Weak approximation for cubic hypersurfaces of large dimension

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We address the problem of weak approximation for general cubic hypersurfaces defined over number fields with arbitrary singular locus. In particular, weak approximation is established for the smooth locus of projective, geometrically integral, nonconical cubic hypersurfaces of dimension at least 17. The proof utilises the Hardy–Littlewood circle method and the fibration method.

1. Introduction

Let \( k \) be an algebraic number field. The possible existence and structure of \( k \)-rational points on hypersurfaces defined over \( k \) is a major theme in number theory and arithmetic geometry. Let \( X \subset \mathbb{P}^{n-1}_k \) be a variety defined over \( k \). Given a place \( \nu \) of \( k \), define \( k_\nu \) to be the completion with respect to that place. If \( X \) is smooth, recall that weak approximation holds for \( X \) if \( X(k) \neq \emptyset \) and the image of the diagonal embedding

\[
X(k) \to \prod_{\nu \in S} X(k_\nu)
\]

is dense for any finite set of places \( S \). Given a possibly singular \( X \), we shall consider weak approximation for \( X_{\text{smooth}} \), the smooth locus of \( X \).

We say that \( X \) is \( k \)-rational if there is a \( k \)-birational map \( \mathbb{P}^{n-1}_k \to X \). Weak approximation is a birational invariant of smooth integral varieties, and since weak approximation holds on \( \mathbb{P}^m_k \) for any positive integer \( m \), it must hold for any smooth \( k \)-rational variety.

A classical observation is that a quadric \( k \)-hypersurface \( Q \) with a nonsingular \( k \)-point will be \( k \)-rational, provided it is geometrically integral. Essentially this is because we can parameterise the surface by lines through the \( k \)-point. In this case, the smooth locus of the quadric satisfies weak approximation. The Hasse–Minkowski theorem implies that \( Q_{\text{smooth}}(k) \neq \emptyset \iff Q_{\text{smooth}}(k_\nu) \neq \emptyset \) for all places \( \nu \) of \( k \).

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For larger degree hypersurfaces, relatively little is known. The emergence of counterexamples to weak approximation, even when rational points are present, is an indication that the situation is much more complex. For instance, with $k = \mathbb{Q}$ we have a cubic surface from [Swinnerton-Dyer 1962],

$$x_1(x_2^2 + x_3^2) = (4x_4 - 7x_1)(x_4^2 - 2x_1^2).$$  \hspace{1cm} (1-1)$$

Ignoring the subvariety $x_1 = x_4 = 0$, this has two components over the reals: one with $x_4/x_1 \geq 7/4$, which contains infinitely many rational points, and the other with $|x_4/x_1| \leq \sqrt{2}$, which contains none. Clearly then this fails weak approximation. This counterexample can be accounted for by the Brauer–Manin obstruction. A conjecture of Colliot-Thélène [2003] states that this will be the only such obstruction for rationally connected varieties such as cubic hypersurfaces of dimension at least 2.

Suppose (once and for all) that $Y \subset \mathbb{P}^{n-1}_k$ is a geometrically integral, nonconical cubic hypersurface given by the zero locus of a cubic form $C \in k[x_1, \ldots, x_n]$. The Brauer group of $Y_{\text{smooth}}$ will be trivial if its dimension is at least 3, and the codimension of the singular locus is at least 4 (see the appendix by Colliot-Thélène in [Browning 2010]). Thus, we expect that these assumptions, together with $Y_{\text{smooth}}(k) \neq \emptyset$, imply that weak approximation holds for $Y_{\text{smooth}}$.

We should note at this point that $Y(k) \neq \emptyset$ for $n \geq 16$ by the Corollary in [Pleasants 1975]. Furthermore, $Y(K) \neq \emptyset \implies Y_{\text{smooth}}(K) \neq \emptyset$ for any field $K$, for example, by [Kollár 2002, Theorem 2.3], and so the Hasse principle holds on $Y_{\text{smooth}}$ as soon as $n \geq 16$.

Let us now consider a few cases where weak approximation for $Y_{\text{smooth}}$ is known. We shall assume that $Y_{\text{smooth}}(k_\nu) \neq \emptyset$ for all $\nu$ since otherwise the matter is trivial (this is in fact guaranteed if $n \geq 10$; see, for example, [Birch and Lewis 1960]). First a classical remark: if $Y_{\text{sing}}(k) \neq \emptyset$, where $Y_{\text{sing}}$ is the singular locus, then we can parameterise $Y$ by means of lines through a rational singular point, so $Y$ is $k$-rational. If $Y$ contains two conjugate singular points and $n \geq 7$, it follows from work of Harari [1995, §5.1; 1994]. If $Y$ contains three conjugate singular points, it is known for $n = 4$ [Coray 1976, Corollary 2; Coray and Tsfasman 1988, Theorem A] and for $n \geq 6$ [Colliot-Thélène and Salberger 1989], but counterexamples exist for $n = 5$ [ibid., Section 8]. If $n \geq 5$ and $Y$ is smooth and contains a $k$-rational line, then it follows from §5.2.2 of [Harari 1995]. Finally, Corollary 2 of [Skinner 1997] shows that if $Y$ is smooth, $n \geq 17$ is sufficient.

Note that all the results mentioned so far rely fundamentally on the shape of the singular locus of $Y$, either that it is empty or contains some arithmetic structure. The aim of this paper is to consider general cubic hypersurfaces $Y$ with arbitrary singular locus and to obtain a reasonable lower bound on the dimension required to guarantee that $Y_{\text{smooth}}$ satisfies weak approximation. Our main result is the following:
Theorem A. Let \( Y \subset \mathbb{P}_k^{n-1} \) be a geometrically integral, nonconical cubic hypersurface defined over \( k \). If \( n \geq 19 \), then \( Y_{\text{smooth}} \) satisfies weak approximation.

In a qualitative sense, this result is best possible in that the geometrically integral and nonconical assumptions cannot be eliminated. For example, we could consider the union of a line and a quadric having no nonsingular rational points, or we could take a cone over the surface (1-1).

2. Structure of the proof

Theorem A is related to the result of Skinner [1997, Corollary 2], which was obtained using the Hardy–Littlewood circle method. This is advantageous when the dimension of the singular locus is small. The circle method can also be an effective tool when the equations involved have large ‘\( h \)-invariant’. This concept was originally introduced by Davenport and Lewis [1964].

Given a cubic form \( C \), define the \( h \)-invariant of \( C \), \( h = h_k(C) \), as follows: \( h \) is the smallest positive integer such that \( C(x) \) is expressible identically as

\[
L_1(x)Q_1(x) + \cdots + L_h(x)Q_h(x),
\]

where \( L_i \) and \( Q_i \) are linear and quadratic forms, respectively, with coefficients in \( k \). Similarly, for the cubic hypersurface \( Y \), we shall define \( h_k(Y) \) to mean the \( h \)-invariant of the underlying cubic form. Finally, for a cubic polynomial \( f \), we define \( h_k(f) \) to be the \( h \)-invariant of the homogeneous cubic part of the polynomial. Clearly it is an invariant with respect to nonsingular linear transformations on \( x \) over \( k \). Also note that \( h_k(Y) \leq n \) with equality if and only if \( Y(k) = \emptyset \). Furthermore, if \( n \geq h_k(Y) + r + 1 \), there is a \( k \)-rational \( r \)-plane contained in \( Y \) given by \( L_i = 0 \) in (2-1).

Our strategy is to show weak approximation for two classes of cubics, the union of which contains all of those considered in Theorem A. The first class is cubics for which the \( h \)-invariant is sufficiently large.

Lemma 1. Suppose we have a geometrically integral, nonconical cubic hypersurface \( Y \) defined over \( k \). If \( h_k(Y) \geq 16 \), then \( Y_{\text{smooth}} \) satisfies weak approximation.

To prove this, we take a cue from the concluding remarks of [Skinner 1997] and note that to find \( k \)-rational points that are \( p \)-adically close to a \( p \)-adic point, it is sufficient to find integer points that are in specific classes modulo \( p^t \) for some integer \( t \); this is equivalent to finding integral solutions to a cubic polynomial \( f \), where the cubic part has the same \( h \)-invariant as the original cubic form. For large \( h \), this problem is tailor-made for the circle method. Indeed, we will use a mild generalisation of a previous result of Pleasants [1975] that obtains an asymptotic expression for the number of integral solutions to \( f \) in an expanding region under the assumption that \( h_k(f) \geq 16 \).
The second class of cubic hypersurfaces we consider are those for which the dimension is somewhat larger than the $h$-invariant.

**Lemma 2.** Suppose we have a geometrically integral, nonconical cubic hypersurface $Y$ defined over $k$. If $n \geq h_k(Y) + 4$, then $Y_{\text{smooth}}$ satisfies weak approximation.

This is based on the fibration method (see, for example, [Colliot-Thélène 2003] for a more general description), which reduces the question to proving weak approximation for the fibres of a particular map involving $Y$. Thanks to the assumptions of the lemma, the fibres in question are quadrics of dimension at least 3. As noted in the introduction, a quadric $Q$ will satisfy weak approximation if $Q_{\text{smooth}}(k_\nu) \neq \emptyset$ for each place $\nu$ of $k$. A well known theorem of Hasse [1923] tells us that this holds if the underlying quadratic form has rank at least 5 and $Q_{\text{smooth}}(k_\nu) \neq \emptyset$ for each real place $\nu$. We must then find conditions on $Y$ under which we can assume that for a generic fibre $Q$, this is the case. This is achieved using an elementary argument.

Lemmas 1 and 2 immediately give Theorem A.

### 3. Proof of Lemma 1

First we introduce some notation. Let $k$ be of degree $d$ over $\mathbb{Q}$, and let $\mathfrak{o}$ be the ring of integers of $k$ with $\mathbb{Z}$-basis $\omega_1, \ldots, \omega_d$. Let $m$ be an integral ideal of $\mathfrak{o}$ with $\mathbb{Z}$-basis $\tau_1, \ldots, \tau_d$.

Define $\sigma_1, \ldots, \sigma_{d_1}$ to be the distinct real embeddings of $k$ and $\sigma_{d_1+1}, \ldots, \sigma_{d_1+2d_2}$ the distinct complex embeddings such that $\sigma_{d_1+i}$ is conjugate to $\sigma_{d_1+d_2+i}$. Put $k_i$ to be the completion of $k$ with respect to the embedding $\sigma_i$ for $i = 1, \ldots, d_1 + d_2$.

Define $V$ to be the commutative $\mathbb{R}$-algebra $\bigoplus_{i=1}^{d_1+d_2} k_i \cong k \otimes_{\mathbb{Q}} \mathbb{R}$ that has dimension $d$. For an element $x \in V$, we write $\pi_i(x)$ for its projection onto the $i$-th summand (so $x = \bigoplus \pi_i(x)$). There is a canonical embedding of $k$ into $V$ given by $\alpha \to \bigoplus \sigma_i(\alpha)$, and we identify $k$ with its image in $V$. Under this image, $m$ forms a lattice in $V$, and $\tau_1, \ldots, \tau_d$ form a real basis for $V$.

We define a distance function $|\cdot|_\tau$ on $V$ as follows:

$$|x|_\tau = |x_1 \tau_1 + \cdots + x_d \tau_d|_\tau = \max_i |x_i|.$$

This extends to $V^n$ in the obvious way: if $x = (x^{(1)}, \ldots, x^{(n)}) \in V^n$, then

$$|x|_\tau = \max_j |x^{(j)}|_\tau.$$

We note that there will be some constant $c$, dependent only on $k$ and the choice of basis $\tau_1, \ldots, \tau_d$, such that

$$|\pi_i(x)| \leq c|x|_\tau \quad (3-1)$$

for all $x \in V$ and $1 \leq i \leq d_1 + d_2$ (since each $\pi_i$ is linear, this is clear).
For any point \( b \in V^n \), let \( \mathcal{B}(b) \) be the box
\[
\mathcal{B}(b) = \{ x \in V^n : |x - b|_\tau < \rho/2 \},
\]
where \( \rho \) will always be a real number \( 0 < \rho < 1 \).

For any set \( \mathcal{A} \subset V^n \) and positive real number \( P \), we define \( P\mathcal{A} \) to be the set
\[
\{ x \in V^n : P^{-1}x \in \mathcal{A} \}.
\]
Given a polynomial \( \psi(x_1, \ldots, x_n) \) defined over \( k \), we shall be interested in the quantity
\[
N_{\psi, \mathcal{A}, m}(P) = \# \{ x \in P\mathcal{A} \cap m^n : \psi(x) = 0 \}
\]
and its asymptotic behaviour as \( P \to \infty \).

We can now state the generalisation of the main theorem of [Pleasants 1975] we shall use.

**Lemma 3 (Pleasants).** Let \( m \) be an integral ideal of \( \mathcal{o} \), and let \( f(x) \) be a cubic polynomial over \( k \) with homogeneous cubic part \( C(x) \) that is not the cube of a linear form. Suppose that \( h_k(f) \geq 16 \) and that for every integral ideal \( a \) of \( \mathcal{o} \), the congruence
\[
f(x) \equiv 0 \pmod{a}
\]
has nonsingular solutions in \( m^n \). Also, let \( \xi_0 = \bigoplus_{i=1}^{d_1+d_2} \xi_i \), where each \( \xi_i \in \pi_i(V)^n \) is a nonsingular solution to \( C(x) = 0 \). Then there exists a set \( \mathcal{R} \subset V^n \) containing \( \xi_0 \) and a real constant \( c_{f, \mathcal{R}, m} > 0 \) such that
\[
N_{f, \mathcal{R}, m}(P) = c_{f, \mathcal{R}, m} P^{(n-3)d} + o(P^{(n-3)d}).
\]

**Proof.** In the case where \( m = \mathcal{o} \), this is equivalent to Lemmas 6.1, 7.1, 7.2, and 7.4 of [Pleasants 1975], which were proved using the circle method. However, all the arguments go through unchanged to prove the generalisation. Indeed, if one just changes the words ‘integral points’ to ‘elements of \( m \)’ and ‘\( \omega_1, \ldots, \omega_d \)’ to ‘\( \tau_1, \ldots, \tau_d \)’ in the relevant places, essentially all the arguments work verbatim in the same way. In terms of the circle method, the commutative algebra \( V \) does not behave differently whether \( m \) is the ring of integers or an arbitrary integral ideal, and the nontrivial algebraic number theory results required [Pleasants 1975, Section 4] were not specific to \( \mathcal{o} \). □

It is straightforward to show that a suitable \( \xi_0 \) exists, for example, the argument following [Pleasants 1975, Lemma 7.2]. Thus, Lemma 3 shows that any such cubic polynomial has infinitely many solutions in \( m^n \).

We now prove Lemma 1. Recall \( Y \) is the hypersurface associated to a rational cubic form \( C \) with \( h_k(C) \geq 16 \). As noted in the introduction, since \( Y \) is not a cone and \( C \) has at least ten variables, the congruences (3-3) have nonsingular solutions in \( \mathcal{o}^n \) for all ideals \( a \). Then it is clear upon taking \( m = \mathcal{o} \) in Lemma 3 that \( Y(k) \neq \emptyset \). Furthermore, this implies that \( Y_{smooth}(k) \neq \emptyset \) as in the introduction. Suppose
we are given $\varepsilon > 0$, any finite set of places $S$, and any set of nonsingular points
\[ \{ x_v = (x_v^{(1)}, \ldots, x_v^{(n)}) \in Y_{\text{smooth}}(k_v) : v \in S \}. \]
To show that weak approximation holds for $Y_{\text{smooth}}$, it suffices to show there exists a point
\[ x = (x^{(1)}, \ldots, x^{(n)}) \in Y_{\text{smooth}}(k) \]
such that $|x^{(i)} - x_v^{(i)}|_v < \varepsilon$ for each $i$ and every $v \in S$ (where $|\cdot|_v$ is the valuation
with respect to $v$). We follow the line of argument of [Skinner 1997, Section 5].

Let $\varepsilon < 1$, $S$, and $\{ x_v \}_{v \in S}$ be given. Write $S = S_\infty \cup S_f$, where $S_\infty$
consists of infinite places and $S_f$ consists of finite places. Without loss of generality, we can
assume that $S_\infty$ consists of all the infinite places of $k$ since there are only finitely
many of them and $Y_{\text{smooth}}(k_v) \supset Y_{\text{smooth}}(k) \neq \emptyset$ for all $v$. We may also assume that
$\text{ord}_v(x_v^{(i)}) \geq 0$ for every $i$ and every $v \in S_f$.

We can find $a = (a^{(1)}, \ldots, a^{(n)}) \in \mathfrak{o}^n$ such that $|a^{(i)} - x_v^{(i)}|_v < \varepsilon/3$ for all $i$ and
$v \in S_f$ (by the Chinese remainder theorem). Let
\[ r_v = \min_i \text{ord}_v(a^{(i)} - x_v^{(i)}), \]
and let $p_v$ be the prime ideal corresponding to $v$. Put
\[ m = \prod_{v \in S_f} p_v^{r_v}. \]

Consider $f(x) = C(x + a)$, a cubic polynomial defined over $k$. Let $t$ be a positive
integer. Choose $D \equiv 1 \pmod{m^t}$ to be a positive integer such that
\[ D > \frac{2c}{\varepsilon} \]
with $c$ as in (3-1). For each infinite place $v$, let
\[ r_v = Dx_v. \]
Put
\[ \zeta_0 = \bigoplus_{i=1}^{d_1+d_2} r_{v_i} \in V^n, \]
where $v_i$ is the infinite place corresponding to the embedding $\sigma_i$. Note that $\zeta_0$
satisfies the conditions of Lemma 3. Take a set $\mathcal{R}$ as in Lemma 3 centred at $\zeta_0$. In
[Pleasants 1975], the region $\mathcal{R}$ is essentially a box-like shape, and the only extra
condition it needs to have is that it is sufficiently small. Therefore, we can take its
‘diameter’ with respect to $|\cdot|_\tau$ to be as small as we like, and we can assume it is
contained inside a box of side length $\rho < 1$ as in (3-2).

We consider the congruence conditions (3-3). For any finite place $v \notin S_f$, we
have $\mathfrak{o}_v = \mathfrak{m}_v$, and so any point in $Y_{\text{smooth}}(k_v)$ will give rise to a nonsingular solution
in $\mathfrak{m}_v^n$ of $f(x) = 0$. On the other hand, if $v \in S_f$, then $x_v - a \in \mathfrak{m}_v^n$ is a nonsingular
solution to $f(x) = 0$. Thus, the conditions hold for all integral ideals $a$ in $k$ by the Chinese remainder theorem.

Finally, we note that the cubic part of $f$ is just $C$ and $h_k(C) \geq 16$. Thus, the conditions for Lemma 3 are met, so for a sufficiently large integer $P \equiv 1 \pmod{m'}$, there exists a point $y \in m'' \cap P\mathfrak{N}$ that is a zero of $f$, and thus, $z = y + a$ is a zero of $C$. Also, since we can have arbitrarily many such points, we can choose one such that $z \neq 0$. We now fix our point $x$ to be $z/(DP)$.

For $\nu \in S_\infty$ we have

$$|DPx^{(i)}_\nu - y^{(i)}|_{\nu} \leq \rho c P < c P,$$

whence

$$|x^{(i)}_\nu - x^{(i)}|_{\nu} = \left| \frac{z^{(i)}}{DP} \right|_{\nu} \leq \left| \frac{y^{(i)}}{DP} \right|_{\nu} + \left| \frac{a^{(i)}}{DP} \right|_{\nu} \leq \frac{c}{D} + \left| \frac{a^{(i)}}{DP} \right| < \varepsilon$$

for $P$ sufficiently large.

For every $\nu \in S_f$ and when $DP \equiv 1 \pmod{m'}$ for sufficiently large $t$, we have

$$|x^{(i)}_\nu - x^{(i)}|_{\nu} = \left| \frac{z^{(i)}}{DP} \right|_{\nu} \leq |x^{(i)}_\nu - z^{(i)}|_{\nu} + \frac{|DP - 1|_\nu}{|DP|_\nu} |z^{(i)}|_{\nu} \leq |x^{(i)}_\nu - z^{(i)}|_{\nu} + 2^{-t} = |(z^{(i)} - a^{(i)}) - (x^{(i)} - a^{(i)})|_{\nu} + 2^{-t} \leq \frac{2}{3} \varepsilon + 2^{-t} < \varepsilon.$$

Finally, we note that by taking $\varepsilon$ sufficiently small, we can make $x$ be arbitrarily close to a nonsingular point on $Y(k_\nu)$ for each $\nu \in S$. In this way, we may clearly assume that $x \in Y_{smooth}(k)$. This proves Lemma 1. $\square$

4. Proof of Lemma 2

Throughout this section, we shall suppose that $h = h_k(Y)$ and that $n \geq h + 4$. In fact, for simplification, we shall only consider the case $n = h + 4$, other cases being handled similarly. After a change of variables if necessary, we can express the cubic form $C$ in terms of variables $(x, y) = (x_1, \ldots, x_4, y_1, \ldots, y_h)$ as follows:

$$Y : C(x, y) = y_1 Q_1(x_1, \ldots, x_4, y_1, \ldots, y_h) + \cdots + y_h Q_h(x_1, \ldots, x_4, y_1, \ldots, y_h) = 0,$$
where the $Q_i$ are quadratic forms defined over $k$. Clearly $Y(k) \neq \emptyset$ since $h < n$, so $Y_{\text{smooth}}(k) \neq \emptyset$ as in the introduction.

Consider the three-dimensional linear space $L$ given by $y_1 = \cdots = y_h = 0$. Take the blow-up $W$ of $Y$ along $L$. Let $z_1, \ldots, z_h$ be coordinates for $\mathbb{A}^h_k$. Then the variety given by the vanishing of $C$ and $y_i z_j - y_j z_i$ for all $i, j \in \{1, \ldots, h\}$ is $L \cup W$. Our plan is to prove weak approximation for $W_{\text{smooth}}$, which, by birationality, will prove it for $Y_{\text{smooth}}$.

Let $\pi : W \to \mathbb{A}^h_k$ be the projection $(x, y, z) \mapsto z$. The generic fibres of this map are quadrics in $\mathbb{P}^4_k$. Abusing our notation slightly, suppose that for $y \in \mathbb{A}^h_k \setminus \{0\}$, the fibre of $W/\mathbb{A}^h_k$ at the point $y$ is given by the quadric

$$X_y : Q_y(x, t) = 0.$$  \hfill (4-1)

Then we see that

$$Q_y(x, t) = C(x, yt).$$  \hfill (4-2)

For an alternative description of the fibres, consider a generic four-dimensional linear space $L_4$ that contains $L$. Then this cuts out on $Y$ the union of $L$ and one such quadric $Q$.

We may write $Y : C(x, y) = \sum_{i=1}^h y_i Q_i(x) + 2 \sum_{j=1}^4 x_j q_j(y) + c(y),$ where $Q_i$ and $q_j$ are quadratic forms and $c$ is a cubic form. Then by (4-2), $Q_y$ takes the form

$$Q_y(x, t) = \sum_{i=1}^h y_i Q_i(x) + 2 \sum_{j=1}^4 x_j q_j(y)t + c(y)t^2.$$  \hfill (4-3)

We rewrite this in terms of a $5 \times 5$ matrix $A = A(y)$ defined by

$$Q_y(x, t) = (x, t)^T A(x, t).$$
Also we consider the $4 \times 4$ submatrices $M_i$ given by $Q_i(x) = x^T M_i x$.

By (4-3), $A$ must take the shape

$$A(y) = \begin{pmatrix} \sum_i y_i M_i & q_1(y) \\ q_1(y) & \vdots & q_4(y) & c(y) \end{pmatrix}. \quad (4-4)$$

Now we can reduce Lemma 2 to the following:

**Lemma 5.** Suppose $Y$ is as in Lemma 2. For $y \in \mathbb{A}_k^h$, we define $A(y)$ as in (4-4). Either $A(y)$ is generically of full rank, or $Y_{\text{sing}}(k) \neq \emptyset$.

If $A(y)$ is generically of full rank at least 5, then $Q_y$ is generically smooth and of dimension at least 3, so we apply Lemma 4 with $Z = W$ and $X = \mathbb{A}_k^h$ to prove weak approximation for $W_{\text{smooth}}$ and hence $Y_{\text{smooth}}$. In the alternative, $Y_{\text{sing}}(k) \neq \emptyset$, which is sufficient for weak approximation as noted in the introduction.

We now establish Lemma 5. Note that the equation $\det[A(y)] = 0$ defines an algebraic set, say, in $\mathbb{A}_k^h$, so if $A(y)$ is not generically of full rank, then $\det[A(y)] = 0$ identically in $y$.

We examine the leading $4 \times 4$ submatrix $M = M(y) = \sum_i y_i M_i$.

**Lemma 6.** Let $M_1, \ldots, M_h$ be $m \times m$ symmetric matrices defined over an arbitrary field $k$ with char$(k) \neq 2$. For $(y_1, \ldots, y_h) \in k^h$, put $M(y_1, \ldots, y_h) = \sum_i y_i M_i$. Either there exists a vector $v \in \mathbb{A}_k^m$ such that $v^T M_i v = 0$ for each $1 \leq i \leq h$, or the polynomial $\det[M(y)]$ is not identically zero in $y$.

**Proof.** This is essentially contained in [Colliot-Thélène et al. 1987, Lemma 1.14], which deals with the case $h = 2$, but there are enough changes to warrant giving detail. We proceed by induction. Suppose that $\det[M(y)] = 0$ identically. If $h = 1$, then the lemma follows from the fact that $\det(M_1) = 0$ implies that $M_1$ has a nonzero null space. Now assume that $h \geq 2$ and that the lemma is true for smaller values of $h$. If $M_1 = \cdots = M_h = 0$, then the lemma is obvious. So we assume that $M_1 \neq 0$. Inserting $y = (1, 0, \ldots, 0)$, we see that $\det(M_1) = 0$. Since char$(k) \neq 2$, we can choose a basis $v_1, \ldots, v_m$ of $\mathbb{A}_k^m$ such that $M_1 = \text{diag}(a_1, \ldots, a_r, 0, \ldots, 0)$ with each $a_i \neq 0$ for some $0 < r < m$. Let $M'_2, \ldots, M'_h$ be the $(m-r) \times (m-r)$ symmetric matrices corresponding to the basis elements $v_{r+1}, \ldots, v_m$. Then for fixed $y_2, \ldots, y_h$, the coefficient of $y'_1$ in $\det[M(y)]$ is $a_1 \cdots a_r \det[\sum_{i=2}^h y_i M'_i]$. Hence, by assumption this must vanish identically. By the induction hypothesis, this implies that there is a nonzero vector $w = (b_{r+1}, \ldots, b_m) \in \mathbb{A}_k^{m-r}$ such that $w^T M'_i w = 0$ for $2 \leq i \leq h$. Then $v = (0, \ldots, 0, b_{r+1}, \ldots, b_m)$ satisfies $v^T M_i v = 0$ for $1 \leq i \leq h$. \qed
Now we prove Lemma 5. First suppose that $\det[M(y)]$ is identically zero. Then we apply the previous lemma with $m = 4$ to show that there is a nonzero $v \in \mathbb{A}^h_k$ such that $v^TM_i v = 0$ for $1 \leq i \leq h$, i.e., $Q_i(v) = 0$. We can assume after a change of variables that this vector is $(1, 0, 0, 0)$. Now note that this implies there is no $x_i^2$ term in any of the $Q_i$. This implies that all terms in $C(x, y)$ are at most linear in $x_1$. But then $Y_{\text{sing}}(k) \neq \emptyset$ since it contains the point $(1, 0, \ldots, 0)$, so we are done.

Next, we suppose that $\det[M(y)]$ is not identically zero. In particular, after a change of variables involving only $y_1, \ldots, y_h$, we can assume that it is not zero at $(1, 0, \ldots, 0)$, so $\det(M_1) \neq 0$. Now applying another change of variables involving only $x_1, \ldots, x_4$, we can assume that $M_1 = \text{diag}(a_1, \ldots, a_4)$ with $a_j \in k^\times$.

We write
\[
q_j(y) = 2d_jy_1^2 + \cdots \quad \text{for } 1 \leq j \leq 4
\]
and
\[
c(y) = ey_1^3 + y_1^2L(y_2, \ldots, y_h) + \cdots,
\]
with $L$ a linear form defined over $k$.

Assume that $\det[A(y)]$ is identically zero. Now
\[
C(x, y_1, 0 \ldots, 0) = y_1\left(\sum_{j=1}^4 a_jx_j^2\right) + \sum_{j=1}^4 2d_jx_jy_1^2 + ey_1^3
\]
\[
= y_1\sum_{j=1}^4 a_j\left(\sum_{j=1}^4 x_j + \frac{d_j}{a_j}y_1\right)^2 + e'y_1^3
\]
for some $e' \in k$. The invertible linear change of variables $x'_j = x_j + (d_j/a_j)y_1$ shows that we can assume that each $d_j = 0$. The coefficient of $y_1^7$ in $\det[A(y)]$ is $a_1 \cdots a_4e$, which must be zero; hence, $e = 0$. The coefficient of $y_1^6$ in $\det[A(y)]$ is $a_1 \cdots a_4L(y_2, \ldots, y_h)$, which is also identically zero in $y_2, \ldots, y_h$; hence, $L = 0$. Now we see that all terms in $C$ are at most linear in $y_1$; consequently, the point $(x_0, y_0) = (0, 0, 0, 0; 1, 0, \ldots, 0)$ lies in $Y_{\text{sing}}(k)$, completing the proof of the lemma.

\[\square\]

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References


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<table>
<thead>
<tr>
<th>Title</th>
<th>Page</th>
</tr>
</thead>
<tbody>
<tr>
<td>On the discrete logarithm problem in elliptic curves II</td>
<td>1281</td>
</tr>
<tr>
<td>Claus Diem</td>
<td></td>
</tr>
<tr>
<td>Identifying Frobenius elements in Galois groups</td>
<td>1325</td>
</tr>
<tr>
<td>Tim Dokchitser and Vladimir Dokchitser</td>
<td></td>
</tr>
<tr>
<td>Weak approximation for cubic hypersurfaces of large dimension</td>
<td>1353</td>
</tr>
<tr>
<td>Mike Swarbrick Jones</td>
<td></td>
</tr>
<tr>
<td>The Picard crossed module of a braided tensor category</td>
<td>1365</td>
</tr>
<tr>
<td>Alexei Davydov and Dmitri Nikshych</td>
<td></td>
</tr>
<tr>
<td>A Gross–Zagier formula for quaternion algebras over totally real fields</td>
<td>1405</td>
</tr>
<tr>
<td>Eyal Z. Goren and Kristin E. Lauter</td>
<td></td>
</tr>
<tr>
<td>Counting rational points over number fields on a singular cubic surface</td>
<td>1451</td>
</tr>
<tr>
<td>Christopher Frei</td>
<td></td>
</tr>
<tr>
<td>On the ample cone of a rational surface with an anticanonical cycle</td>
<td>1481</td>
</tr>
<tr>
<td>Robert Friedman</td>
<td></td>
</tr>
<tr>
<td>Commuting involutions of Lie algebras, commuting varieties, and simple Jordan algebras</td>
<td>1505</td>
</tr>
<tr>
<td>Dmitri I. Panyushev</td>
<td></td>
</tr>
</tbody>
</table>