module $R^{n}$ is a sesquilinear map $h: R^{n} \times R^{n} \rightarrow R$ such that

$$
\begin{equation*}
h(x, y)=\lambda^{\sigma} h(y, x) \tag{21}
\end{equation*}
$$

for all $x, y \in R^{n}$. The hermitian form $h$ is nondegenerate if the map $\hat{h}: R^{n} \rightarrow \operatorname{Hom}_{R}\left(R^{n}, R\right),(\hat{h} x)(y)=h(x, y)$, is an isomorphism of abelian groups [Sch,7.1.3]. Further on, by an hermitian form we will mean always a nondegenerate hermitian form.
5.2 Unless something else is specified, in the subsections $5.2-5.5$ we will denote by $D$ a central division algebra of degree $r$ over an arbitrary infinite field $L$ of characteristic $\neq 2$. As in the Introduction, $K$ is a subfield of $L$ such that either $L=K$ or $L$ is a quadratic extension of $K$. Let $K_{1}$ be any field extension of $K, D_{1}=D \bigotimes_{K} K_{1}$ and $L_{1}=L \bigotimes_{K} K_{1}$. (In the applications, $D$ will be a division algebra over a number field $L$, $K_{1}$ will stand for the completion $K_{v}$ of $K$ at a valuation $v$ of $K, L_{1}=L_{v}$ and $D_{1}=D_{v}$.) We will assume that $D_{1}$ admits an involution $\tau$ which is a $L_{1} / K_{1}$-involution, that is, $K_{1}=\left\{\left.x \in L_{1}\right|^{\tau} x=x\right\}$. (Recall that $\tau$ is an involution of first (respectively, second) kind if $L_{1}=K_{1}$ (respectively, $\left.L_{1} \neq K_{1}\right)$.)

There are two possibilities : either $L_{1}$ is a field or $L_{1}=K_{1} \bigoplus K_{1}$. Let first $L_{1}$ be a field. Then $D_{1}$ coincides with a matrix algebra $\mathrm{M}_{s}(\Delta)$ with entries from a central division algebra $\Delta$ over $L_{1}$. It is known [K1, Theorem, p.37] that $\Delta$ admits an involution ${ }^{-}: \Delta \rightarrow \Delta$ which is of the same kind as $\tau$. We can define a standard involution $\sigma$ on $D_{1}$ as follows: ${ }^{\sigma}\left(a_{i j}\right)=\left(\bar{a}_{j i}\right)$ for all $\left(a_{i j}\right) \in \mathrm{M}_{s}(\Delta)$.

Now let $L_{1}=K_{1} \bigoplus K_{1}$. Then the restriction of $\tau$ on $L_{1}$ transposes the direct summands of $L_{1}$. This implies that $D_{1}=\mathrm{M}_{s}(\Delta) \bigoplus \mathrm{M}_{s}\left(\Delta^{o}\right)$ where $\Delta$ is a division algebra with center $K_{1}$ and $\Delta^{o}$ is the division algebra opposite to $\Delta$ (i.e. $\Delta^{o}$ coincides with $\Delta$ as abelian group and has the multiplication $x . y=y x)$. In this case ${ }^{\tau}(x, y)=(y, x)$ for all $(x, y) \in D_{1}$.

The relation between the different involutions of the same kind on $D_{1}$ is given by the following Proposition. Its proof is similar to [Sch, 8.7.4].

Proposition. Let $\sigma$ and $\tau$ be $L_{1} / K_{1}$-involutions on $D_{1}$. Then $\tau=\sigma \circ \operatorname{Int}(d)$, where $d$ is an invertible element of $D_{1}$ such that ${ }^{\sigma} d= \pm d$ when $\sigma$ and $\tau$ are involutions of first kind, and ${ }^{\sigma} d=d$ when $\sigma$ and $\tau$ are involutions of second kind.

Proof. By the Scolem-Noether theorem, $\tau=\sigma \circ \operatorname{Int}(b)$ where $b \in$ $D_{1}$. A simple direct argument shows that

$$
i d=\tau^{2}=\operatorname{Int}\left(\left({ }^{\sigma} b^{-1}\right) b\right)
$$

So, ${ }^{\sigma} b=l b$ where $l \in L_{1}$. This implies that ${ }^{\sigma} l l=1$. Hence if $\sigma$ is of first kind then $l= \pm 1$ and we can choose $d=b$. Otherwise there exists $c \in L_{1}$ such that $l=\frac{c}{\sigma_{c}}, c \in L$. (The existence of $c$ follows from the Hilbert 90 theorem if $L_{1}$ is a field. If $L_{1}=K_{1} \bigoplus K_{1}$ and $l=\left(s, s^{-1}\right)$ then we can choose $c=(s, 1)$ since $\sigma$ acts on $L$ by interchanging the two coordinates.) Put $d=c b$. It is easy to check that ${ }^{\sigma} d=d$. The proposition is proved.
5.3 Let $h$ be a $\lambda$-hermitian form on $D_{1}^{n}, n \geq 1$, with respect to $\tau$. Since ${ }^{\tau} \lambda \lambda=1$, we get that $\lambda= \pm 1$ if $\tau$ is of first kind. Let $\tau$ be of second kind and let $c$ be an invertible element in $L_{1}$. It follows from (21) that $c h$ is a $\lambda^{\prime}$-hermitian form with respect to $\tau$ where $\lambda^{\prime}=\lambda\left(\frac{c}{\tau_{c}}\right)$. By Hilbert 90 theorem (and its simple analog when $L_{1}=K_{1} \bigoplus K_{1}$ ) c can be chosen in such a way that $\lambda=\frac{{ }^{\tau} c}{c}$. In this case, $c h$ is a 1 -hermitian form.

Let $d \in D_{1}$ be such that ${ }^{\tau} d=\varepsilon d$ where $\varepsilon= \pm 1$. An easy computation shows that $\tau^{\prime}=\tau \circ \operatorname{Int}(d)$ is an involution on $D_{1}$ and $h^{\prime}=d h$ is an $\varepsilon \lambda$-hermitian form with respect to $\tau^{\prime}[\mathrm{Sch}, 7.6 .7]$. The converse of this assertion is given by the following proposition.

Proposition. Let $h$ be a $\lambda$-hermitian form on $D_{1}^{n}$ with respect to a $L_{1} / K_{1}$-involution $\tau$ and $h^{\prime}$ be a $\lambda^{\prime}$-hermitian form on $D_{1}^{n}$ with respect to a $L_{1} / K_{1}$-involution $\tau^{\prime}$. Assume that $h^{\prime}=a h$ where $a \in D_{1}$. Then there exist $\alpha \in L_{1}$ and $d \in D_{1}$ such that $a=\alpha d, \tau^{\prime}=\tau \circ \operatorname{Int}(d), \lambda^{\prime}=\lambda\left(\frac{\alpha}{\tau_{\alpha}}\right)$ and ${ }^{\tau} d= \pm d$ if $\tau$ is of first kind and ${ }^{\tau} d=d$ if $\tau$ is of second kind.

Proof. Since $h$ and $h^{\prime}$ are nondegenerate hermitian forms $a$ is an invertible element in $D_{1}$. Let $x, y \in D_{1}^{n}$ and $c \in D_{1}$. Then

$$
\left(\tau^{\tau^{\prime}} c\right) h^{\prime}(x, y)=a\left({ }^{\tau} c\right) h(x, y)=a\left({ }^{\tau} c\right) a^{-1} h^{\prime}(x, y)
$$

Hence

$$
\begin{equation*}
\tau^{\prime} c=a\left({ }^{\tau} c\right) a^{-1} \tag{22}
\end{equation*}
$$

On the other hand, $\tau^{\prime}=\tau \circ \operatorname{Int}(d)$ with $d$ as in Proposition 5.2. Then $d h$ is a $\lambda$-hermitian form with respect to $\tau^{\prime}$. Applying (22) with $d h$ instead of $h$ we get that $a d^{-1}$ is in the center of $D_{1}$ i.e. $a=\alpha d$ where $\alpha \in L_{1}$. The fact that $\lambda^{\prime}=\lambda\left(\frac{\alpha}{\tau \alpha}\right)$ follows from (21). The proposition is proved.
5.4 Let $h$ be a $\lambda$-hermitian form on $D_{1}^{n}$ with respect to an involution $\tau$ such that $\lambda= \pm 1$ if $\tau$ is of first kind and $\lambda=1$ if $\tau$ is of second kind (see 5.3). Denote by Nrd : $\mathrm{M}_{n}\left(D_{1}\right) \rightarrow L_{1}$ the usual reduced norm on $\mathrm{M}_{n}\left(D_{1}\right)$ if $D_{1}$ is a simple algebra and the direct sum of the reduced norms on $\mathrm{M}_{n s}(\Delta)$ and $\mathrm{M}_{n s}\left(\Delta^{o}\right)$ if $D_{1}=\mathrm{M}_{s}(\Delta) \bigoplus \mathrm{M}_{s}\left(\Delta^{o}\right)$ (see 5.1). The special unitary group corresponding to $h$ is defined by $\mathrm{SU}(h)=\{g \in$ $\mathrm{M}_{n}\left(D_{1}\right) \mid \operatorname{Nrd}(g)=1$ and $h(x, y)=h(g x, g y)$ for all $\left.x, y \in D_{1}^{n}\right\}$. The group $\mathrm{SU}(h)$ coincides with the group of $K_{1}$-rational points of a $K_{1-}$ algebraic group $\mathbf{S U}(h)$. The hermitian form $h$ is isotropic if $\mathbf{S U}(h)$ is $K_{1}$-isotropic, equivalently, if $\mathrm{SU}(h)$ contains a diagonizable over $K_{1}$ infinite subgroup. It is known ( and follows easily from the classification results in [Ti2] and [K1]) that any classical algebraic group defined over an infinite field $K_{1}$ of characteristic $\neq 2$ can be realized as $\mathbf{S U}(h)$ for certain $h$. If $\tau$ is of first kind then $\mathbf{S U}(h)$ gives all $K_{1}$-algebraic groups of types $\mathrm{B}_{m}, \mathrm{C}_{m}$ and $\mathrm{D}_{m}$. If $\tau$ is of second kind then we get all $K_{1^{-}}$ algebraic groups of type $\mathrm{A}_{m}$. Note that if $L_{1}$ is a field then $h$ is isotropic if and only if $h$ represents nontrivially 0 . If $L_{1}=K_{1} \bigoplus K_{1}$ (i.e. $D_{1}=$ $\left.\mathrm{M}_{s}(\Delta) \bigoplus \mathrm{M}_{s}\left(\Delta^{o}\right)\right)$ then the description of $\tau$ in 5.2 implies that $\mathrm{SU}(h)$ coincides with the image of $\mathrm{SL}_{n s}(\Delta)$ in $\mathrm{M}_{n s}(\Delta) \bigoplus \mathrm{M}_{n s}\left(\Delta^{o}\right)$ under the embedding $g \rightarrow\left(g, g^{-1}\right)$. In particular, $\mathbf{S U}(h)$ is $K_{1}$-isotropic if and only if $n s>1$.

Let us summarize the last observations in the following lemma.

Lemma. Assume that either $n \neq 1$ or $D_{1} \neq \Delta \bigoplus \Delta^{o}$ where $\Delta$ is a division algebra. Then $h$ is isotropic if and only if $h$ represents nontrivially zero.

Remark. Let $K$ be a number field, $K_{1}=K_{v}$ be the completion of $K$ at a valuation $v$ of $K$ and $h$ be isotropic. Recall that if $\mathbf{G}$ is a simple $K_{v}$-isotropic $K_{v}$-algebraic group then the subgroup of $\mathbf{G}\left(K_{v}\right)$ generated by its unipotent elements has finite index [BTi1, 6.14]. Also, if $L_{v}$ is a field then the only central division algebras over $L_{v}$ with an involution of the first kind (respectively, an $L_{v} / K_{v}$-involution of second kind) are $L_{v}$ itself and the unique quaternion division algebra (respectively, $L_{v}$ itself) [Sch, 10.2.2]. Using these facts and the above description of $\mathbf{S U}(h)$, one can easily see that the unipotent elements in $\mathrm{SU}(h)$ generate a subgroup of finite index (and, therefore, Zariski dense subgroup) if and only if (a) $r n \geq 3$ if $\tau$ is of first kind and $\tau \neq \mathrm{id}$, (b) $r n \geq 2$ if $\tau$ is of second kind, and (c) $n \geq 3$ if $\tau=$ id. If the inequality in some of the cases (a)-(c) is not fulfilled then $\mathrm{SU}(h)$ is abelian, it consists of semisimple elements, and Theorem 5 (as well as Theorem 1) is not true (see 5.8).
5.5 Proposition. With the notations from 5.4, let $\Sigma$ be a subgroup of $S U(h) \cap S L_{n}(D)$ which is Zariski dense in $\mathbf{S U}(h)$. Assume that (a) $r n \geq 3$ if $\tau$ is of first kind and $\tau \neq i d$, (b) $r n \geq 2$ if $\tau$ is of second kind, and (c) $n \geq 3$ if $\tau=i d$. Then there exist, defined by $\Sigma$, an involution $\sigma$ on $D$ of the same kind as $\tau$ and an hermitian form $h_{\circ}$ on $D^{n}$ with respect to $\sigma$ such that $h=a h_{\circ}, a \in D_{1}$. In particular, $S U(h) \cap$ $S L_{n}(D)=S U\left(h_{\circ}\right)$.

Proof. Let $M=\left(h\left(e_{i}, e_{j}\right)\right)$ be the matrix of $h$ relatively to the standard basis $e_{1}, \ldots, e_{n}$ of $D_{1}^{n}$. Denote by $\rho: \mathrm{M}_{n}\left(D_{1}\right) \rightarrow \mathrm{M}_{n}\left(D_{1}\right)$ the involution ${ }^{\rho}\left(a_{i j}\right)=\left({ }^{\tau} a_{j i}\right),\left(a_{i j}\right) \in \mathrm{M}_{n}\left(D_{1}\right)$. Then ${ }^{\rho} M=\lambda M$ with $\lambda \in L_{1}$ and $\mathrm{SU}(h)=\left\{\left.g \in \mathrm{SL}_{n}\left(D_{1}\right)\right|^{\rho} g M g=M\right\}$. Let $I: \mathrm{M}_{n}\left(D_{1}\right) \rightarrow \mathrm{M}_{n}\left(D_{1}\right)$ where ${ }^{I} a=M^{-1}\left({ }^{\rho} a\right) M$ for all $a \in \mathrm{M}_{n}\left(D_{1}\right)$. Since $L_{1}$ coincides with the center of $D_{1}$, we get that $I$ is an involution of the same kind as $\rho$. Let $G$ be the Zariski closure of $\Sigma$ in $\mathrm{SL}_{n}(D)$ and $L[\Sigma]$ be the $L$-subalgebra of $\mathrm{M}_{n}(D)$ generated by $\Sigma$ over $L$. It is easy to see that $L$ is $I$-invariant and ${ }^{I} g=g^{-1}$ for each $g \in \Sigma$. Therefore the restriction of $I$ on $L[\Sigma]$ induces an involution which will be denoted also by $I$. It follows from the assumptions (a)-(c) in the formulation of Theorem 5 (see also 5.4) that the algebra $\mathrm{M}_{n}\left(D_{1}\right)$ is generated by $\mathrm{SU}(h)$ over $L_{1}$. Since $\Sigma$ is Zariski dense in $\mathrm{SU}(h)$, this implies that $\mathrm{M}_{n}\left(D_{1}\right)$ is generated by $\Sigma$ over $L_{1}$. Therefore $L[\Sigma]=\mathrm{M}_{n}(D)$. Now, the existence of the involution $I$ on $\mathrm{M}_{n}(D)$ implies the existence of an involution $\sigma$ on $D$ of the same kind as $\rho$ (and $\tau)\left[\mathrm{K} 1\right.$, Theorem, p.37]. Let $J: \mathrm{M}_{n}(D) \rightarrow \mathrm{M}_{n}(D),{ }^{J}\left(a_{i j}\right)=\left({ }^{\sigma} a_{j i}\right)$. According to Proposition 5.2, $I=J \circ \operatorname{Int}(N)$ where $N \in \mathrm{GL}_{n}(D)$ and ${ }^{J} N=\lambda N, \lambda= \pm 1$. Since $\mathrm{SU}(h)=\left\{g \in \mathrm{SL}_{n}\left(D_{1}\right) \mid{ }^{I} g=g^{-1}\right\}$, we have that $G=\left\{g \in \mathrm{SL}_{n}(D) \mid{ }^{J} g N g=N\right\}$. Therefore $G=\mathrm{SU}\left(h_{\circ}\right)$ where $h_{\circ}$ is the $\lambda$-hermitian form with respect to $\sigma$ having matrix $N$.

In view of 5.3 , having replaced $h$ by a suitable multiple, we may (and will) assume that $\tau=\sigma$ (equivalently, $I=J$ ). Therefore

$$
g^{-1}=N^{-1}\left({ }^{I} g\right) N=M^{-1}\left({ }^{I} g\right) M
$$

for all $g \in G$. So, $M N^{-1}$ commutes with each $g \in G$. Therefore $M N^{-1}$ is in the center of $M_{n}\left(D_{1}\right)$ i.e. $h=a h_{\circ}$ where $a \in D_{1}$. The proposition is proved.
5.6 Proof of Theorem 5. We will use the notations from the formulation of Theorem 5 in the Introduction. Let $\mathbf{G}_{1}$ be the $L$-algebraic group corresponding to $\mathrm{SL}_{n}(D)$, i.e $\mathbf{G}_{1}(L)=\mathrm{SL}_{n}(D)$. Let $\mathbf{G}=R_{L / K} \mathbf{G}_{1}$ where $R_{L / K}$ is Weil's restriction of scalars. Then $\mathbf{G}\left(K_{v}\right)=\operatorname{SL}_{n}\left(D_{v}\right)$ for each $v \in S$. Put $G=\mathbf{G}\left(K_{S}\right), \Gamma=\operatorname{SL}_{n}(\Lambda)$ and $H=\prod_{v \in S} \operatorname{SU}\left(h_{v}\right)$. It follows from [BTi1, 3.18] that, under the natural action of $G$ on the space of all hermitian forms on $D_{S}^{n}$, the orbit $G h$ is closed and, therefore, homeomorphic to $G / H$. Hence, the almost $S$-integer equivalence implies the proper equivalence and, also, the assertion (ii) from the formulation of Theorem 5 is equivalent to the density of $H \Gamma$ in $G$.

Let us prove that (i) implies (ii). In view of the above remark, it is enough to show that $H_{u} \Gamma$ is dense in $G$. Since $\Gamma$ is a lattice in $G$, it follows from Theorem 1 that there exists a connected $K$-algebraic group $\mathbf{L}$ of $\mathbf{G}$ and a subgroup of finite index $L^{\prime}$ in $L=\mathbf{L}\left(K_{S}\right)$ such that the closure of $H_{u} \Gamma / \Gamma$ in $G / \Gamma$ coincides with $L^{\prime} \Gamma / \Gamma, \Sigma=L^{\prime} \cap \Gamma$ is Zariski dense in $\mathbf{L}$, and $L^{\prime}$ contains $H_{u}$. Note that $\mathbf{S U}\left(h_{v}\right)$ is a maximal connected algebraic subgroup of $\mathbf{S} L_{n}(D)$ and $L^{\prime} \cap \mathrm{SU}\left(h_{v}\right)$ is Zariski dense in $\mathrm{SU}\left(h_{v}\right)$ for all $v \in S_{o}$. This implies that either $\mathbf{L}\left(K_{v}\right)=\mathrm{SU}\left(h_{v}\right)$ for all $v \in S_{o}$ or $\mathbf{L}=\mathbf{G}$. In the first case, it follows from Proposition 5.5 that there exists an hermitian form $h_{o}$ on $D^{n}$ determined by $\Sigma$ and such that $h_{S_{o}}$ is multiple of $h_{o}$. This contradicts our hypothesis. Let $\mathbf{L}=\mathbf{G}$. We will show that $L^{\prime} \Gamma=G$. This is clear when $n>1$ because $\mathbf{G}\left(K_{v}\right)$, $v \in S$, does not contain subgroups of finite index and, therefore, $L^{\prime}=G$. Let $n=1$. It follows from the general description of the orders [W1, ch. 5] that $\Gamma$ is the intersection of $G$ with an open compact subgroup of $\mathbf{G}\left(\mathcal{V}_{S}\right)$ where $\mathcal{V}_{S}$ is the $S$-adele ring (i.e. $\mathcal{V}_{S}$ is the restricted topological product of all fields $K_{v}, v \notin S$, relative to the rings of integers $\mathcal{O}_{v} \subset K_{v}$, $v \notin S)$. By the strong approximation theorem, the projection of $\Gamma$ into $\mathbf{G}\left(K_{S-S_{o}}\right)$ is dense. Therefore $L^{\prime} \Gamma=G$, which proves the implication.

Next we will prove that (ii) implies (i). Assume the contrary, that is,$h_{S_{o}}$ is a multiple of a rational form $h_{o}$ on $D^{n}$. Let $\mathbf{L}_{1}=\mathbf{S U}\left(h_{o}\right)$, $L_{1}=\mathbf{L}_{1}\left(K_{S}\right)$ and $C=\prod_{v \in S-S_{o}} \mathrm{SU}\left(h_{v}\right)$. Since $H_{u} \Gamma$ is dense in $G, L_{1} \Gamma$ is closed, $L_{1} \supset H_{u}$ and $C$ is compact, we get that $G=C L_{1} \Gamma$. Note that $H_{u}$ commutes elementwise with $C$. Therefore every $H_{u}$-ergodic
component of the Haar measure on $G / \Gamma$ is concentrated on $g L_{1} \Gamma / \Gamma$ for some $g \in C$. In particular, $H_{u}$ does not act ergodically on $G / \Gamma$. On the other hand, by a generalization of a theorem of Moore about the Mautner property [MTo2, 2.1], every $H_{u}$-invariant $L^{2}$-function on $G / \Gamma$ is invariant under the action of the smallest normal subgroup $G_{o} \subset G$ containing $H_{u}$. It is clear that $G_{o}=\mathbf{G}\left(K_{S_{o}}\right)$. Since $G_{o} \Gamma$ is dense in $G$ by the strong approximation, $G_{o}$ acts ergodically on $G / \Gamma$. Therefore, the action of $H_{u}$ is ergodic. Contradiction. The theorem is proved.
5.7 Proof of Corollary 1 Since $D^{n}$ is dense in $D_{S}^{n}$ (weak approximation), we can approximate each $x_{i}$ by a vector $y_{i} \in D^{n}$ in such a way that $y_{1}, y_{2}, \ldots, y_{n-1}$ are linearly independent over $D$ and $\left|h_{S}\left(x_{i}, x_{j}\right)-h_{S}\left(y_{i}, y_{j}\right)\right|<\frac{\varepsilon}{2}$ for all $i, j=1,2, \ldots, n-1$. Denote by $e_{1}, e_{2}, \ldots, e_{n}$ the standard basis of $D^{n}$. There exists $g \in \operatorname{SL}_{n}(D)$ such that $g e_{i}=y_{i}$ for all $i=1,2, \ldots, n-1$. Put $y_{n}=g e_{n}$ and $h^{\prime}=h_{S}^{g}$. In view of Theorem $1, h_{S}$ is almost $S$-integer equivalent to $h^{\prime}$. Hence there exists $\gamma \in \operatorname{SL}_{n}(\Lambda)$ such that $\left|h_{S}\left(x_{i}, x_{j}\right)-h_{S}\left(\gamma e_{i}, \gamma e_{j}\right)\right|<\frac{\varepsilon}{2}$ for all $i, j=1,2, \ldots, n-1$. This implies (1) with $z_{i}=\gamma e_{i}$. The Corollary is proved.
5.8 Examples and concluding remarks. 1. Let us show that the assumptions in the formulation of Theorem 5 are essential and can not be relaxed. Let $a, b \in \mathbb{Z}-\{0\}$ and $D=\{a, b\}$ be the quaternion algebra over $\mathbb{Q}$ defined by $a$ and $b$, i.e. $D=\mathbb{Q}+\mathbb{Q} i+\mathbb{Q} j+\mathbb{Q} k$ where $i^{2}=a, j^{2}=b$ and $k=i j-j i$. Assume that $D_{\infty}=D \bigotimes \mathbb{Q} \mathbb{R}$ is isomorphic to $\mathrm{M}_{2}(\mathbb{R})$. Let $\tau: D \rightarrow D$ be the standard involution of $D$, i.e. ${ }^{\tau}(x+y i+z j+t k)=$ $x-y i-z j-t k)$, and let $\Lambda=\mathbb{Z}+\mathbb{Z} i+\mathbb{Z} j j+\mathbb{Z} k$. Recall that $\tau$ acts on $D_{\infty}$ as follows :

$$
\tau\left(\begin{array}{ll}
x & y \\
z & t
\end{array}\right)=\left(\begin{array}{cc}
t & -y \\
-z & x
\end{array}\right)
$$

[Sch, p.361]. Denote by $\mathbf{G}$ the $\mathbb{Q}$-algebraic group corresponding to $\mathrm{SL}_{1}(D)\left(\right.$ i.e $\mathbf{G}(\mathbb{Q})=\mathrm{SL}_{1}(D)$ ) and put $\Gamma=\mathrm{SL}_{1}(D) \cap \Lambda$. Then $\mathbf{G}(\mathbb{R})=\mathrm{SL}_{2}(\mathbb{R})$ and, in view of a classical result of Borel and Harish-Chandra [M4, I.3.2.4], $\mathrm{SL}_{2}(\mathbb{R}) / \Gamma$ is compact if and only if $D$ is a division algebra. Let $T$ be the subgroup of all diagonal matrices in $\mathrm{SL}_{2}(\mathbb{R})$ and $K=\mathrm{SO}_{2}(\mathbb{R})$. Then $X=K \backslash \mathrm{SL}_{2}(\mathbb{R}) / \Gamma$ can be regarded as a Riemannian surface with constant curvature -1 , and the action of $T$ by left transformations on $\mathrm{SL}_{2}(\mathbb{R}) / \Gamma$ induces the geodesic flow on $X$. It is a standard fact that there exists a relatively compact, non-compact, and non-dense $T$-orbit on $X$. (We refer to [St1, Lemma 2] for a more general result due to Margulis.) Thus, there exists a $g \in \mathrm{SL}_{2}(\mathbb{R})$ such that $\Gamma g T$ is neither
dense nor closed in $\mathrm{SL}_{2}(\mathbb{R})$. Put $\alpha_{o}=\left(\begin{array}{cc}1 & 0 \\ 0 & -1\end{array}\right) \quad$ and $\alpha=g \alpha_{o} g^{-1}$. Then ${ }^{\tau} \alpha=-\alpha$ and $h(x, y)={ }^{\tau} x \alpha y$ is an isotropic -1-hermitian form on $D_{\infty}$. Since $\Gamma g T$ is not closed, $h$ is not rational, and since $\Gamma g T$ is not dense in $\mathrm{SL}_{2}(\mathbb{R})$, there are hermitian forms which are properly but not almost $S$-integer equivalent to $h$. In the case when $D$ is a division algebra (say, $\{a, b\}=\{-1,3\}$ ), we get an example showing that if $n=1$ and $r=2$ (i.e. the assumption (a) in the formulation of Theorem 1 is not fulfilled) then (i) does not imply (ii). If $D=\mathrm{M}_{2}(\mathbb{Q})$ and $\Lambda=\mathrm{M}_{2}(\mathbb{Z})$ then $\alpha=\left(\begin{array}{cc}-\gamma & -\delta \\ -\beta & \gamma\end{array}\right)$, where $\beta, \gamma, \delta \in \mathbb{R}$. It is easy to see that if the quadratic form $f(x, y)=\beta x^{2}+2 \gamma x y+\delta y^{2}$ is isotropic and irrational then the closure of $f\left(\mathbb{Z}^{2}\right)$ in $\mathbb{R}$ does not contain 0 and there exists a properly equivalent to $f$ quadratic form which is not almost $S$-integer equivalent to $f$. This shows that the assumption (c) in Theorem 5 is essential. (We refer to [M3, 1.2] and [G, 4.2] for explicit examples of quadratic forms with the same properties.) Concerning (b), note that if $r=n=1$ then $h_{S}$ is always rational.
2. We use the notations from Corollary 1. It is easy to see that if $n \geq 2$ then $h_{S}$ is anisotropic if and only if the map $D_{S}^{n} \rightarrow D_{S}$, $x \rightarrow h_{S}(x, x)$, is proper. This implies that if $h_{S}$ is anisotropic then the subset $\left\{h_{S}(z, z) \mid z \in \Lambda^{n}\right\} \subset D_{S}$ is discrete. Let $n=1$. Then $\operatorname{Nrd}\left\{h_{S}(z, z) \mid z \in \Lambda^{n}\right\}$ is discrete in $L_{S}$. This means that the assertion analogous to Corollary 1 is not true for $n=1$. Similar arguments show that Theorem 5 can not be modified to be true for "equivalent" instead of "properly equivalent" hermitian forms. (Two hermitian forms $h_{S}$ and $h_{S}^{\prime}$ are equivalent if they are conjugated by an element from $\left.\mathrm{GL}_{n}\left(D_{S}\right)=\prod_{v \in S} \mathrm{GL}_{n}\left(D_{v}\right).\right)$
3. Almost the same proofs allow to establish similar results to Theorem 5 and its corollaries when considering finite dimensional central simple algebras with involutions $\tau_{S}$ of "mixed" type (i.e. $\tau_{S}=\bigoplus_{v \in S} \tau_{v}$ where $\tau_{v}, v \in S$, is an involution of first or second type).
4. Recently Eskin, Margulis and Mozes proved the quantitative version of the Oppenheim conjecture for real quadratic forms [EMM]. It is plausible to obtain quantitative results in the general framework of the hermitian forms over division algebras in the $S$-arithmetic case.
5. Another very interesting application of the dynamical approach to the number theory is the recent proof by Kleinbock and Margulis [KM] of conjectures of Bauer and Sprindzhuk from the theory of the Diophantine approximation on manifolds. It is of interest to generalize these results to the $S$-arithmetic setting as well.

## References

[B1] Borel, A., Some finitness properties of adele groups over number fields. Publ. Math. IHES 16 (1963), 1-30.
[B2] Borel, A., Values of indefinite quadratic forms at integral points and flows on spaces of lattices. Bulletin of AMS 32 (1995), 184-204.
[B3] Borel, A., Linear Algebraic Groups. New York: W.A. Benjamin (1969).
[BPr] Borel, A., Prasad, G., Values of isotropic quadratic forms at $S$-integral points. Compositio Math. 83 (1992), 347-372.
[BTi1] Borel, A., Tits, J., Homomorphismes "abstraits" de groupes algébriques simples. Ann. of Math. 97 (1973), 499-571.
[BTi2] Borel, A., Tits, J., Groupes reductifs. Publ. Math. IHES 27 (1965), 55-150.
[CaF] Cassels, J.W.S., Fröhlich A. (editors), Algebraic Number Theory, Academic Press, (1967).
[C] Chevalley, C., Théorie des groupes de Lie II. Groupes algébraiques. Hermann, Paris, (1951).
[D1] Dani, S.G., On invariant measures, minimal sets, and a lemma of Margulis. Invent. Math. 51 (1979), 239-260.
[D2] Dani, S.G., Invariant measures and minimal sets of horospherical flows. Invent. Math. 64 (1981), 357-385.
[D3] Dani, S.G., On orbits of unipotent flows on homogeneous spaces. Ergod. Th. Dyn. Syst. 4 (1984), 25-34.
[D4] Dani, S.G., On orbits of unipotent flows on homogeneous spaces - II. Ergod. Th. Dyn. Syst. 6 (1986), 167-182.
[D5] Dani, S.G., Flows on homogeneous spaces: A review, Preprint.
[DM1] Dani, S.G., Margulis, G.A., Values of quadratic forms at primitive integral points. Invent. Math. 98 (1989), 405-424.
[DM2] Dani, S.G., Margulis, G.A., Limit distributions of unipotent flows and values of quadratic forms. Advances in Soviet Math. 16 (1993), 91-137.
[EMM] Eskin,A., Margulis, G.A., Mozes, S., Upper bounds and asymptotics in a quantitative version of the Oppenheim conjecture, Ann. of Math. 147(1998), 93-141.
[G] Ghys, E., Dynamique des flots unipotents sur les espaces homogenes. Asterisque 206 (1992), 93-136.
[KIM] Kleinbock, D., Margulis, G.A., Flows on homogeneous spaces and Diophantine approximation on manifolds, Ann. of Math.(2) 148(1998), 339-360.
[K1] Kneser, M., Lectures on Galois Cohomology. Bombay: Tata Inst. of Fund. Research, (1969).
[K2] Kneser, M., Schwache Approximation in algebraichen Gruppen. Colloque sur la théorie des groupes algébraiques, CBRM Brussels (1962), 41-52.
[L] Landherr, W., Liesche Ringe vom Typus A uber einem algebraischen Zahlkorper und hermitesche Formen uber einem Schiefkorper. Abh. Math. Sem. Hamburg 12 (1938), 200-241.
[M1] Margulis, G.A., On the actions of unipotent groups on the space of lattices.In: Gelfand I.M. (ed) Proc of the summer school on group representations. Bolyai Janos Math. Soc., Budapest, 1971, Budapest: Akadémiai Kiado (1975) 365-370.
[M2] Margulis, G.A., Formes quadratiques indéfinies et flots unipotents sur les espaces homogènes. C.R.Acad.Sci. Paris Ser 1304 (1987), 247253.
[M3] Margulis, G.A., Discrete subgroups and ergodic theory. Proc. of the conference "Number theory, trace formulas and discrete groups" in honour of A. Selberg (Oslo, 1987), 377-398.
[M4] Margulis, G.A., Discrete subgroups of semisimple Lie groups. Springer-Verlag Berlin Heidelberg (1991).
[M5] Margulis, G.A., Dynamical and ergodic properties of subgroup actions on homogeneous spaces with applications to number theory, Proc. of ICM, Kyoto, 1990. Math.Soc.ofJapan and Springer (1991), 193215.
[M6] Margulis, G.A., Oppenheim Conjecture. Preprint (1996).
[MTo1] Margulis, G.A., Tomanov, G.M., Invariant measures for actions of unipotent groups over local fields on homogeneous spaces. Invent. Math. 116 (1994), 347-392.
[MTo2] Margulis, G.A., Tomanov, G.M., Measure rigidity for almost linear groups and its applications. Journal d'Analyse Mathématique 69 (1996), 25-54.
[Op1] Oppenheim, A., The minima of indefinite quaternary quadratic forms. Proc. Nat. Acad. Sci. USA 15 (1929), 724-727.
[Op2] Oppenheim, A., The minima of indefinite quaternary quadratic forms. Ann. of Math. 32 (1931), 271-298.
[Pl] Platonov, V.P., The problem of strong approximation and the KneserTits conjecture. Math USSR Izv. 3 (1969), 1139-1147; Addendum, ibid. 4 (1970), 784-786.
[Pr1] Prasad, G., Strong approximation for semisimple groups over function fields. Ann. of Math. 105 (1977), 553-572.
[Pr2] Prasad, G., Ratner's theorem in S-arithmetic setting. In: Workshop on Lie Groups, Ergodic Theory and Geometry. Math. Sci. Res. Inst. Publ., Springer-Verlag, new York (1992).
[Rag] Raghunathan, M.S., Discrete subgroups of Lie groups. SpringerVerlag Berlin Heidelberg, (1972).
[Rat1] Ratner, M., On Raghunathan's measure conjecture. Ann. of Math. 134 (1991), 545-607.
[Rat2] Ratner, M., Raghunathan's topological conjecture and distributions of unipotent flows. Duke Math. J. 63 (1991), 235-280.
[Rat3] Ratner, M., Raghunathan's conjectures for cartesian products of real and p-adic Lie groups. Duke Math. J. 77 (1995), 275-382.
[Rat4] Ratner, M., Interactions between ergodic theory, Lie groups and number theory . Proc. of ICM, Zurich 1994, Birkhauser 94 (1995), 157182.
[Sch] Scharlau, W., Quadratic and Hermitian Forms, Springer-Verlag Berlin Heidelberg, (1985).
[Se] Serre, J.-P., A course in arithmetics. Graduate Texts in Math., Springer, New York (1973).
[St1] Starkov, A.N., Orbit structures of homogeneous flows and Raghunathan's conjecture. Usp. Mat. Nauk 45 (1990), 219-220.
[St2] Starkov, A.N., New progress in the theory of the homogeneous flows, Preprint (1996).
[Ti1] Tits, J., Algebraic and abstract simple groups. Ann. of Math. 80 (1964), 313-329.
[Ti2] Tits, J., Classification of algebraic semisimple groups, algebraic groups, and discontinuous groups. Symposium Colorado, Boulder 1965, AMS Proc. Symp. Pure Math. IX.
[To] Tomanov, G.M., Irrational Hermitian forms over division algebras. Prepublications de l'Institut Girard Desargues (URES-A 5028) 6, (1997).
[W1] Weil, A., Basic Number Theory. Springer-Verlag Berlin Heidelberg, (1967).
[W2] Weil, A., Adeles and algebraic groups. Lecture Notes, Princeton (1961).

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