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# Rational points on some pencils of conics with 6 singular fibres ${ }^{(*)}$ 

Sir Peter Swinnerton-Dyer ${ }^{(1)}$

Resumé. - Soient $k$ un corps de nombres et $c \in k$ non carré. Soient $f_{4}, f_{2}$ des polynômes homogènes en $X, Y$, de degré 4 et 2 respectivement. On donne des conditions nécessaires et suffisantes pour que l'équation

$$
U^{2}-c V^{2}=f_{4}(X, Y) f_{2}(X, Y)
$$

ait des solutions dans $k$.

$$
\begin{aligned}
& \text { AbSTRACT. - Let } k \text { be an algebraic number field, let } c \text { be a non-square } \\
& \text { in } k \text { and let } f_{4}, f_{2} \text { be homogeneous polynomials in } X, Y \text { of degrees } 4 \text { and } \\
& 2 \text { respectively. Necessary and sufficient conditions are obtained for the } \\
& \text { solubility in } k \text { of } \\
& \qquad U^{2}-c V^{2}=f_{4}(X, Y) f_{2}(X, Y) .
\end{aligned}
$$

Let $\mathcal{Y} \rightarrow \mathbf{P}^{1}$ be a pencil of conics defined over an algebraic number field $k$. It is conjectured that the only obstruction to the Hasse principle on $\mathcal{Y}$, and also to weak approximation, is the Brauer-Manin obstruction; and it was shown in [3] that this follows from Schinzel's Hypothesis. Descriptions of the Brauer-Manin obstruction and of Schinzel's Hypothesis can be found in [3]. It is of interest that arguments which show that the Brauer-Manin obstruction is the only obstruction to the Hasse principle for particular classes of $\mathcal{Y}$ normally fall into two parts:

[^0](i) the proof that some comparatively down-to-earth obstruction is the only obstruction to the Hasse principle;
(ii) the identification of that obstruction with the Brauer-Manin obstruction.

The theorem in this paper is entirely concerned with (i); the equivalence of the obstruction in the theorem with the Brauer-Manin one has already been proved in a much more general context in [1], $\S 2.6 \mathrm{~b}$ and Chapter 3.

If one does not assume Schinzel's Hypothesis, little is known. The only promising-looking line of attack is through the geometry of the universal torseurs on $\mathcal{Y}$; and these are much easier to study when $\mathcal{Y}$ has the special form

$$
\begin{equation*}
U^{2}-c V^{2}=P(W) \tag{1}
\end{equation*}
$$

where $c$ is a non-square in $k$ and $P(W)$ is a separable polynomial in $k[W]$. By writing $W=X / Y$ we can take the solubility of (1) into the equivalent (though ungeometric) problem of the solubility of

$$
\begin{equation*}
U^{2}-c V^{2}=f(X, Y) \tag{2}
\end{equation*}
$$

in $k$, where $f$ is homogeneous of even degree; here $\operatorname{deg} f$ is $1+\operatorname{deg} P$ or $\operatorname{deg} P$. The simplest non-trivial case is that of Châtelet surfaces, when $P(W)$ has degree 3 or 4 ; in this case the conjecture was proved in [2]. The object of this paper is to prove the conjecture when $\operatorname{deg} f=6$ and $f=f_{4} f_{2}$ over $k$, where $\operatorname{deg} f_{4}=4$ and $\operatorname{deg} f_{2}=2$.

Until the statement of the main theorem, we make no assumption about (2) other than that $f(X, Y)$ has even degree $n$ and no repeated factor. After multiplying $X, Y$ by suitable integers in $k$, we can assume that

$$
f(X, Y)=a \prod_{1}^{n}\left(X+\lambda_{i} Y\right)
$$

where $a$ is an integer in $k$ and the $\lambda_{i}$ are integers in $\bar{k}$; the $\lambda_{i}$ form complete sets of conjugates over $k$. For convenience we write $\gamma=\sqrt{c}$. We can assume that $\gamma$ does not lie in any $k\left(\lambda_{i}\right)$; for otherwise $f(X, Y)$ would have a nontrivial factor of the form $F^{2}-c G^{2}$ with $F, G$ in $k[X, Y]$ and we could instead consider the simpler equation

$$
U^{2}-c V^{2}=g(X, Y)=f(X, Y) /\left(F^{2}-c G^{2}\right) .
$$

We can clearly also assume that the $\lambda_{i}$ are all distinct; for otherwise we can remove a squared factor from $f(X, Y)$ and reduce to a simpler problem
which has already been solved in [2]. To avoid trivialities, we shall also rule out solutions for which each side of (2) vanishes.

Let $\mathfrak{a}_{1}, \ldots, \mathfrak{a}_{h}$ be a set of representatives for the ideal classes in $k$; then it is enough to look for solutions $u, v, x, y$ of (2) for which $x, y$ are integers whose highest common factor is some $a_{m}$. (To move from rational to integral solutions may appear unnatural; but in fact it greatly simplifies the argument which follows, because it means that our intermediate equations do not have to be homogeneous.)

LEMMA 1. - There is a finite computable list of $n$-tuples $\left(\alpha_{1}^{(r)}, \ldots, \alpha_{n}^{(r)}\right)$ not depending on $u, v, x, y$, where $\alpha_{i}^{(r)}$ is in $k\left(\lambda_{i}\right)$ and conjugacy between $\lambda_{i}$ and $\lambda_{j}$ extends to conjugacy between $\alpha_{i}^{(r)}$ and $\alpha_{j}^{(r)}$, with the following property. If (2) has a solution with $x, y$ integers whose highest common factor is some $\mathfrak{a}_{m}$, then for some $r$ the system

$$
\begin{equation*}
u_{i}^{2}-c v_{i}^{2}=\alpha_{i}^{(r)}\left(x+\lambda_{i} y\right) \quad(1 \leqslant i \leqslant n) \tag{3}
\end{equation*}
$$

has solutions with $u_{i}, v_{i}$ in $k\left(\lambda_{i}\right)$ for each $i$.
Proof. - We postulate once for all that the manipulations which follow are to be carried out in such a way as to preserve conjugacy. A prime factor $\mathfrak{p}$ of $x+\lambda_{i} y$ in $k\left(\lambda_{i}\right)$ which also divides $f(x, y) /\left(x+\lambda_{i} y\right)$ must divide

$$
a \prod_{j \neq i}\left(-\lambda_{i} y+\lambda_{j} y\right)=y^{5} a \prod_{j \neq i}\left(\lambda_{j}-\lambda_{i}\right)
$$

and for similar reasons it must divide

$$
a \prod_{j \neq i}\left(\lambda_{i} x-\lambda_{j} x\right)=-x^{5} a \prod_{j \neq i}\left(\lambda_{j}-\lambda_{i}\right)
$$

Hence it divides $a a_{m} \prod_{j \neq i}\left(\lambda_{j}-\lambda_{i}\right)$ and must therefore belong to a finite computable list; and any prime ideal not in this list which divides some $x+\lambda_{i} y$ to an odd power must split or ramify in $k\left(\lambda_{i}, \gamma\right) / k\left(\lambda_{i}\right)$. As ideals, $\left(x+\lambda_{i} y\right)=\mathfrak{b}_{i} \mathfrak{c}_{i}$ where $\mathfrak{b}_{i}$ only contains the prime ideals which either lie in the finite computable list above or ramify in $k\left(\lambda_{i}, \gamma\right) / k\left(\lambda_{i}\right)$, and every prime ideal which occurs to an odd power in $c_{i}$ must split in $k\left(\lambda_{i}, \gamma\right) / k\left(\lambda_{i}\right)$. By transferring squares from $\mathfrak{b}_{i}$ to $\boldsymbol{c}_{i}$ we can assume that each $\mathfrak{b}_{i}$ is square-free. Each $\mathfrak{b}_{i}$ belongs to a finite list independent of $x, y$, and conorm $\mathfrak{c}_{i}=\mathfrak{C}_{i} . \sigma \mathfrak{C}_{i}$ where $\mathfrak{C}_{i}$ is an ideal in $k\left(\lambda_{i}, \gamma\right)$ and $\sigma$ is the non-trivial automorphism of $k\left(\lambda_{i}, \gamma\right)$ over $k\left(\lambda_{i}\right)$. Let $\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{H}$ be a set of representatives for the ideal classes in $k\left(\lambda_{i}, \gamma\right)$; then for some $\mathfrak{A}^{(i)}$ from this list $\mathfrak{A}^{(i)} \mathfrak{C}_{i}$ is principal, say $\mathfrak{A}^{(i)} \mathfrak{C}_{i}=\left(\xi_{i}+\gamma \eta_{i}\right)$ with $\xi_{i}, \eta_{i}$ in $k\left(\lambda_{i}\right)$. Thus

$$
\left(\xi_{i}^{2}-c \eta_{i}^{2}\right)=\mathfrak{b}_{i}^{-1} \mathfrak{A}^{(i)} \sigma \mathfrak{A}^{(i)}\left(x+\lambda_{i} y\right)
$$

as ideals. This implies $\xi_{i}^{2}-c \eta_{i}^{2}=\alpha_{i}\left(x+\lambda_{i} y\right)$ where the ideal $\left(\alpha_{i}\right)$ belongs to a finite computable list; and as we can clearly vary $\alpha_{i}$ by any squared factor, this ensures the same property for $\alpha_{i}$.

Strictly speaking, the elements of our list consist of equivalence classes of $n$-tuples (where the formulation of the equivalence relation is left to the reader); but we shall need to fix which representatives we choose. However, in what follows we shall also need to know that we can take the $u_{i}, v_{i}$ to be integers without thereby imposing an uncontrolled extra factor in the $\alpha_{i}^{(r)}$. For this purpose we need the following result:

Lemma 2. - Let $K$ be an algebraic number field and $C$ a non-square in $\mathfrak{D}_{K}$. Then there exists $A=A(K, C)$ in $\mathfrak{D}_{K}$ such that if $D$ is in $\mathfrak{D}_{K}$ with

$$
\begin{equation*}
U^{2}-C V^{2}=D \tag{4}
\end{equation*}
$$

soluble in $K$, and if $A^{2} \mid D$, then (4) is soluble with $U, V$ in $\mathfrak{D}_{K}$.
Proof. - Write $L=K(\sqrt{C})$, let $\sigma$ be the non-trivial automorphism of $L / K$ and let $\mathfrak{A}_{1}, \ldots, \mathfrak{A}_{H}$ be a set of integral representatives for the ideal classes of $L$. Let $d$ be any non-zero integer of $K$ such that $u^{2}-C v^{2}=d$ for some $u, v$ in $K$, and write

$$
\left(u+C^{1 / 2} v\right)=\mathfrak{m} / \mathfrak{n}
$$

where $\mathfrak{m}, \mathfrak{n}$ are coprime ideals in $L$. Thus $\left(u-C^{1 / 2} v\right)=\sigma \mathfrak{m} / \sigma \mathfrak{n}$, so that $\sigma \mathfrak{n} \mid \mathfrak{m}$. Choose $r$ so that $\mathfrak{A}_{r} \mathfrak{n}$ is principal - say equal to (B). If $u_{1}, v_{1}$ are defined by

$$
u_{1}+C^{1 / 2} v_{1}=B\left(u+C^{1 / 2} v\right) / \sigma B
$$

then the denominator of $u_{1}+C^{1 / 2} v_{1}$ divides $\sigma \mathfrak{A}_{r}$ and $u_{1}^{2}-C v_{1}^{2}=d$. If $A$ in $K$ is divisible by $2 C^{1 / 2} \mathfrak{A}_{r} . \sigma \mathfrak{A}_{r}$ for every $r$, then $A^{2} d=\left(A u_{1}\right)^{2}-C\left(A v_{1}\right)^{2}$ where $A u_{1}$ and $A v_{1}$ are integers.

Since we can multiply each $\alpha_{i}^{(r)}$ by the square of any nonzero integer in $k\left(\lambda_{i}\right)$, subject to the preservation of conjugacy, we can assume that $\alpha_{i}^{(r)}$ is divisible by $\left(A\left(k\left(\lambda_{i}\right), c\right)\right)^{2}$ in the notation of Lemma 2 ; thus if (3) is soluble at all for given integers $x, y$ then it is soluble in integers. Moreover

$$
\prod_{i=1}^{n} \alpha_{i}^{(r)}=\left(a \prod_{i=1}^{n}\left(u_{i}^{2}-c v_{i}^{2}\right)\right) /\left(u^{2}-c v^{2}\right)
$$

so that $\prod \alpha_{i}^{(r)}=a\left(u_{(r)}^{2}-c v_{(r)}^{2}\right)$ for some $u_{(r)}, v_{(r)}$ in $k$. Conversely, any solution of (3) gives rise to a solution of (2); and for this we do not require any condition on $(x, y)$.

If the system (3) has solutions at all, it has solutions for which conjugacy between $\lambda_{i}$ and $\lambda_{j}$ extends to conjugacy between $u_{i}, v_{i}$ and $u_{j}, v_{j}$; such solutions have the form

$$
\begin{equation*}
u_{i}=\lambda_{i}^{n-1} \xi_{0}+\ldots+\xi_{n-1}, \quad v_{i}=\lambda_{i}^{n-1} \eta_{0}+\ldots+\eta_{n-1} \tag{5}
\end{equation*}
$$

for some $\xi_{\nu}, \eta_{\nu}$ in $k$. Thus we can replace (3) by the system

$$
\begin{equation*}
\left(\lambda_{i}^{n-1} X_{0}+\ldots+X_{n-1}\right)^{2}-c\left(\lambda_{i}^{n-1} Y_{0}+\ldots+Y_{n-1}\right)^{2}=\alpha_{i}^{(r)}\left(X+\lambda_{i} Y\right) \tag{6}
\end{equation*}
$$

which is to be solved in $k$. If we eliminate $X, Y$ these become $n-2$ homogeneous quadratic equations in $2 n$ variables, which give a variety $\mathcal{Y}^{(r)}$ defined over $k$. In the special case $n=4$ it was shown in [2], $\S 7$ that the $\mathcal{Y}^{(r)}$ are factors of the universal torseurs for (1); and the same argument works for all even $n>2$. However, we shall not need to know this.

In the following theorem all the statements about (2) can be trivially translated into statements about (1).

Theorem 1. - Suppose that $n=6$ and that $f(X, Y)$ in (2) has the form

$$
\begin{equation*}
f(X, Y)=f_{4}(X, Y) f_{2}(X, Y) \tag{7}
\end{equation*}
$$

where $f_{4}, f_{2}$ are defined over $k$ and have degrees 4,2 respectively. Assume also that $f(X, Y)$ has no repeated factor. If there is a $\mathcal{Y}^{(r)}$ which is soluble in every completion of $k$ then that $\mathcal{Y}^{(r)}$ is soluble in $k$; and if this holds for some $\mathcal{Y}^{(r)}$ then (2) contains a Zariski dense set of points defined over $k$.

Proof. - We first rewrite the equations for $\mathcal{Y}^{(r)}$ in a form which makes better use of the decomposition (7). We can suppose that the linear factors of $f_{4}$ are the $X+\lambda_{i} Y$ with $i=1,2,3,4$. The system (6) is equivalent to (3); but instead of (5) we now make the substitution

$$
\begin{array}{lll}
u_{i}=\lambda_{i}^{3} \xi_{0}+\ldots+\xi_{3}, & v_{i}=\lambda_{i}^{3} \eta_{0}+\ldots+\eta_{3} & \\
u_{i}=\lambda_{i} \xi_{4}+\xi_{5}, & v_{i}=\lambda_{i} \eta_{4}+\eta_{5} & \\
\hline
\end{array}
$$

in (3). Correspondingly we replace (6) by

$$
\begin{align*}
U_{i}^{2}-c V_{i}^{2}=\alpha_{i}^{(r)}\left(X+\lambda_{i} Y\right) & (i=1,2,3,4)  \tag{8}\\
\left(\lambda_{i} X_{4}+X_{5}\right)^{2}-c\left(\lambda_{i} Y_{4}+Y_{5}\right)^{2}=\alpha_{i}^{(r)}\left(X+\lambda_{i} Y\right) & (i=5,6) \tag{9}
\end{align*}
$$

where we have written

$$
U_{i}=\lambda_{i}^{3} X_{0}+\ldots+X_{3}, V_{i}=\lambda_{i}^{3} Y_{0}+\ldots+Y_{3} \quad(i=1,2,3,4)
$$

By eliminating $X, Y$ between the four equations (8), we obtain two homogeneous quadratic equations in the eight variables $U_{i}, V_{i}$; we treat these as
defining a projective variety $\mathcal{X}_{1} \subset \mathbf{P}^{7}$. The $U_{i}, V_{i}$ are not defined over $k$, but it is clear how $\operatorname{Gal}(\bar{k} / k)$ acts on them.

We can now outline the proof of the theorem. It falls naturally into three steps.
(i) $\mathcal{X}_{1}$ contains a large enough supply of lines defined over $k$.
(ii) We can choose a Zariski dense set of lines each of whose inverse images in $\mathcal{Y}^{(r)}$ is everywhere locally soluble.
(iii) $\mathcal{Y}^{(r)}$ contains a Zariski dense set of points defined over $k$.

The map $\mathcal{Y}^{(r)} \rightarrow \mathcal{Y}$ then gives the theorem.
By hypothesis, $\mathcal{X}_{1}$ has points in every completion of $k$; hence as in [2], Theorem A, there is a point $P_{0}$ in $\mathcal{X}_{1}(k)$, and we can take $P_{0}$ to be in general position on $\mathcal{X}_{1}$. Indeed, we have weak approximation on $\mathcal{X}_{1}$ because $\mathcal{X}_{1}$ contains two conjugate $\mathbf{P}^{3}$ given by

$$
X_{i} \pm \gamma Y_{i}=0 \quad(i=1,2,3,4)
$$

for either choice of sign, and these have no common point. To a general $k$-point $P$ of $\mathcal{X}_{1}$ we can in an infinity of ways find a $k$-plane which contains $P_{0}$ and $P$ and which meets both these $\mathbf{P}^{3}$; for we need only choose a $k$-point $P^{\prime}$ on $P P_{0}$ and note that since $P^{\prime}$ does not lie on either $\mathbf{P}^{3}$ there is a unique transversal from $P^{\prime}$ to the two $\mathbf{P}^{3}$. Conversely, a general $k$-plane through $P_{0}$ which meets both these $\mathbf{P}^{3}$ will meet $\mathcal{X}_{1}$ in just one more point, which must therefore be defined over $k$. In this way we obtain a map $\mathrm{P}^{6}(k) \rightarrow \mathcal{X}_{1}(k)$ which is surjective, and this implies weak approximation.

Now let $\Lambda_{0}$, which is a $\mathbf{P}^{5}$, be the tangent space to $\mathcal{X}_{1}$ at $P_{0}$, and write $\mathcal{X}_{2}=\mathcal{X}_{1} \cap \Lambda_{0}$, so that $\mathcal{X}_{2}$ is a cone whose vertex is $P_{0}$ and whose base $\mathcal{X}_{3}$ is a Del Pezzo surface of degree 4. (The fact that there are 16 lines on a nonsingular Del Pezzo surface, and the incidence relations between them, can be read off from [4], Theorem 26.2.) We can give a rather explicit description of $\mathcal{X}_{3}$, and in particular we can identify the 16 lines on it, which turn out to be distinct. Drawing on Cayley's exhaustive classification of singular cubic surfaces, a sufficiently erudite reader can derive a painless proof that $\mathcal{X}_{3}$ is actually nonsingular. (What we actually use is the much weaker statement that $\mathcal{X}_{3}$ is absolutely irreducible and not a cone, which is not hard to verify.) For $\mathcal{X}_{2}$ contains the line which is the intersection of

$$
\begin{align*}
U_{i}-\epsilon_{i} \gamma V_{i} & =0 \quad(i=1,2,3,4)  \tag{10}\\
& -336-
\end{align*}
$$

with $\Lambda_{0}$, where each $\epsilon_{i}$ is $\pm 1$. (This intersection is proper because $P_{0}$ is in general position.) We denote this line by $L^{*}\left(\epsilon_{1}, \epsilon_{2}, \epsilon_{3}, \epsilon_{4}\right)$ and its projection onto $\mathcal{X}_{3}$ by $L\left(\epsilon_{1}, \epsilon_{2}, \epsilon_{3}, \epsilon_{4}\right)$. The latter clearly meets the four lines which are obtained by changing just one sign, because this already happens for the corresponding lines in $\mathcal{X}_{2}$; so by symmetry the fifth line which it meets must be obtained by changing all four signs. This can be checked directly; for if we temporarily drop the notation of (3) and write

$$
P_{0}=\left(u_{1}, v_{1}, u_{2}, v_{2}, u_{3}, v_{3}, u_{4}, v_{4}\right) \text { in } \mathcal{X}_{0} \subset \mathbf{P}^{7}
$$

then the join of the two points $\left(\epsilon_{1} c v_{1} \pm \gamma u_{1}, \epsilon_{1} u_{1} \pm \gamma v_{1}, \ldots\right)$ passes through $P_{0}$, and each point lies on the corresponding $L^{*}\left( \pm \epsilon_{1}, \pm \epsilon_{2}, \pm \epsilon_{3}, \pm \epsilon_{4}\right)$. Since

$$
u_{i}\left(\epsilon_{i} c v_{i} \pm \gamma u_{i}\right)-c v_{i}\left(\epsilon_{i} u_{i} \pm \gamma v_{i}\right)= \pm \gamma\left(u_{i}^{2}-c v_{i}^{2}\right)
$$

and the equations for $\mathcal{X}_{1}$ are given by the vanishing of linear combinations of the $U_{i}^{2}-c V_{i}^{2}$, these two points also lie on $\Lambda_{0}$. The point

$$
P_{1}=\left(\epsilon_{1} c v_{1}, \epsilon_{1} u_{1}, \epsilon_{2} c v_{2}, \epsilon_{2} u_{2}, \epsilon_{3} c v_{3}, \epsilon_{3} u_{3}, \epsilon_{4} c v_{4}, \epsilon_{4} u_{4}\right)
$$

lies on the join of these two points; $P_{1}$ is distinct from $P_{0}$ unless $P_{0}$ lies on the $\mathbf{P}^{3}$ given by (10) or the $\mathbf{P}^{3}$ derived from it by changing the sign of $\gamma$. Because $P_{0}$ is in general position, we can assume that neither of these happens. Now a straightforward calculation, using the fact that we can describe $\mathcal{X}_{1}$ by equations which express $U_{1}^{2}-c V_{1}^{2}$ and $U_{2}^{2}-c V_{2}^{2}$ as linear combinations of $U_{3}^{2}-c V_{3}^{2}$ and $U_{4}^{2}-c V_{4}^{2}$, shows that $P_{1}$ is nonsingular on $\mathcal{X}_{2}$ unless $P_{0}$ lies on one of 12 lines, a typical one of which is given by

$$
U_{1}=V_{1}=U_{2}=V_{2}=0, U_{3}=\epsilon_{3} \gamma V_{3}, U_{4}=-\epsilon_{4} \gamma V_{4} .
$$

Under the same condition, the point induced on $\mathcal{X}_{3}$ is nonsingular.
The lines $L(++++)$ and $L(----)$ are defined over $k(\gamma)$ and conjugate over $k$; thus their intersection is defined over $k$ and $\mathcal{X}_{3}$ does contain a point defined over $k$. Moreover the $u_{i}^{2}-c v_{i}^{2}$ cannot all vanish because $\gamma$ is not in any $k\left(\lambda_{i}\right)$; so $P_{1}$ is nonsingular on $\mathcal{X}_{2}$ and $k$-points are Zariski dense on $\mathcal{X}_{3}$. (See [4], Theorems 30.1 and 29.4.) Henceforth $P_{2} \neq P_{0}$ will always denote a point on $\mathcal{X}_{2}$ defined over $k$ and $P_{3}$ will denote the corresponding point on $\mathcal{X}_{3}$.

Once we have chosen $P_{2}$, the general point of the line $P_{0} P_{2}$ is given by setting the $X_{i}, Y_{i}$ for $i=0,1,2,3$ equal to linear forms in $Z_{1}, Z_{2}$; and we can suppose that $P_{0}$ corresponds to $(1,0)$ and $P_{2}$ to $(0,1)$. The equations for $\mathcal{X}_{1}$ are then satisfied identically, and (8) expresses $X, Y$ as quadratic forms in $Z_{1}, Z_{2}$. There remain the equations (9), which now take the form

$$
\begin{equation*}
\left(\lambda_{i} X_{4}+X_{5}\right)^{2}-c\left(\lambda_{i} Y_{4}+Y_{5}\right)^{2}=\phi_{i}\left(Z_{1}, Z_{2}\right) \quad(i=5,6) \tag{11}
\end{equation*}
$$

for certain quadratic forms $\phi_{5}, \phi_{6}$. In view of the remarks in the previous paragraph we can certainly assume that $\phi_{5}, \phi_{6}$ are linearly independent and each has rank 2. We need to check that we can choose the line $P_{0} P_{2}$ so that the system (11) is everywhere locally soluble. This is of course the crucial step in the proof of the Theorem; but in order not to disrupt the flow of the argument, we postpone the proof of it and of an auxiliary result to Lemma 3 below. Given this, we would like to conclude the argument by appealing to Theorem A of [2]; but unfortunately we are in the exceptional case ( $E_{5}$ ) of that theorem. Some discussion of this exceptional case can already be found in the literature (for example in [2]); but it is not clear that any published result meets our needs. We therefore proceed as follows.

Suppose first that $\lambda_{5}, \lambda_{6}$ are in $k$ and write

$$
U_{i}=\lambda_{i} X_{4}+X_{5}, \quad V_{i}=\lambda_{i} Y_{4}+Y_{5} \quad(i=5,6)
$$

The equation (11) for $i=5$ is $U_{5}^{2}-c V_{5}^{2}=\phi_{5}\left(Z_{1}, Z_{2}\right)$, which is everywhere locally soluble, and therefore soluble by the Hasse-Minkowski theorem. Its general solution is given by homogeneous quadratic forms in three variables $W_{1}, W_{2}, W_{3}$. The equation (11) with $i=6$ now reduces to

$$
\begin{equation*}
U_{6}^{2}-c V_{6}^{2}=g\left(W_{1}, W_{2}, W_{3}\right) \tag{12}
\end{equation*}
$$

where $g$ is quartic. This is everywhere locally soluble; so all we have to do is to set $W_{3}$ equal to $e_{1} W_{1}+e_{2} W_{2}$ where $e_{1}, e_{2}$ are integers in $k$ such that

$$
\begin{equation*}
U_{6}^{2}-c V_{6}^{2}=g\left(W_{1}, W_{2}, e_{1} W_{1}+e_{2} W_{2}\right) \tag{13}
\end{equation*}
$$

is everywhere locally soluble and has no Brauer-Manin obstruction. This is not difficult. Let $\mathcal{S}$ consist of the places in $k$ which are either infinite or divide $6 c$ or either of the polynomials $g\left(W_{1}, 0, W_{3}\right)$ or $g\left(0, W_{2}, W_{3}\right)$; by means of a linear transformation on the $W_{i}$ if necessary, we can assume that neither of these expressions vanishes identically and hence $\mathcal{S}$ is finite. Solubility of (13) at the places in $\mathcal{S}$ can be ensured by local conditions on $e_{1}, e_{2}$. Choose $e_{1}$ to satisfy all these local conditions and also $g\left(1,0, e_{1}\right) \neq 0$. For the local solubility of (13) all we now have to consider are the primes in $\mathcal{S}$ and the primes $\mathfrak{p}$ which divide $g\left(1,0, e_{1}\right)$. For the former, we need only impose local conditions on $e_{2}$; for the latter it is enough to ensure that $\mathfrak{p} \mathrm{Xg}\left(0,1, e_{2}\right)$, which we can do because Norm $\mathfrak{p}>3$. Finally, $g\left(W_{1}, W_{2}, W_{3}\right)$ is the product of two absolutely irreducible quadratic forms defined over $\bar{k}$ which correspond to the linear factors of $\phi_{6}$; so it is irreducible over $k$ by Lemma 3. By Hilbert irreducibility we can ensure that $g\left(W_{1}, W_{2}, e_{1} W_{1}+e_{2} W_{2}\right)$ is irreducible over $k$; so the Châtelet equation (13) is soluble, by Theorem B of [2].

If instead $\lambda_{5}, \lambda_{6}$ are not in $k$, it follows from Lemma 3 and the linear independence of $\phi_{5}$ and $\phi_{6}$ that $\phi_{5} \phi_{6}$ is irreducible over $k$. Hence (11) is
soluble in $k$ by Theorem 12.1 of [2]. The reader can easily check that the solutions thus constructed are in general position, and therefore Zariski dense on (2).

All that remains to do is to prove the following:
Lemma 3. - If $\mathcal{Y}^{(r)}$ is everywhere locally soluble there are lines $P_{0} P_{2}$ such that (11) is everywhere locally soluble and $\phi_{i}\left(Z_{1}, Z_{2}\right)$ is irreducible over $k\left(\lambda_{i}\right)$ for $i=5,6$.

Proof. - We note first that in general $\phi_{i}$ is irreducible over $k\left(\lambda_{i}\right)$. For if we take $P_{2}$ to be $P_{1}$ and $P_{0}, P_{1}$ to have $Z$-coordinates $(1,0),(0,1)$ respectively, each $U_{i}^{2}-c V_{i}^{2}$ with $i=1,2,3,4$ is a multiple of $Z_{1}^{2}-c Z_{2}^{2}$; hence the same is true of $X$ and $Y$, and therefore of $\phi_{5}$ and $\phi_{6}$. The general assertion now follows from Hilbert's Irreducibility Theorem.

The main complication in the proof of this Lemma is that we cannot assume weak approximation on $\mathcal{X}_{3}$; indeed weak approximation is probably not even true, since the Brauer group of $\mathcal{X}_{3}$ is non-trivial. (See [5].) Let $\mathcal{S}_{1}$ be a finite set of places in $k$ containing the infinite places, all small primes and all primes dividing $2 c$, any $\mathfrak{a}_{m}$, the discriminant of $f$ or any of the $\alpha_{i}^{(r)}$. Then we can choose $P_{0}$ to be in the image of $\mathcal{Y}^{(r)}\left(k_{v}\right)$ under the map $\mathcal{Y}^{(r)} \rightarrow \mathcal{X}_{1}$ for each $v$ in $\mathcal{S}_{1}$, by weak approximation on $\mathcal{X}_{1}$. Denote by $u_{i}, v_{i}, x, y$ the values of $U_{i}, V_{i}, X, Y$ at $P_{0}$; these values depend on the particular coordinate representation of $P_{0}$ which we choose, so that we can still multiply the $u_{i}, v_{i}$ by an arbitrary $\mu \neq 0$ in $k$ and multiply $x, y$ by $\mu^{2}$. We can therefore ensure that $x, y$ are integers and that the ideal $(x, y)$ is not divisible by the square of any prime ideal outside $\mathcal{S}_{1}$. We then re-choose the $u_{i}, v_{i}$ for $i=1,2,3,4$ to satisfy (8) and be integral, which we can do by the remark immediately after the proof of Lemma 2 . This of course alters $P_{0}$, but since it leaves $x, y$ unchanged the equations (9) remain locally soluble at every place in $\mathcal{S}_{1}$. Because the old $P_{0}$ was in general position on $\mathcal{X}_{1}$, we can assume that the right hand sides of the two equations (9) do not vanish at $P_{0}$.

We do not know the quadratic forms $\phi_{5}$ and $\phi_{6}$ until we have chosen $P_{2}$. But the values of $\phi_{5}(1,0)$ and $\phi_{6}(1,0)$ as elements of $k^{*} / k^{* 2}$ only depend on $P_{0}$, for they are simply the values of the right hand sides of the two equations (9) at $P_{0}$. We can therefore properly involve these values in the argument in advance of the choice of $P_{2}$. We now have local solubility of (11) for $i=5,6$ for $Z_{2}=0$ except perhaps at primes which are not in $\mathcal{S}_{1}$ but which divide $\phi_{5}(1,0) \phi_{6}(1,0)$; let $\mathcal{S}_{2}$ be the finite set of such primes. We can delete from $\mathcal{S}_{2}$ any primes for which $c$ is a quadratic residue, for (11) is
certainly soluble at such primes. To prove the Lemma, we need only show that we can choose $P_{2}$ so that no prime $\mathfrak{p}$ in $\mathcal{S}_{2}$ divides $\phi_{5}(0,1) \phi_{6}(0,1)$.

Now let $\mathfrak{p}$ be in $\mathcal{S}_{2}$ and $\mathfrak{P}$ be any prime ideal in $k\left(\lambda_{1}, \ldots, \lambda_{4}, \gamma\right)$ which divides $\mathfrak{p}$, and use a tilde to denote reduction $\bmod \mathfrak{P}$; we have $\mathfrak{P} \| \mathfrak{p}$ because all the primes which ramify lie in $\mathcal{S}_{1}$. The two $\mathbf{P}^{3}$ given by $U_{i} \pm \tilde{\gamma} V_{i}=0$ $(i=1,2,3,4)$ are also given by $X_{i} \pm \tilde{\gamma} Y_{i}=0(i=1,2,3,4)$; so if $\tilde{P}_{0}$ lies on either of them then $\tilde{\gamma}$ would be equal to the reduction mod $\mathfrak{P}$ of the value of $\mp X_{i} / Y_{i}$ at $P_{0}$. Since the latter is an element of $k$, this would mean that $c$ would be a quadratic residue $\bmod \mathfrak{p}$ - a case which we have already ruled out. Again, if for example $\tilde{u}_{1}=\tilde{v}_{1}=\tilde{u}_{2}=\tilde{v}_{2}=0$ then $x, y$ would be divisible by $\mathfrak{P}^{2}$ and hence by $\mathfrak{p}^{2}$; and this too we have ruled out. The calculations following (10) now show that $\tilde{P}_{1}$ is nonsingular on $\tilde{\mathcal{X}}_{2}$, where $P_{1}$ is as in those calculations.

At most one pair of $\tilde{u}_{i}, \tilde{v}_{i}$ vanish; if there is such a pair, we can suppose it is given by $i=4$. The equations for $\tilde{\mathcal{X}}_{2}$ are

$$
\begin{gather*}
U_{1}^{2}-\tilde{c} V_{1}^{2}=\text { homogeneous quadratic form in } U_{3}, V_{3}, U_{4}, V_{4}  \tag{14}\\
U_{1} \tilde{u}_{1}-\tilde{c} V_{1} \tilde{v}_{1}=\text { linear form in } U_{3}, V_{3}, U_{4}, V_{4}
\end{gather*}
$$

and two similar ones involving $U_{2}$ and $V_{2}$. The equation (14) is equivalent to the vanishing of a quadratic form of rank 6 , so it cannot have a hyperplane section which is not absolutely irreducible; and it now follows easily that $\tilde{\mathcal{X}}_{2}$ is absolutely irreducible. The projection from $\tilde{\mathcal{X}}_{2}$ to the $\mathbf{P}^{3}$ with coordinates $U_{3}, V_{3}, U_{4}, V_{4}$ is generically onto. Hence there are at most $O\left(q^{2}\right)$ points in $\tilde{\mathcal{X}}_{2}\left(\mathbf{F}_{q}\right)$ for which the right hand side of (9) vanishes for $i=5$ or $i=6$. The implied constant here, like $\lambda$ below, is absolute because it depends only on the degrees of the various maps and varieties involved. Now let $P$ be the point on $\mathcal{X}_{3}$ corresponding to $P_{1}$ on $\mathcal{X}_{2}$; thus $P$ is the intersection of two lines on $\mathcal{X}_{3}$. We have already shown that $\tilde{P}$ is nonsingular for all the $p$ which still concern us. The construction in the proof of [4], Theorem 30.1 specifies a non-constant map $\psi: \mathbf{P}^{1} \rightarrow \mathcal{X}_{3}$; and the reduction $\bmod p$ of the image of $\psi$ is obtained by carrying out the corresponding construction using $\tilde{\psi}$ and $\tilde{\mathcal{X}}_{3}$, so this image has good reduction. Hence there is a point $Q$ in the image of $\psi$, defined over $k$ and such that $\tilde{Q}$ is nonsingular on $\tilde{\mathcal{X}}_{3}$ and does not lie on any of the lines of $\tilde{\mathcal{X}}_{3}$. Repeating this process using this time the construction in the proof of [4], Theorem 29.4, we obtain a map $\mathbf{P}^{2} \rightarrow \mathcal{X}_{3}$ which has good reduction mod $p$ for all relevent $\mathfrak{p}$. This lifts back to a map $\mathbf{P}^{3} \rightarrow \mathcal{X}_{2}$ which is generically onto and has good reduction $\bmod p$ for all relevent $\mathfrak{p}$. Hence there exists an absolute constant $\lambda>0$ such that $\tilde{\mathcal{X}}_{2}$ has at least $\lambda q^{3}$ points which can be lifted back to points of $\tilde{\mathcal{X}}_{0}\left(\mathbf{F}_{q}\right)$. Provided that $q$ is large enough, which we ensure by putting all small primes into $\mathcal{S}_{1}$,

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we can choose such a point $\tilde{P}_{2}$ for which the right hand sides of (9) for $i=5$ and $i=6$, reduced $\bmod \mathfrak{p}$, do not vanish. We lift this $\tilde{P}_{2}$ back to $\tilde{Q}$ on $\tilde{\mathcal{X}}_{0}$. But we have weak approximation on $\mathcal{X}_{0}$. Hence we can choose a rational point $Q$ on $\mathcal{X}_{0}$ whose reduction $\bmod p$ is $\tilde{Q}$ for each of the finitely many primes in $\mathcal{S}_{2}$. If we choose $P_{2}=\phi(Q)$ this will satisfy all our conditions.

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