# An extension to the Hermite-Joubert problem 

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## Abstract

Let $E / F$ be a field extension of degree $n$. A classical problem is to find a generating element in $E$ whose characteristic polynomial over $F$ is as simple is possible. An 1861 theorem of Ch. Hermite [5] asserts that for every separable field $E / F$ of degree $n$ there exists an element $a \in E$ whose characteristic polynomial is of the form

$$
f(x)=x^{5}+b_{2} x^{3}+b_{4} x+b_{5}
$$

or equivalently, $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{3}\right)=0$. A similar result for extensions of degree 6 was proven by P. Joubert in 1867; see [6].

In this thesis we ask if these results can be extended to field extensions of larger degree. Specifically, we give a necessary and sufficient condition for a field $F$, a prime $p$ and an integer $n \geq 3$ to have the following property: Every separable field extension $E / F$ of degree $n$ contains an element $a \in E$ such that $a$ generates $E$ over $F$, and $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{p}\right)=0$.

As a corollary we show for infinitely many new values of $n$ that the theorems of Hermite and Joubert do not extend to field extensions of degree $n$. We conjecture the same for more values of $n$ and provide computational evidence for a large number of these.

## Preface

Much of this thesis is based on an earlier joint paper [1] by Reichstein and myself. Theorem 1.3 is covered in [1] in a higher generality and some of the introduction is borrowed from [1]. Chapters 2 and 3 of this thesis are taken and modified from [1]. Theorem 4.1 appears in [1]. The rest of Chapter 4 is originial to this thesis. Some of the results of Chapter 5 appear in [1], others are original to this thesis and expand upon these results. Further results in Chapter 6 to support Conjecture 1.5 are original to this thesis, as well as the code appearing in Appendix A,

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## Chapter 1

## Introduction

Let $E / F$ be a field extension of degree $n$. A classical problem is to find a generating element in $E$ whose characteristic polynomial

$$
f(x)=x^{n}+b_{1} x^{n-1}+\cdots+b_{n-1} x+n_{n}
$$

over $F$ is as simple as possible. For example, any separable quadratic extension $F[x] /\left(x^{2}+b x+c\right)$ over a field $F$ of characteristic $\neq 2$ is generated by $\sqrt{b^{2}-4 c}$, which has characteristic polynomial $x^{2}+\left(4 c-b^{2}\right)$. Similarly, any separable cubic extension of a field of characteristic $\neq 3$ is generated by an element whose characteristic polynomial is of the form $x^{3}+b x+c$ for some $b, c \in F$.

An 1861 theorem of Ch. Hermite [5] asserts that for every separable field extension $E / F$ of degree 5 there exists an element $a \in E$ whose characteristic polynomial is of the form

$$
f(x)=x^{5}+b_{2} x^{3}+b_{4} x+b_{5} .
$$

A similar result for extensions of degree 6 was proven by P. Joubert in 1867; see [6]. For modern proofs of these results, see $[3,8]$.

Let $a \in E$ and $a_{1}, \ldots, a_{n}$ be the conjugates of $a$ (with possible repetitions if the degree of $a$ is less than $n$ ). Let $s_{i}$ denote the elementary symmetric polynomials in $n$ variables, that is, the sum of monomials of degree $i$. Then the characteristic polynomial of $a$ over $F$ is given by

$$
x^{n}-s_{1}\left(a_{1}, \ldots, a_{n}\right)+\cdots+(-1)^{n} s_{n}\left(a_{1}, \ldots, a_{n}\right)
$$

The above results of Hermite and Joubert can then be phrased as the existence of an element $a$ in
$E$ of degree $n$ over $F$ such that its conjugates satisfy $s_{1}\left(a_{1}, \ldots, a_{n}\right)=s_{3}\left(a_{1}, \ldots, a_{n}\right)=0$. Denote by $t_{i}$ the $i^{\text {th }}$ power sum in $n$ variables, that is,

$$
t_{i}\left(a_{1}, \ldots a_{n}\right)=a_{1}^{i}+\cdots+a_{n}^{i}
$$

Newton's identities yield the following:

$$
\begin{aligned}
& t_{1}=s_{1} \\
& t_{2}=s_{1}^{2}-2 s_{2} \\
& t_{3}=s_{1}^{3}-3 s_{1} s_{2}+3 s_{3} .
\end{aligned}
$$

This implies we have $s_{1}=s_{3}=0$ if and only if $t_{1}=t_{3}=0$. If $E=F[\alpha]$, for each $a \in E$ we consider the linear transformation $x \mapsto a x$ with basis $1, \alpha, \alpha^{2}, \ldots, \alpha^{n-1}$. We have $t_{i}\left(a_{1}, \ldots, a_{n}\right)=$ $a_{1}^{i}+\cdots+a_{n}^{i}=\operatorname{tr}_{E / F}\left(a^{i}\right)$, where the trace here is the trace of the linear transformation $x \mapsto a x$. Note that the trace depends only on the field extension $E / F$ and not on the choice of basis. Coray's approach to the Hermite-Joubert problem studied the existence of $F$-points in the projective variety in $\mathbb{P}^{n-1}$ defined by $\operatorname{tr}_{E / F}(x)=\operatorname{tr}_{E / F}\left(x^{3}\right)=0$, with coordinates being a basis $1, \alpha, \ldots, \alpha^{n-1}$ of $E$ over $F$.

In this thesis we will consider a generalization of the Hermite-Joubert problem: For a fixed $n \geq 7$, does every separable field extension $E / F$ of degree $n$ contain an element $a \in E$ which generates $E$ over $F$ and satisfies $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{3}\right)=0$ ? The answer is "no" if $n$ is of the form $3^{k}$ or $3^{k_{1}}+3^{k_{2}}$ where $k_{1}>k_{2} \geq 0$, see [11, Theorem 1.3]. We will see that the answer remains "no" for $n$ of the form $3^{k_{1}}+3^{k_{2}}+3^{k_{3}}$ where $k_{1}>k_{2}>k_{3} \geq 0$, under additional assumptions on the integers $k_{1}, k_{2}, k_{3}$, and we conjecture that the additional assumptions are unneeded. For other values of $n$ (in particular, for $n=7$ ), the question remains open. One can also ask a similar question for an arbitrary prime $p$.

Question 1.1. Let $n \geq 2$ be an integer, $p$ be a prime, and $F_{0}$ be a base field. Which triples ( $F_{0}, p, n$ ) have the property that for every field extension $F / F_{0}$ and every étale algebra $E / F$ of degree $n$, there exists an element $a \in E$ such that $a$ generates $E$ as an $F$-algebra and $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{p}\right)=0$ ?

Here an étale algebra $E / F$ is defined to be a finite product of separable field extensions over $F$. The degree of an étale algebra $E / F$ is the dimension of $E$ as an $F$-algebra. For example, any separable polynomial $f$ in $x$ gives rise to an étale algebra $F[x] / f$. If $f=f_{1} \cdots f_{n}$, with each $f_{i}$ irreducible and separable, then $F[x] / f \cong F[x] / f_{1} \times \cdots \times F[x] / f_{n}$.

In this paper we will show that this question becomes tractable if we restrict our attention to the case where $F$ is a $p$-field. A field $F$ is called a $p$-field if the degree of every finite field extension of $F$ is a power of $p$. If $F$ is a $p$-field we say that $F$ is $p$-closed. Then our question becomes:

Question 1.2. Let $n \geq 2$ be an integer, $p$ be a prime, and $F_{0}$ be a base field. Which triples ( $F_{0}, p, n$ ) have the property that for every $p$-field $F$ containing $F_{0}$, and every étale algebra $E / F$ of degree $n$, there exists an element $a \in E$ such that $a$ generates $E$ as an $F$-algebra and $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{p}\right)=0$ ?

Before stating our main results, we recall the definition of the "general field extension" $E_{n} / F_{n}$ of degree $n$. Let $F_{0}$ be a base field and $x_{1}, \ldots, x_{n}$ be independent variables. Set $L_{n}:=F_{0}\left(x_{1}, \ldots, x_{n}\right)$, $F_{n}:=L_{n}^{S_{n}}$ and $E_{n}:=L_{n}^{S_{n-1}}=F_{n}\left(x_{1}\right)$, where $\mathrm{S}_{n}$ acts on $L_{n}$ by permuting $x_{1}, \ldots, x_{n}$ and $\mathrm{S}_{n-1}$ by permuting $x_{2}, \ldots, x_{n}$.

Let $p$ be a prime. We will say that $n=p^{k_{1}}+p^{k_{2}}+\cdots+p^{k_{m}}$ is the base $p$ expansion of $n$ if each power of $p$ appears in the sum at most $p-1$ times. It is well known that the base $p$ expansion of $n$ is unique.

Theorem 1.3. Let $p$ be a prime, $F_{0}$ be a field of characteristic $\neq 2,3$ or $p$ containing a primitive $p^{\text {th }}$ root of unity $\zeta_{p}$, and $n=p^{k_{1}}+\cdots+p^{k_{m}}$ be the base $p$ expansion of an integer $n \geq 3$. Then the following conditions are equivalent.
(1) For every p-field $F$ containing $F_{0}$ and every $n$-dimensional étale algebra $E / F$, there exists an element $0 \neq a \in E$ such that $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{p}\right)=0$.
(2) There exists a finite field extension $F^{\prime} / F_{n}$ of degree prime to $p$ and an element $0 \neq a \in E^{\prime}:=$ $E_{n} \otimes_{F_{n}} F^{\prime}$ such that $\operatorname{tr}_{E^{\prime} / F^{\prime}}(a)=\operatorname{tr}_{E^{\prime} / F^{\prime}}\left(a^{p}\right)=0$.
(3) The system of equations

$$
\begin{align*}
p^{k_{1}} y_{1}+\cdots+p^{k_{m}} y_{m} & =0 \\
p^{k_{1}} y_{1}^{p}+\cdots+p^{k_{m}} y_{m}^{p} & =0 \tag{1.1}
\end{align*}
$$

has a solution $y=\left(y_{1}: \cdots: y_{m}\right) \in \mathbb{P}^{m-1}\left(F_{0}\right)$.
(*) Moreover, assume that condition (3) holds and one of the following additional conditions is met:
(i) There is a solution $y=\left(y_{1}: \cdots: y_{m}\right) \in \mathbb{P}^{m-1}\left(F_{0}\right)$ to (1.1) such that $y \neq(1: \ldots: 1)$, or
(ii) $p>2$.

Then the element a in parts (1) and (2) can be chosen so that $E=F[a]$ and $E^{\prime}=F^{\prime}[a]$, respectively.
Note that conditions $(i)$ and $(i i)$ of assertion $(*)$ are very mild. In particular, condition $(i)$ is satisfied whenever char $F_{0}$ does not divide $n$, since in this case $(1: \cdots: 1)$ is not a solution to (1.1). We will prove Theorem 1.3 in Chapter 2. Sections 2.1, 2.2 and 2.3 will cover the main statement of the theorem. The assertion $(*)$ will be proven in section 2.6 , after establishing preliminary results in sections 2.4 and 2.5 .

It is natural to ask for which $n, p$ and $F_{0}$ there exist solutions to the system (1.1). We will give some partial results to this question in Chapter 3. In particular, we show that for $p \geq 3$ the system (1.1) has a solution over any field $F_{0}$, if when we write $n=\left[n_{d} n_{d-1} \ldots n_{0}\right]_{p}$ in base $p$, one of the digits $n_{i}$ is $\geq 2$, or if the number of non zero digits is $\geq p+3$. (Here $n_{i}$ is the number of times $p^{i}$ occurs in the base $p$ expansion $n=p^{k_{1}}+\cdots+p^{k_{m}}$. If each $n_{i}$ is 0 or 1 , then the number of non zero digits is $m$.)

We will also see that the system (1.1) has no solutions when $n=p^{k}$ for any $k \geq 1$, or when $n=p^{k_{1}}+p^{k_{2}}$ and char $F_{0}$ does not divide $p^{\left(k_{1}-k_{2}\right)(p-1)}+(-1)^{p}$, for any $k_{1}>k_{2} \geq 0$.

We will give more detailed results for particular cases of $n$ when $p=2$ and $p=3$ in Chapters 4 and 5 respectively. Due to a theorem of Springer theorem, proven in [9, VII. Theorem 2.7], we can remove the requirement of $F$ being a 2 -field in condition (1) in the case of $p=2$. This leads to a nearly complete solution of Question 1.1 when $p=2$.

If we assume an analogous conjecture of Cassels and Swinnerton-Dyer [2] then we could also remove the requirement of $F$ being a 3 -field in the case of $p=3$. We obtain the following result for $p=3$ when $n$ has 3 digits base 3 .

Theorem 1.4. Let $E_{n} / F_{n}$ be the general field extension of degree $n$, over the base field $F_{0}=\mathbb{Q}$. Suppose $n=3^{k_{1}}+3^{k_{2}}+3^{k_{3}}$, where $k_{1}>k_{2}>k_{3} \geq 0$. If $k_{1}+k_{2}+k_{3} \not \equiv 2(\bmod 3)$ or $k_{1} \not \equiv k_{2}$ $(\bmod 3)$, then for any finite field extension $F^{\prime} / F_{n}$ of degree prime to 3 there does not exist an element $0 \neq a \in E_{n}^{\prime}:=E_{n} \otimes_{F_{n}} F^{\prime}$ such that $\operatorname{tr}_{E^{\prime} / F^{\prime}}(a)=\operatorname{tr}_{E^{\prime} / F^{\prime}}\left(a^{3}\right)=0$.

We provide a proof of this theorem in Chapter 5. This yields new examples where the HermiteJoubert problem has a negative solution, the smallest being $n=13=3^{2}+3^{1}+3^{0}$. In Chapter 6 we provide some evidence in support of the following conjecture.

Conjecture 1.5. Theorem 1.4 remains true for all triples $k_{1}>k_{2}>k_{3} \geq 0$.
The maple code used to establish this evidence is provided in Appendix A.

Much of this thesis is based on an earlier joint paper of Reichstein and myself. Theorem 1.3 is covered in [1] in a higher generality, and Chapters 2 and 3 of this thesis are taken from [1]. Chapters 4 and 5 expand upon the paper, and further results to support Conjecture 1.5 are original to this thesis.

## Chapter 2

## Proof of Theorem 1.3

### 2.1 Proof of Theorem 1.3: $(1) \Longrightarrow(2)$

For every field $F$ there exists an algebraic extension $F^{(p)} / F$, such that $F^{(p)}$ is $p$-closed, and for every finite subextension $F \subseteq F^{\prime} \subseteq F^{(p)}$, the degree $\left[F^{\prime}: F\right]$ is prime to $p$. We call this field $F^{(p)}$ the $p$-closure of $F$. This field is unique with respect to these conditions up to $F$-isomorphism. For a proof of its existence, see [4, Proposition 101.16].

Lemma 2.1.1. Let $E / F$ be an étale algebra of degree $n$.
(a) Every element $a \in E \otimes_{F} F^{(p)}$ lies in the image of the natural map

$$
\phi: E \otimes_{F} F^{\prime} \hookrightarrow E \otimes_{F} F^{(p)}
$$

for some intermediate field $F \subset F^{\prime} \subset F^{(p)}$ (depending on a), such that $\left[F^{\prime}: F\right]$ is finite (and thus prime to $p$ ).
(b) $a \in E^{\prime}:=E \otimes_{F} F^{\prime}$ generates $E^{\prime}$ over $F^{\prime}$ if and only if $\phi(a)$ generates $E \otimes_{F} F^{(p)}$ over $F^{(p)}$.

Proof. (a) Let $e_{1}, \ldots, e_{n}$ be a basis of $E$, viewed as an $F$-vector space. Then

$$
a=f_{1}\left(e_{1} \otimes 1\right)+\cdots+f_{n}\left(e_{n} \otimes 1\right)
$$

for some $f_{1}, \ldots, f_{n} \in F^{(p)}$, and we can take $F^{\prime}=F\left(f_{1}, \ldots, f_{n}\right)$.
(b) Let $e_{1}, \ldots, e_{n}$ be a basis for $E$ over $F$. Then $e_{1} \otimes 1, \ldots, e_{n} \otimes 1$ forms a basis for $E \otimes_{F} F^{\prime}$ over $F^{\prime}$ as well as for $E \otimes_{F} F^{(p)}$ over $F^{(p)}$. Working in these bases we then have $1, a, \ldots, a^{n-1}$ are linearly dependent over $F^{\prime}$ if and only if $1, \phi(a), \ldots, \phi(a)^{n-1}$ are linearly dependent over $F^{(p)}$.

Lemma 2.1.2. The condition (1) of Theorem 1.3 is equivalent to the following:
(1') For every field $F$ containing $F_{0}$ and every $n$-dimensional étale algebra $E / F$, there exists a finite field extension $F^{\prime} / F$ of degree prime to $p$ and an element $0 \neq a \in E^{\prime}:=E \otimes_{F} F^{\prime}$ such that $\operatorname{tr}_{E^{\prime} / F^{\prime}}(a)=\operatorname{tr}_{E^{\prime} / F^{\prime}}\left(a^{p}\right)=0$.

Proof. For $\left(1^{\prime}\right) \Longrightarrow(1)$ the result follows immediately, as $F^{\prime}=F$ is the only finite field extension of a $p$-field $F$ such that $\left[F^{\prime}: F\right]$ is prime to $p$.

For the other direction, apply (1) to $E^{\prime \prime}:=E \otimes_{F} F^{(p)}$ over $F^{(p)}$ to obtain an element $0 \neq a \in E^{\prime \prime}$ satisfying $\operatorname{tr}_{E^{\prime \prime} / F^{(p)}}(a)=\operatorname{tr}_{E^{\prime \prime} / F^{(p)}}\left(a^{p}\right)=0$. Then by Lemma 2.1.1 there exists a subfield $F \subseteq F^{\prime} \subseteq$ $F^{(p)}$ such that $\left[F^{\prime}: F\right]$ is prime to $p$ and $a$ lies in the natural map

$$
E^{\prime}:=E \otimes_{F} F^{\prime} \hookrightarrow E^{\prime \prime}=E \otimes_{F} F^{(p)}
$$

This $a$ as an element of $E^{\prime}$ then satisfies $\operatorname{tr}_{E^{\prime} / F^{\prime}}(a)=\operatorname{tr}_{E^{\prime} / F^{\prime}}\left(a^{p}\right)=0$.

The implication $(1) \Longrightarrow(2)$ of Theorem 1.3 follows from Lemma 2.1.2 by setting $F=F_{n}$ and $E=E_{n}$.

### 2.2 Proof of Theorem 1.3: $(2) \Longrightarrow$ (3)

Let $Z$ be an $F_{0}$-variety whose $F_{0}$-algebra of rational functions is isomorphic to $L^{\prime}:=L_{n} \otimes_{F_{n}} F^{\prime}$. Define an $\mathrm{S}_{n}$-action on $Z$ to be the action induced by the $\mathrm{S}_{n}$-action on $L^{\prime}$, so that $\sigma(f)(z)=f(\sigma(z))$ for all $z \in Z, f \in L^{\prime}$. The $\mathrm{S}_{n}$-equivariant inclusion

$$
L_{n} \hookrightarrow L^{\prime}=L_{n} \otimes_{F_{n}} F^{\prime}
$$

then gives rise to a dominant $S_{n}$-equivariant rational map $Z \rightarrow \mathbb{A}^{n}$ of degree $d:=\left[F^{\prime}: F_{n}\right]$.
Let $\sigma_{i} \in \mathrm{~S}_{n}$ be such that $\sigma_{i}(1)=i$. These represent the $n$ cosets of $\mathrm{S}_{n-1}$ in $\mathrm{S}_{n}$. As $a \in E^{\prime}=$ $\left(L^{\prime}\right)^{S_{n-1}}, \sigma_{i}(a)$ are the conjugates of $a$ in $L^{\prime}$ and depend only on the coset $h_{i} S_{n-1}$ and not on the particular choice of $h_{i}$ in this coset. Now we define the rational map

$$
\begin{aligned}
f_{a}: & Z \longrightarrow \mathbb{P}^{n-1} \\
& z \mapsto\left(\sigma_{1}(a)(z): \cdots: \sigma_{n}(a)(z)\right) .
\end{aligned}
$$

This map is well-defined, as $\sigma_{1}(a)=a \neq 0$.
Let $\mathrm{S}_{n}$ act on $\mathbb{P}^{n-1}$ in the natural way. For any $\tau \in \mathrm{S}_{n}$ and $z \in Z$ we have

$$
\begin{aligned}
\tau\left(f_{a}(z)\right) & =\tau\left(\sigma_{1}(a)(z): \cdots: \sigma_{n}(a)(z)\right) \\
& =\tau\left(a\left(\sigma_{1}(z)\right): \cdots: a\left(\sigma_{n}(z)\right)\right) \\
& =\left(a\left(\sigma_{\tau^{-1}(1)}(z)\right): \cdots: a\left(\sigma_{\tau^{-1}(n)}(z)\right)\right) \\
& =f_{a}(\tau(z))
\end{aligned}
$$

and so $f_{a}$ is $\mathrm{S}_{n}$-equivariant.
Let $X_{n, p} \subseteq \mathbb{P}^{n-1}$ be the variety defined by

$$
\begin{aligned}
& y_{1}+\cdots+y_{n}=0 \\
& y_{1}^{p}+\cdots+y_{n}^{p}=0 .
\end{aligned}
$$

As $\operatorname{tr}_{E^{\prime} / F^{\prime}}(a)=\sigma_{1}(a)+\cdots+\sigma_{n}(a)=0$ and $\operatorname{tr}_{E^{\prime} / F^{\prime}}\left(a^{p}\right)=0$, the image of $f_{a}$ lies in $X_{n, p}$. In summary, we have the following diagram of $S_{n}$-equivariant rational maps:


Let $H=(\mathbb{Z} / p \mathbb{Z})^{k_{1}} \times \cdots \times(\mathbb{Z} / p \mathbb{Z})^{k_{m}}$, viewed as a subgroup of $\mathrm{S}_{n}$ by embedding into $\mathrm{S}_{p^{k_{1}}} \times \cdots \times$ $\mathrm{S}_{p^{k} m}$ via regular representations on each factor.

We note that $\mathbb{A}^{n}$ contains a smooth $H$-fixed point (e.g. $(0, \ldots, 0)$ ) and apply the going up theorem of J. Kollár and E. Szabó [12, Proposition A.2]. The assumptions of the theorem are satisfied, being that $X$ is complete, $f_{a}$ is rational and that every linear representation of $H$ has a 1-dimensional invariant subspace (as $H$ is abelian and $F_{0}$ contains all $p^{\text {th }}$ roots of unity). Hence $X_{n, p}$ contains a $H$-fixed point defined over $F_{0}\left(\zeta_{p}\right)$.

Let $z \in \mathbb{P}^{n-1}$ be a $H$-fixed point. Then $z$ corresponds to a 1-dimensional subspace of $\mathbb{A}^{n}$ fixed by $H$. As $H \subseteq \mathrm{~S}_{p^{k_{1}}} \times \cdots \times \mathrm{S}_{p^{k_{m}}}$, its representation on $\mathbb{A}^{n}=F_{0}[H]$ decomposes as $F_{0}[H]=W_{1} \oplus$ $\cdots \oplus W_{m}$, where $W_{i}=\mathbb{A}^{p_{i}}$ and we have $H$ acts invariantly on each $W_{i}$.

The action of $H$ on each $W_{i}$ is given by the regular representation of $(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$. As $(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ is abelian, its action on $W_{i}$ is decomposable as a direct sum of irreducible representations of dimension 1 so that $W_{i}=V_{1} \oplus \cdots \oplus V_{p^{k_{i}}}$ with $(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ acting invariantly on each $V_{j}$.

Denote an element of $(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ by $\left(\gamma_{1}, \ldots, \gamma_{k_{i}}\right)$. Then the characters of $(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ can be described as

$$
\begin{aligned}
\chi_{\gamma_{1} \ldots \gamma_{k_{i}}} & :(\mathbb{Z} / p \mathbb{Z})^{k_{i}} \\
\left(\beta_{1}, \ldots, \beta_{k_{i}}\right) & \mapsto \zeta_{p}^{\gamma_{1} \beta_{1}} \ldots \zeta_{p}^{\gamma_{k_{i}} \beta_{k_{i}}}
\end{aligned}
$$

where we have $\gamma_{j}, \beta_{j} \in\{0, \ldots, p-1\}$. Then the irreducible subrepresentation of $W_{i}$ corresponding to $\left(\gamma_{1}, \ldots \gamma_{k}\right)$ is spanned by $\left(\chi_{\gamma_{1} \ldots \gamma_{i}}\left(g_{1}\right), \ldots, \chi_{\gamma_{1} \ldots \gamma_{k_{i}}}\left(g_{p^{k_{i}}}\right)\right)$, where $g_{j}$ are the elements of $(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$. In particular, the coordinates are all $p^{\text {th }}$ roots of unity.

Let $V_{h}$ be the subspace spanned by $\left(\chi_{h}\left(g_{1}\right), \ldots, \chi_{h}\left(g_{p^{k_{i}}}\right)\right)$. We check that these $V_{h}$ are the only 1-dimensional invariant subspaces of $W_{i}$. Let $w \in W_{i}$ and write $w=\sum_{h \in(\mathbb{Z} / p \mathbb{Z})^{k_{i}}} \nu_{h}$, with $v_{h} \in V_{h}$ for each $h$. For each $g \in(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ we have $g \cdot\left(\sum v_{h}\right)=\sum \chi_{h}\left(g^{-1}\right) v_{h}$. If $w$ lies in a 1-dimensional subspace of $W_{i}$ and $v_{h}, v_{h}^{\prime}$ are nonzero, we must have $\chi_{h}(g)=\chi_{h^{\prime}}(g)$ for all $g \in(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ and so $h=h^{\prime}$ and $w \in V_{h}$. Thus the 1-dimensional invariant subspaces of $W_{i}$ are exactly these $V_{h}$.

Write $z=\left(z_{1}, \ldots, z_{m}\right) \in W_{1} \oplus \cdots \oplus W_{m} \cong \mathbb{A}^{n}$. Then for each $i, j$ such that $z_{i}, z_{j} \neq 0$, there exist $h_{i} \in(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ and $h_{j} \in(\mathbb{Z} / p \mathbb{Z})^{k_{j}}$ corresponding to characters $\chi_{i}$ of $W_{i}$ and $\chi_{j}$ of $W_{j}$ such that
$\chi_{i}\left(g_{i}\right)=\chi_{j}\left(g_{j}\right)$ for all $g_{i} \in(\mathbb{Z} / p \mathbb{Z})^{k_{i}}$ and $g_{j} \in(\mathbb{Z} / p \mathbb{Z})^{k_{j}}$. If $z_{i} \neq 0$ for at least 2 different $i$, then this requires $\chi_{i}$ being the trivial character for all $i$ such that $z_{i} \neq 0$. Otherwise $z_{i}=0$ for all $i \neq j$ for a particular $j$.

Fix $j \in\{1, \ldots, m\}$ and suppose $z_{i}=0$ for all $i \neq j$. Then the coordinates of $z_{j}$ are all $\lambda \zeta_{p}^{\ell}$ for some $p^{\text {th }}$ roots of unity $\zeta_{p}^{\ell}$ and for some fixed scalar $\lambda \neq 0$. As $z \in X$ we thus have $\left(\lambda \zeta_{p}^{\ell_{1}}\right)^{p}+\cdots+$ $\left(\lambda \zeta_{p}{ }^{{ }_{p}{ }^{k_{i}}}\right)^{p}=\lambda^{p} p^{k_{i}}=0$, which is impossible in characteristic $\neq p$.

Thus $z$ is acted on by $H$ by the trivial character. Writing $z=\left(z_{1}: \cdots: z_{n}\right)$ as an element of $\mathbb{P}^{n-1}$ we then have $z_{i}=z_{j}$ when $i$ and $j$ lie in the same $H$-orbit and so there are at most $m$ different coordinates of $z$. Using these coordinates, and as $z \in X$, we have thus obtained a solution $y=\left(y_{1}: \cdots: y_{n}\right) \in \mathbb{P}^{m-1}\left(F_{0}\right)$ to the system of equations

$$
\begin{aligned}
& p^{k_{1}} y_{1}+\cdots+p^{k_{m}} y_{m}=0 \\
& p^{k_{1}} y_{1}^{p}+\cdots+p^{k_{m}} y_{m}^{p}=0
\end{aligned}
$$

as required.

### 2.3 Proof of Theorem 1.3: $(3) \Longrightarrow$ (1)

For the implication $(3) \Longrightarrow$ (1) it suffices to prove the following:
Lemma 2.3.1. Let $p$ be a prime and $n \geq 1$ be a positive integer with base $p$ expansion $n=p^{k_{1}}+$ $\cdots+p^{k_{m}}$. Suppose $F$ is a $p$-field and $E$ is an étale algebra of degree $n$ over $F$. Then we can write $E \cong E_{1} \times \cdots \times E_{m}$, where each $E_{i}$ is an étale algebra of degree $p^{k_{i}}$ over $F$.

Indeed, suppose this lemma holds, and let $y=\left(y_{1}: \cdots: y_{m}\right) \in \mathbb{P}^{m-1}\left(F_{0}\right)$ be a solution to (1.1) as in condition (3). Let $F$ be a $p$-closed field and $E / F$ be an $n$-dimensional étale algebra. From the lemma, we can write $E \cong E_{1} \times \cdots \times E_{m}$ with each $E_{i}$ an étale algebra of degree $p^{k_{i}}$ over $F$. We view $a:=\left(y_{1} 1_{E_{1}}, \ldots, y_{m} 1_{E_{m}}\right)$ as a nonzero element of $E_{1} \times \cdots \times E_{m}$. As each $y_{i}$ lies in $F_{0} \subseteq F$ we have $\operatorname{tr}_{E_{i} / F}\left(y_{i}\right)=p^{k_{i}} y_{i}$ and $\operatorname{tr}_{E_{i} / F}\left(y_{i}^{p}\right)=p^{k_{i}} y_{i}^{p}$. Then we have

$$
\begin{aligned}
\operatorname{tr}_{E / F}(a) & =\operatorname{tr}_{E_{1} / F}\left(y_{i}\right)+\cdots+\operatorname{tr}_{E_{m} / F}\left(y_{i}\right)=p^{k_{1}} y_{1}+\cdots+p^{k_{m}} y_{m}=0 \\
\operatorname{tr}_{E / F}\left(a^{p}\right) & =\operatorname{tr}_{E_{1} / F}\left(y_{i}^{p}\right)+\cdots+\operatorname{tr}_{E_{m} / F}\left(y_{i}^{p}\right)=p^{k_{1}} y_{1}^{p}+\cdots+p^{k_{m}} y_{m}^{p}=0
\end{aligned}
$$

and so

$$
\begin{equation*}
a=\left(y_{1} 1_{E_{1}}, \ldots, y_{m} 1_{E_{m}}\right) \tag{2.3.1}
\end{equation*}
$$

gives us a solution to (1).
We now consider Lemma 2.3.1. Let $F$ be a $p$-field and $E / F$ be an étale algebra of degree $n$ over $F$. As the only finite field extensions of $F$ have degree a power of $p$, we can write $E \cong F_{1} \times \cdots \times F_{t}$ with each $F_{i}$ being a field extension of $F$ of degree a power of $p$, and the sum of the degrees being $n$.

To prove the lemma, it suffices to arrange these fields into $m$ distinct groups $G_{1}, \ldots, G_{m}$, such that $E_{i}:=\prod_{F_{j} \in G_{i}} F_{j}$ is an étale algebra of degree $p^{k_{i}}$ over $F$, for each $i$. This is satisfied if and only if $\sum_{F_{j} \in G_{i}}\left[F_{j}: F\right]=p^{k_{i}}$. Thus, it suffices to rearrange the integers $\left[F_{j}: F\right]$, which partition $n$, into distinct groups such that the sum of the integers in each group $i$ is $p^{k_{i}}$.

A partition $\lambda=\left[\lambda_{1}, \ldots, \lambda_{r}\right]$ of $n$ is an unordered collection of (possibly repeated) positive integers $\lambda_{1}, \ldots, \lambda_{r}$ such that $\lambda_{1}+\cdots+\lambda_{r}=n$. For two partitions $\lambda=\left[\lambda_{1}, \ldots, \lambda_{r}\right]$ and $\mu=\left[\mu_{1}, \ldots, \mu_{s}\right]$ of $n$ we say that $\lambda$ is a refinement of $\mu$, and write $\mu \preceq \lambda$, if $\mu$ can be obtained from $\lambda$ by partitioning each of the numbers $\lambda_{1}, \lambda_{r}$. Equivalently, we have $\mu \preceq \lambda$ if $\mu_{1}, \ldots, \mu_{s}$ can be arranged into $r$ distinct collections such that the sum of the numbers in each group $i$ is $\lambda_{i}$. For example, we have $[3,2,1] \prec[3,3] \prec[6]$ and $[3,2,1] \prec[5,1] \prec[6]$. This operation defines a partial order on the partitions of $n$.

Let $p$ be a prime. We call a partition $\lambda=\left[\lambda_{1}, \ldots, \lambda_{r}\right]$ a $p$-partition of $n$ if each $\lambda_{i}$ is a power of $p$. Every integer $n \geq 1$ has a unique $p$-partition $\left[p^{k_{1}}, \ldots, p^{k_{m}}\right]$ corresponding to the base $p$ expansion $n=p^{k_{1}}+\cdots=p^{k_{m}}$. We denote this partition by $[n]_{p}$. As the partition arising from the decomposition $E \cong F_{1} \times \cdots \times F_{r}$ is a $p$-partition of $n$, we reduce Lemma 2.3.1 to the following:

Lemma 2.3.2. Let $\lambda=\left[\lambda_{1}, \ldots, \lambda_{r}\right]$ be a p-partition of $n$. Then $\lambda \preceq[n]_{p}$.
Let $\mathscr{P}$ be the set of all $p$-partitions of $n$. As there are finitely many $p$-partitions of $n$, to prove this lemma it suffices to show that $[n]_{p}$ is the unique maximal element of $\mathscr{P}$ with respect to the refinement ordering. Again, as $\mathscr{P}$ is finite, we know that there exists at least one maximal partition, so it suffices to show that if $\lambda \neq[n]_{p}$, then $\lambda$ is not maximal.

We assume $\lambda \neq[n]_{p}$. The base $p$ expansion $n=p^{k_{1}}+\ldots p^{k_{m}}$ is unique with the property that each power of $p$ appears in the sum at most $p-1$ times, so the partition $[n]_{p}=\left[p^{k_{1}}, \ldots, p_{k_{m}}\right]$ is unique with the same property. Thus some power of $p$ appears in $\lambda$ at least $p$ times, say $\lambda_{1}=\cdots=$ $\lambda_{p}=p^{\ell}$. Then we have

$$
\lambda=[\underbrace{p^{\ell}, \ldots, p^{\ell}}_{p \text { times }}, \lambda_{p+1}, \ldots, \lambda_{r}] \prec\left[p^{\ell+1}, \lambda_{p+1}, \ldots, \lambda_{r}\right]
$$

with the latter partition again lying in $\mathscr{P}$. Thus $\lambda$ is not maximal and so $[n]_{p}$ is the unique maximal element of $\mathscr{P}$, proving Lemma 2.3.2, and hence the implication $(3) \Longrightarrow(1)$.

### 2.4 Density of $F$-points on hypersurfaces

Recall that a closed subvariety $X$ of projective space is called $a$ cone over a point $c \in X$ if $X$ contains the line through $c$ and $c^{\prime}$ for every $c \neq c^{\prime} \in X$. We will say that $X$ is a cone if it is a cone over one of its points.
Lemma 2.4.1. Let $X$ be a degree $d$ hypersurface in $\mathbb{P}^{\ell}$ defined over a field $F$, which is not a cone over $c \in X$. If $H \simeq \mathbb{P}^{\ell}$ is a hyperplane defined over $F$, then the projection $\pi: X \backslash\{c\} \rightarrow H$ is dominant.

Proof. Without loss of generality, after a change of basis we may assume $c=(1: 0: \cdots: 0)$ and $H$ is defined by $x_{0}=0$. Let $X$ be defined by $f$ of degree $d$. Let $U=\left\{\left(0: a_{1}: \ldots a_{\ell}\right) \mid a_{i} \neq 0\right.$ for each $\left.a_{i}\right\}$, which is dense and open in $H$.

For a contradiction we suppose there exists $a=\left(0: a_{1}: \cdots: a_{\ell}\right) \in U$ such that $a \notin \pi(X \backslash\{c\})$. Then $f\left(x_{0}: a_{1}: \cdots: a_{\ell}\right)$ is nonzero of degree zero as a polynomial in $x$. As each of the $a_{i}$ are nonzero we thus have $x_{0}$ does not appear in $f$. But then for all $b=\left(b_{0}: \cdots: b_{\ell}\right) \in X$ we have

$$
f\left(\lambda_{1}+\lambda_{2} b_{0}: \lambda_{2} b_{1}: \cdots: \lambda_{2} b_{\ell}\right)=f\left(\lambda_{2} b_{0}: \lambda_{2} b_{1}: \cdots: \lambda_{2} b_{\ell}\right)=f(b)=0 .
$$

Hence $X$ contains the line through $c$ and $b$. As $b \in X$ was arbitrary we have a contradiction. Thus $U \subseteq \pi(X \backslash\{c\})$ and hence $\pi$ is dominant.

An effective zero cycle in $X$ is a subvariety (possibly reducible) of $X$ of dimension 0 . The degree of a zero cycle $Z$ in $X$ is defined as $\operatorname{dim}_{F}(F[Z])$. Our main lemma for this section is the following.

Lemma 2.4.2. Let $F$ be a p-closed field of characteristic $\neq p$. Suppose $X \subset \mathbb{P}^{\ell}$ is a closed hypersurface of degree $\leq p$ defined over a $p$-closed field $F$ of characteristic $\neq p$. Assume that $X$ has an $F$-point $c$ such that $X$ is not a cone over $c$. Then $F$-points are dense in $X$.

Proof. Case 1: $X$ is a hypersurface of degree $d<p$. Note that effective zero cycles of degree $d$ are dense in $X$ (these can be obtained by intersecting $X$ with lines defined over $F$ in $\mathbb{P}^{\ell}$ ). Since $F$ is a $p$-closed field, every effective zero cycle of degree $d<p$ splits over $F$, that is, is a sum of $d$ $F$-points. Consequently, $F$-points are dense in $X$.

Case 2: $X$ is reducible over $F$, that is, its irreducible components, $X_{1}, \ldots, X_{r}$ are defined over $F$ and $r \geq 2$. Here each $X_{i}$ is a hypersurface of degree $<p$. By Case 1, $F$-points are dense in each $X_{i}$; hence, they are dense in $X$.

Case 3: $X$ is irreducible over $F$ but reducible over $\bar{F}$. Note that since $F$ is a $p$-closed field, and char $F \neq p, F$ is perfect. Hence, the irreducible components $X_{1}, \ldots, X_{r}$ of $X$ are transitively permuted by the Galois group $\operatorname{Gal}(\bar{F} / F)$, which is a pro- $p$ group. Thus $r \geq 2$ is a power of $p$. Moreover, since $\operatorname{deg}(X)=\operatorname{deg}\left(X_{1}\right)+\cdots+\operatorname{deg}\left(X_{r}\right) \leq p$, we conclude that $r=p$ and $\operatorname{deg}\left(X_{1}\right)=$ $\cdots=\operatorname{deg}\left(X_{p}\right)=1$. In other words, $X$ is a union of the hyperplanes $X_{1}, \ldots, X_{p}$. Now observe that $c \in X(F)$ is fixed by $\operatorname{Gal}(\bar{F} / F)$. After relabeling the components, we may assume that $c \in X_{1}$. Translating $X_{1}$ by $\operatorname{Gal}(\bar{F} / F)$, we see that $c$ lies on every translate of $X_{1}$, that is, on every $X_{i}$ for $i=1, \ldots, r$. Since each $X_{i}$ is a hyperplane, we conclude that $X$ is a cone over $c$, contradicting our assumption.

Case 4: $X$ is absolutely irreducible. Choose a hyperplane $H \simeq \mathbb{P}^{\ell-1}$ in $\mathbb{P}^{\ell}$ such that $H$ is defined over $F$ and $c \notin H$. Let $\pi: X \backslash\{c\} \rightarrow H$ be projection from $c$. Since $X$ is not a cone over $c$, this map is dominant, by Lemma 2.4.1. In particular, there is a dense open subset $U \subset H$ such that $\pi$ is finite over $U$. The preimage $\pi^{-1}(u)$ of any $F$-point $u \in U(F)$ is then an effective zero cycle of degree $\leq p-1$. Once again, every such zero cycle splits over $F$, that is, $\pi^{-1}(u)$ is a union of $F$-points. Taking the union of $\pi^{-1}(u)$, as $u$ varies over $U(F)$, we obtain a dense set of $F$-points in $X$.

If $p=2$ or 3 , then Lemma 2.4.2 remains true for all infinite fields $F$ (not necessarily $p$-closed), under mild additional assumptions. We say $X$ is rational over $F$ if $X$ is birationally isomorphic to some projective space $\mathbb{P}^{k}$ over $F$.

Lemma 2.4.3. (a) Let $X \subset \mathbb{P}^{\ell}$ be a hypersurface of degree 2 defined over a field $F$ of characteristic $\neq 2$. Assume $X$ has an $F$-point and $X$ is not a cone. Then $X$ is rational over $F$. In particular, $F$-points are dense in $X$.
(b) (J. Kollár) Let $X \subset \mathbb{P}^{\ell}$ be an absolutely irreducible cubic hypersurface of dimension $\geq 2$ defined over a field $F$. Assume $X$ has an $F$-point, and $X$ is not a cone. If $X$ is singular, assume also that $\operatorname{char} F \neq 2$ or 3 . Then $X$ is unirational over $F$. In particular, if $F$ is infinite, then $F$-points are dense in $X$.

Proof. (a) Suppose $X$ is given by $q=0$, where $q$ is a quadratic form on $F^{\ell+1}$. If $X$ were reducible, then $q$ would be a product of two linear forms, say $q=l_{1} l_{2}$. Then $X$ would be a cone over the intersection of the lines defined by $l_{1}$ and $l_{2}$. This is not the case, hence $X$ is irreducible. As $q$ is irreducible, $X$ is also smooth. The stereographic projection from a point $c \in X(F)$ to a hyperplane $H \subset \mathbb{P}^{\ell}$ defined over $F$ and not passing through $c$, gives rise to a birational isomorphism between $X$ and $H$.
(b) If $X$ is smooth, see [7, Theorem 1.1]. If $X$ is singular and $F$ is perfect, see [7, Theorem 1.2]. Finally, if $X$ is singular and $F$ is an imperfect field of characteristic $\neq 2,3$, see the remark after the statement of Theorem 1.2 on [7, p. 468].

### 2.5 Geometry of the hypersurface $X_{n, p}$

In this section we will prove some geometric properties of the hypersurface

$$
X_{n, p}:=\left\{\left(x_{1}: \cdots: x_{n}\right) \mid x_{1}+\cdots+x_{n}=x_{1}^{p}+\cdots+x_{n}^{p}=0\right\} \subset \mathbb{P}^{n-1}
$$

defined over the base field $F_{0}$. These will be used in the next section to prove assertion $(*)$.
Let $\Delta_{n}$ be the union of the "diagonal" hyperplanes $x_{i}=x_{j}$, over all $1 \leq i<j \leq n$.
Lemma 2.5.1. Assume char $F_{0} \neq 2,3$ or $p$. Then
(a) The singular locus of $X_{n, p}$ is $X_{n, p} \cap\left\{\left(x_{1}: \cdots: x_{n}\right) \mid x_{1}^{p-1}=\cdots=x_{n}^{p-1}\right\}$.
(b) $X_{n, p}$ is absolutely irreducible if $n \geq 5$.
(c) $X_{n, p}$ is not contained in $\Delta_{n}$ if $n \geq 3$.
(d) Let $(1: \cdots: 1) \neq c \in X_{n, p}$. Then $X_{n, p}$ is not a cone over $c$.
(e) $X_{n, p}$ is not a cone if $p>2$.

Proof. In order to prove the lemma we may, without loss of generality, pass to the algebraic closure of $F_{0}$, that is, assume that $F_{0}$ is algebraically closed.
(a) follows from the Jacobian criterion.
(b) Assume the contrary: $X_{n, p}$ has at least two irreducible components, $X_{1}$ and $X_{2}$. Since $X_{n, p}$ is a hypersurface in $\mathbb{P}^{n-1}, \operatorname{dim}\left(X_{1}\right)=\operatorname{dim}\left(X_{2}\right)=n-3$, and $\operatorname{dim}\left(X_{1} \cap X_{2}\right)=n-4$ with $X_{1} \cap X_{2}$ contained in the singular locus of $X$. By part (a), the singular locus of $X$ is 0 -dimensional. Thus $n-4 \leq 0$, a contradiction, and the result follows.
(c) First assume $n \geq 5$ and assume the contrary. We have $X_{n, p}$ is irreducible by part (b) and hence lies in one the hyperplanes defined by $x_{i}=x_{j}$. Since $X_{n, p}$ is invariant under the action of $\mathrm{S}_{n}$, it is contained in every hyperplane $H_{i j}$. That is,

$$
X_{n, p} \subset \bigcap_{1 \leq i<j \leq n} H_{i j}=\{(1: \cdots: 1)\},
$$

contradicting $\operatorname{dim}\left(X_{n, p}\right)=n-3 \geq 2$.
In the remaining cases, where $n=3$ or 4 , we will exhibit a point $y \in X_{n, p}$ which does not lie in $\Delta_{n}$.

If $n=4$ we can take $y:=\left(1: \zeta_{4}: \zeta_{4}^{2}: \zeta_{4}^{3}\right)$, where $\zeta_{4} \in F_{0}$ is a primitive 4 th root of unity (recall that we are assuming $F_{0}$ to be algebraically closed and char $F_{0} \neq 2$ ).

Now suppose $n=3$. If $p \neq 2$, then we can take $y:=(1:-1: 0)$. Otherwise we can take $y:=\left(1: \zeta_{3}: \zeta_{3}^{2}\right)$. This completes the proof of (c).
(d) Assume the contrary. Since $\mathrm{S}_{n}$ acts on $X_{n, p} \subset \mathbb{P}^{n-1}$ by permuting coordinates, $X_{n, p}$ is a cone over $g \cdot c$ for every $g \in \mathrm{~S}_{n}$. From this it follows that $X_{n, p}$ contains the linear span of $\left\{g \cdot c \mid g \in \mathrm{~S}_{n}\right\}$. Denote this linear span by $L$. Let $H \simeq \mathbb{P}^{n-2}$ be defined by

$$
H:=\left\{\left(x_{1}: \cdots: x_{n}\right) \mid x_{1}+\cdots+x_{n}=0\right\} .
$$

Then $L$ is an $\mathrm{S}_{n}$-invariant linear subspace of $H$. If $\operatorname{char} F_{0}$ does not divide $n$, then the $\mathrm{S}_{n}$-representation on $H$ is irreducible. Hence, the only $\mathrm{S}_{n}$-invariant subspace of $H$ is $H$ itself. Thus $H=L \subset X_{n, p}$, a contradiction. If char $F_{0}$ divides $n$, then the only other possibility is $L=\{(1: 1 \cdots: 1)\}$. This is ruled out by our assumption that $c \neq(1: \cdots: 1)$.
(e) Assume the contrary, that $X_{n, p}$ is a cone over some point $c \in X_{n, p}$. By part (d), $c=(1: \cdots: 1)$. Note here that this implies char $F_{0}$ divides $n$. Whenever $X_{n, p}$ contains a point $y=\left(y_{1}: \cdots: y_{n}\right)$, it contains the entire line through $c$ and $y$. That is,

$$
\left(1+t y_{1}\right)^{p}+\cdots+\left(1+t y_{n}\right)^{p}=0
$$

as a polynomial in $t$. In particular, $p\left(y_{1}^{p-1}+\cdots+y_{n}^{p-1}\right)$, which is the coefficient of $t^{p-1}$ in this polynomial, should vanish for every $y=\left(y_{1}: \cdots: y_{n}\right) \in X_{n, p}$. Set $y:=(-1: 1: 0: \cdots: 0)$. Note that $y \in X_{n, p}$, because $p>2$. For this $y, p\left(y_{1}^{p-1}+\cdots+y_{n}^{p-1}\right)=0$ reduces to $2 p=0$, contradicting our assumptions that char $F_{0} \neq 2$ or $p$.

### 2.6 Proof of Theorem 1.3: (*)

In this section we complete the proof of Theorem 1.3 by proving assertion (*).
Given an étale algebra $E / F$ of degree $n$, we define $X_{E / F, p}$ as the degree $p$ hypersurface in $\mathbb{P}(E)=\mathbb{P}_{F}^{n-1}$, given by $\operatorname{tr}_{E / F}(x)=\operatorname{tr}_{E / F}\left(x^{p}\right)=0$. Let $\Delta_{E / F}$ be the discriminant locus in $\mathbb{P}(E)$, that is, the closed subvariety of $\mathbb{P}(E)$ defined by the condition $1, a, \ldots, a^{n-1}$ are linearly dependent over $F$ (here $a \in E$ ).

Lemma 2.6.1. $X_{E / F, p}$ and $\Delta_{E / F}$ are $F$-forms of the varieties $X_{n, p}$ and $\Delta_{n}$ defined above.
Proof. It suffices to show that $X_{E / F, p} \cong X_{n, p}$ and $\Delta_{E / F} \cong \Delta_{n}$ when we pass to the algebraic closure of $F$. That is, it suffices to show that the varieties are isomorphic when we assume $F$ to be algebraically closed.

When $F$ is algebraically closed, there exist no finite field extensions and hence $E \cong F^{n}$. Then we have $\operatorname{tr}_{E / F}\left(x_{1}, \ldots, x_{n}\right)=x_{1}+\cdots+x_{n}$ and $\operatorname{tr}_{E / F}\left(x_{1}^{p}, \ldots, x_{n}^{p}\right)=x_{1}^{p}+\cdots+x_{n}^{p}$ and hence $X_{E / F, p} \cong$ $X_{n, p}$. In the standard basis for $F^{n}$ over $F$, the discriminant locus becomes the determinant of the Vandermonde matrix of an element $\left(a_{1}, \ldots, a_{n}\right)$. This is nonzero if and only if the coordinates are distinct. That is, we have $\Delta_{E / F}$ is isomorphic to $\Delta_{n}$, being the union of the hyperplanes $x_{i}=x_{j}$.

We now prove the assertion $(*)$ of Theorem 1.3, which we do by appealing to Lemma 2.4.2. Assume (3) holds. Our goal is to show that $a$ can be chosen so that $(\dagger) E=F[a]$ in part (1), and $(\dagger \dagger) E^{\prime}=F^{\prime}[a]$ in part (2). We will see that it suffices to prove ( $\dagger$ ) in part (1).

Suppose ( $\dagger$ ) holds in part (1), and in particular that $a$ generates $E_{n} \otimes_{F_{n}} F_{n}^{(p)}$ over $F_{n}^{(p)}$. Then by Lemma 2.1.1(a), $a$ lies in the image of the natural map $\phi: E \otimes_{F} F^{\prime} \rightarrow E \otimes_{F} F^{(p)}$ for some field extension $F^{\prime} / F_{n}$ such that $\left[F^{\prime}: F\right]$ is prime to $p$. By Lemma 2.1.1(b) we have $\phi^{-1}(a)$ generates $E_{n} \otimes_{F_{n}} F^{\prime}$ over $F^{\prime}$. Thus ( $\dagger \dagger$ ) follows from ( $\dagger$ ) and so it suffices to prove ( $\dagger$ ) for part (1), which can be restated as follows:
If $X_{E / F, p}$ has an $F$-point, then $X_{E / F, p}$ has an $F$-point away from $\Delta_{E / F}$.
As $X_{E / F, p}$ is an $F$-form of $X_{n, p}$ and $n \geq 3$, Lemma 2.5.1(c) tells us that $X_{E / F, p}$ is not contained in $\Delta_{E / F}$. Assuming that condition (3) of Theorem 1.3 holds, we have constructed an $F$-point $a$ of $X_{E / F, p}(F)$, as in (2.3.1). Thus it suffices to prove that $F$-points are dense in $X_{E / F, p}$.

We claim that under the assumptions of Theorem $1.3(*), X_{E / F, p}$ is not a cone over $a$.
Indeed, if condition (i) of Theorem $1.3(*)$ holds, that is, $a=\left(y_{1}: \ldots: y_{m}\right) \neq(1: \cdots: 1)$, then formula (2.3.1) tells us that $a$ is not a scalar (that is, $a \notin F \cdot 1_{E}$ ). If $X_{E / F, p}$ were a cone over $a$, then it
would remain a cone over $a$ after passing to the algebraic closure $\bar{F}$ of $F$. When we pass from $F$ to $\bar{F}, E$ becomes split, and so isomorphic to $F^{n}$; so $\mathbb{P}(E)$ reduces to $\mathbb{P}\left(F^{n}\right) \simeq \mathbb{P}^{n-1}, X_{E / F, p}$ reduces to $X_{n, p}$, and the condition that $a$ is not a scalar in $E$ translates into $a \neq(1: \cdots: 1)$ in $\mathbb{P}^{n-1}$. By Lemma 2.5.1(d), $X_{n, p}$ cannot be a cone over $a$, a contradiction. We conclude that $X_{E / F, p}$ is not a cone over $a$, as claimed.

On the other hand, if condition (ii) of Theorem 1.3(*) holds, that is if $p>2$ then by Lemma 2.5.1(e), $X_{n, p}$ is not a cone over any of its points, and hence, neither is $X_{E / F, p}$. This proves the claim.

Thus in order to prove assertion $(*)$, it suffices to establish the following lemma.
Lemma 2.6.2. Suppose $p$ is a prime, $F$ is a p-closed field, and $E / F$ is an étale algebra of degree $n \geq$ 3. If $X_{E / F, p}$ has an $F$-point, and $X_{E / F, p}$ is not a cone over this point, then $F$-points are dense in $X_{E / F, p}$.

Lemma 2.6.2 is a special case of Lemma 2.4.2. This completes the proof of assertion $(*)$ of Theorem 1.3, and the proof of the theorem.

## Chapter 3

## When are there solutions to (1.1)?

Before exploring the existence of solutions to (1.1) in the cases of $p=2$ and $p=3$ in more detail we will prove some more general results.

Lemma 3.1. Let $F_{0}$ be a field such that $\operatorname{char} F_{0} \neq p$. Then the system (1.1) has solution in $\mathbb{P}^{m-1}\left(F_{0}\right)$ if any of the following conditions hold:
(1) $k_{i} \equiv k_{j}(\bmod p)$ and $k_{i^{\prime}} \equiv k_{j^{\prime}}(\bmod p)$ for distinct $i, i^{\prime}, j, j^{\prime}$, and either $p$ is odd, or $\sqrt{-1} \in F_{0}$.
(2) $m \geq p+3$ and either $p$ is odd or $\sqrt{-1} \in F_{0}$.
(3) $\operatorname{char} F_{0} \neq 0$ and $m \geq p+2$

Proof. (1) Let $y=\left(y_{1}: \cdots: y_{m}\right)$, where $y_{i}=1, y_{j}=\sqrt[p]{-p^{k_{i}-k_{j}}}$ and $y_{h}=0$ for every other $h \neq i, j$ and $y^{\prime}=\left(y_{1}^{\prime}: \ldots y_{j}^{\prime}\right)$, where $y_{i^{\prime}}=1, y_{j^{\prime}}=\sqrt[p]{-p^{k_{i^{\prime}}-k_{j^{\prime}}}}$ and $y_{h^{\prime}}=0$ for every other $h^{\prime} \neq i^{\prime}, j^{\prime}$. Then $y$ and $y^{\prime}$ lie in $\mathbb{P}^{m-1}\left(F_{0}\right)$, as $\sqrt[p]{-1} \in F_{0}$ and $k_{i}-k_{j} \equiv k_{i^{\prime}}-k_{j^{\prime}} \equiv 0(\bmod p)$. These points, and also the line between them, lie in the hypersurface defined by $p^{k_{1}} y_{1}^{p}+\cdots+y^{k_{m}} y_{m}^{p}$. Intersecting this line with the hyperplane defined by $p^{k_{1}} y_{1}+\cdots+p^{k_{m}} y_{m}$, we obtain a solution to (1.1).
(2) This follows from (1) by two applications of the pigeonhole principle.
(3) Let $q=\operatorname{char} F_{0}>0$. Then $F_{0}$ contains a copy of $\mathbb{Z} / q \mathbb{Z}$. By Chevalley's theorem, finite fields are $C_{1}$ (see [10, Theorem 5.2.1]), which implies that any quadratic form of dimension $m-1$ is isotropic when $m-1>2$, that is, when $m \geq 4$. Thus (1.3) has a nontrivial solution in $\mathbb{Z} / q \mathbb{Z} \subseteq F_{0}$ when $m \geq 4$.

Lemma 3.2. Let $m=2$ and $F_{0}$ be a field such that char $F_{0} \neq p$. The system (1.1) has a solution in $\mathbb{P}^{1}\left(F_{0}\right)$ if and only if $p^{\left(k_{1}-k_{2}\right)(p-1)}+(-1)^{p}=0$ in $F_{0}$.

Proof. A solution to the system when $m=2$ is equivalent to a nonzero solution in $F_{0}$ to $p^{k_{1}-k_{2}} y_{1}^{p}+$ $\left(-p^{k_{1}-k_{2}} y_{1}\right)^{p}=0$. This occurs precisely when $p^{k_{1}-k_{2}}\left(1+(-p)^{\left(k_{1}-k_{2}\right)(p-1)}\right)=0$ in $F_{0}$.

## Chapter 4

## The Hermite-Joubert problem for $p=2$

For $p=2$ we can strengthen Theorem 1.3 to the following:
Theorem 4.1. Let $F_{0}$ be a field of characteristic $\neq 2$ and $n=2^{k_{1}}+\cdots+2^{k_{m}} \geq 3$ where $k_{1}>\cdots>$ $k_{m} \geq 0$.
(a) Conditions (1), (2) and (3) of Theorem 1.3 (with $p=2$ ) are equivalent to:
(4) For every field $F$ containing $F_{0}$ and every $n$-dimensional étale algebra $E / F$, there exists an element $0 \neq a \in E$ such that $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{2}\right)=0$.
(b) Moreover, if (4) holds, char $F_{0}$ does not divide $n$, and $F$ is an infinite field, then the element $a \in E$ in (4) can be chosen so that $E=F[a]$.

Proof. To prove (a), by applying Lemma 2.1.2 it suffices to prove the equivalence of conditions ( $1^{\prime}$ ) and (4).

Let $F^{\prime}$ be a field of finite odd degree over $F$. By a theorem of Springer, a quadratic form defined over a field $F$ of characteristic $\neq 2$ has a solution in $F$ if and only if it has a solution in $F^{\prime}$, see [ 9 , Theorem VII.2.7].

Let $E^{\prime}=E \otimes_{F} F^{\prime}$. The equations $\operatorname{tr}_{E^{\prime} / F^{\prime}}(x)=\operatorname{tr}_{E^{\prime} / F^{\prime}}\left(x^{2}\right)=0$ and $\operatorname{tr}_{E / F}(x)=\operatorname{tr}_{E / F}\left(x^{2}\right)=0$ define the same quadratic form over $F^{\prime}$ and $F$ respectively. This and Springer's theorem give ( $1^{\prime}$ ) is equivalent to (4)

To prove (b), we wish to show that $F$-points are dense in $X_{E / F, 2}$, by applying Lemma 2.4.3. As we are assuming char $F_{0} \neq 2$, we have $X_{n, 2}$ is a smooth hypersurface of degree 2, by Lemma 2.5.1(a). As $X_{E / F, 2}$ is an $F$-form of $X_{n, 2}$ we then have $X_{E / F, 2}$ is a smooth hypersurface of degree 2, and hence not a cone. As we are assuming (4) holds, $X_{E / F, 2}$ has an $F$-point and hence by Lemma 2.4 .3 we have $F$-points are dense in $X_{E / F, 2}$.

By Lemma 2.5.1(c) we have $X_{n, 2}$ is not contained in $\Delta_{n}$ and hence $X_{E / F, 2}$ is not contained in $\Delta_{E / F}$. Thus $X_{E / F, 2}$ contains an $F$-point away from $\Delta_{E / F}$. This point generates $E$ over $F$.

In the the rest of the chapter we will describe and prove necessary and sufficient conditions on the field $F_{0}$ to satisfy condition (4) of Theorem 4.1 for a given $n$, by appealing to condition (3) of Theorem 1.3. We provide a summary here, and go in to more detail for particular cases of $m$ in the following subsections.

Theorem 4.2. Let $k_{1}>\cdots>k_{m} \geq 0$ be integers and $F_{0}$ be a field of characteristic $\neq 2$. Consider the system of equations:

$$
\begin{align*}
& 2^{k_{1}} y_{1}+\cdots+2^{k_{m}} y_{m}=0 \\
& 2^{k_{1}} y_{1}^{2}+\cdots+2^{k_{m}} y_{m}^{2}=0 . \tag{4.1}
\end{align*}
$$

The following statements hold:

- If $m=1$ then (4.1) has no solution in $\mathbb{P}^{0}\left(F_{0}\right)$.
- If $m=2$ then (4.1) has a solution in $P^{1}\left(F_{0}\right)$ if and only if $\operatorname{char} F_{0} \mid 2^{k_{1}-k_{2}}+1$.
- If $m=3$ then (4.1) has a solution in $P^{2}\left(F_{0}\right)$ if and only if $-2^{k_{1}+k_{2}+k_{3}} n$ is a square in $F_{0}$.
- If $m \geq 4, \sqrt{-1} \in F_{0}$ and $\sqrt{2} \in F_{0}$, then (4.1) has a solution in $P^{m-1}\left(F_{0}\right)$.
- If $m \geq 4$ and char $F_{0} \neq 0$, then (4.1) has a solution in $P^{m-1}\left(F_{0}\right)$.
- If $m \geq 5$ and $\sqrt{-1} \in F_{0}$, then (4.1) has a solution in $P^{m-1}\left(F_{0}\right)$.
- If $m \geq 6$ and $F_{0}$ is a nonreal number field, then (4.1) has a solution in $P^{m-1}\left(F_{0}\right)$.
- In all cases of $m$, it is necessary for $F_{0}$ to be nonreal for (4.1) to have a solution in $P^{m-1}\left(F_{0}\right)$.

In the statements throughout the rest of this chapter we assume $k_{1}>k_{2}>\cdots>k_{m} \geq 0$ and char $F_{0} \neq 2$. The subsections are sorted by $m$, being the number of digits in the base 2 expansion of $n$.

## $4.1 m \leq 2$

## Proposition 4.1.1.

If $m=1$, then (4.1) has no solution in $\mathbb{P}^{0}\left(F_{0}\right)$.
If $m=2$, then (4.1) has a solution in $\mathbb{P}^{1}\left(F_{0}\right)$ if and only if $\operatorname{char} F_{0} \mid 2^{k_{1}-k_{2}}+1$.
Proof. If $m=1$ then $\mathbb{P}^{0}\left(F_{0}\right)$ consists of a point, which is not a solution to (4.1). If $m=2$, the result follows from Lemma 3.2.

## $4.2 m=3$

Proposition 4.2.1. Suppose $m=3$. Then (4.1) has a solution in $\mathbb{P}^{2}\left(F_{0}\right)$ if and only if $-2^{k_{1}+k_{2}+k_{3}} n$ is a square in $F_{0}$.

Proof. Let $\ell_{1}=k_{1}-k_{3}$ and $\ell_{2}=k_{2}-k_{3}$. Then (4.1) has a solution in $\mathbb{P}^{2}\left(F_{0}\right)$ if and only if the quadratic form

$$
Q\left(y_{1}, y_{2}\right)=2^{\ell_{1}} y_{1}^{2}+2^{\ell_{2}} y_{2}^{2}+\left(2^{\ell_{1}} y_{1}+2^{\ell_{2}} y_{2}\right)^{2}
$$

is isotropic. We have

$$
\begin{aligned}
Q\left(y_{1}, y_{2}\right) & =2^{\ell_{1}}\left(1+2^{\ell_{1}}\right) y_{1}^{2}+2^{\ell_{2}}\left(1+2^{\ell_{2}}\right) y_{2}^{2}+2^{\ell_{1}+\ell_{2}+1} y_{1} y_{2} \\
& =2^{\ell_{1}}\left(1+2^{\ell_{1}}\right)\left(y_{1}+\frac{2^{\ell_{1}+\ell_{2}}}{2^{\ell_{1}}\left(1+2^{\ell_{1}}\right)} y_{2}\right)^{2}+\left(2^{\ell_{2}}\left(1+2^{\ell_{2}}\right)-\frac{2^{2 \ell_{1}+2 \ell_{2}}}{2^{\ell_{1}}\left(1+2^{\ell_{1}}\right)}\right) y_{2}^{2}
\end{aligned}
$$

and so

$$
Q \cong\left\langle 2^{\ell_{1}}\left(1+2^{\ell_{1}}\right), \frac{2^{\ell_{1}+\ell_{2}}\left(1+2^{\ell_{1}}\right)\left(1+2^{\ell_{2}}\right)-2^{2 \ell_{1}+2 \ell_{2}}}{2^{\ell_{1}}\left(1+2^{\ell_{1}}\right)}\right\rangle .
$$

This gives us

$$
\begin{aligned}
2^{\ell_{1}}\left(1+2^{\ell_{1}}\right) Q & \cong\left\langle 1,2^{\ell_{1}+\ell_{2}}\left(\left(1+2^{\ell_{1}}\right)\left(1+2^{\ell_{2}}\right)-2^{\ell_{1}+\ell_{2}}\right)\right\rangle \\
& =\left\langle 1,2^{\ell_{1}+\ell_{2}}\left(1+2^{\ell_{1}}+2^{\ell_{2}}\right)\right\rangle \\
& =\left\langle 1,2^{k_{1}+k_{2}-2 k_{3}} \cdot 2^{-k_{3}}\left(2^{k_{1}}+2^{k_{2}}+2^{k_{3}}\right)\right\rangle \\
& =\left\langle 1,2^{k_{1}+k_{2}-3 k_{3}} n\right\rangle \\
& \cong\left\langle 1,2^{k_{1}+k_{2}+k_{3}} n\right\rangle
\end{aligned}
$$

which is isotropic if and only if $Q$ is isotropic, if and only if $-2^{k_{1}+k_{2}+k_{3}} n$ is a square in $F_{0}$.

## $4.3 m=4$

Proposition 4.3.1. Suppose $m=4$. If either $\operatorname{char} F_{0} \neq 0$, or both $\sqrt{-1} \in F_{0}$ and $\sqrt{2} \in F_{0}$, then (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$.

Here a nonreal field is one in which -1 is a sum of squares. We will prove a more general lemma here, giving more detailed necessary and sufficient conditions for the existence of solutions to (4.1) by considering different classes of quadratic forms.

In the following lemma we will use the following result: If the quadratic form $\left\langle 2^{k_{1}}, 2^{k_{2}}, \ldots, 2^{k_{m}}\right\rangle$ is equivalent to a scalar multiple of a form with Schur index 2 (a form with a 2-dimensional totally isotropic subspace), then (4.1) has a solution. This is implied by the 2-dimensional subspace having nontrivial intersection with the line defined by $2^{k_{1}} y_{1}+\cdots+2^{k_{m}} y_{m}=0$.

We will also use the result that in any field, the smallest $r$ such that -1 is a sum of $r$ squares is a power of 2 (see [10, Theorem 3.1.3]).

Lemma 4.3.2. Let $m=4$. The quadratic form $\left\langle 2^{k_{1}}, 2^{k_{2}}, 2^{k_{3}}, 2^{k_{4}}\right\rangle$ is equivalent to a scalar multiple of one of:

- $\langle 1,1,1,1\rangle$. In this case (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$ if and only if -1 is a sum of 2 squares in $F_{0}$.
- $\langle 1,1,1,2\rangle$ or $\langle 1,1,2,2\rangle$. In these cases, if $\sqrt{2} \in F_{0}$ and -1 is a sum of 2 squares in $F_{0}$, then (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$.

If (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$, then -1 is a sum of 4 squares in $F_{0}$.
Proof. By removing squares, reordering and possibly multiplying by 2 , the quadratic form $\left\langle 2^{k_{1}}, 2^{k_{2}}, 2^{k_{3}}, 2^{k_{4}}\right\rangle$ is equivalent to a scalar multiple of either $\langle 1,1,1,1\rangle,\langle 1,1,1,2\rangle$ or $\langle 1,1,2,2\rangle$. We split these in to two cases.

- Suppose $\left\langle 2^{k_{1}}, 2^{k_{2}}, 2^{k_{3}}, 2^{k_{4}}\right\rangle$ is equivalent to a multiple of $\langle 1,1,1,1\rangle$, which is a Pfister form $\langle 1,1,1,1\rangle=\langle 1,1\rangle \otimes\langle 1,1\rangle$. Suppose -1 is a sum of 2 squares in $F_{0}$. Then $\langle 1,1,1,1\rangle$ is isotropic. By [9, X. Theorem 1.7], an isotropic Pfister form is hyperbolic and hence has Schur index 2 , thus (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$.

Conversely, suppose (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$. Then $\langle 1,1,1,1\rangle$ is isotropic, and so -1 is a sum of 3 , and hence 2 squares in $F_{0}$.

- Suppose $\left\langle 2^{k_{1}}, 2^{k_{2}}, 2^{k_{3}}, 2^{k_{4}}\right\rangle$ is equivalent to a multiple of $\langle 1,1,1,2\rangle$ or $\langle 1,1,2,2\rangle$. Suppose $\sqrt{2} \in F_{0}$ and -1 is a sum of 2 squares in $F_{0}$. Then $\left\langle 2^{k_{1}}, 2^{k_{2}}, 2^{k_{3}}, 2^{k_{4}}\right\rangle \cong\langle 1,1,1,1\rangle$. By the above case, this then has Schur index 2 and so (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$.
Conversely, suppose (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$. Then either $\langle 1,1,1,2\rangle$ or $\langle 1,1,2,2\rangle$ is isotropic. Both cases imply that -1 is a sum of 5 , and hence 4 squares in $F_{0}$.

We now consider Proposition 4.3.1. If both $\sqrt{-1} \in F_{0}$ and $\sqrt{2} \in F_{0}$ then (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$ as a special case of Lemma 4.3.2. If $\operatorname{char} F_{0} \neq 0$, then (4.1) has a solution by Lemma 3.1(3).

## $4.4 m=5$

Proposition 4.4.1. Suppose $m=5$. If $\sqrt{-1} \in F_{0}$, then (4.1) has a solution in $\mathbb{P}^{4}\left(F_{0}\right)$.
Again, we will prove a more general lemma here, giving more detailed necessary and sufficient conditions for the existence of solutions to (4.1).

Lemma 4.4.2. Let $m=5$. The quadratic form $\left\langle 2^{k_{1}}, 2^{k_{2}}, 2^{k_{3}}, 2^{k_{4}}, 2^{k_{5}}\right\rangle$ is equivalent to a scalar multiple of one of:

- $\langle 1,1,1,2,2\rangle$. In this case, if $\sqrt{-2} \in F_{0}$ or $\sqrt{-1} \in F_{0}$, then (4.1) has a solution in $\mathbb{P}^{4}\left(F_{0}\right)$.
$\cdot\langle 1,1,1,1,1\rangle$ or $\langle 1,1,1,1,2\rangle$. In these cases, if $\sqrt{-1} \in F_{0}$, then (4.1) has a solution in $\mathbb{P}^{4}\left(F_{0}\right)$. If (4.1) has a solution in $\mathbb{P}^{3}\left(F_{0}\right)$, then -1 is a sum of 4 squares in $F_{0}$.

Proof. By removing squares, reordering and possibly multiplying by 2 , the quadratic form $\left\langle 2^{k_{1}}, \ldots, 2^{k_{5}}\right\rangle$ is equivalent to a scalar multiple of either $\langle 1,1,1,1,1\rangle,\langle 1,1,1,1,2\rangle$ or $\langle 1,1,1,2,2\rangle$. We consider each case in turn.

- Suppose $\left\langle 2^{k_{1}}, \ldots, 2^{k_{5}}\right\rangle$ is equivalent to a scalar multiple of $\langle 1,1,1,2,2\rangle$. If $\sqrt{-2} \in F_{0}$, then $\langle 1,1,1,2,2\rangle \cong\langle 1,-1,1,-1,1\rangle$ which has Schur index 2 and so there exists a solution to (4.1). If $\sqrt{-1} \in F_{0}$, then by Lemma 3.1(2) there exists a solution to (4.1).
- Otherwise $\left\langle 2^{k_{1}}, \ldots, 2^{k_{5}}\right\rangle$ is equivalent to a multiple of $\langle 1,1,1,1, c\rangle$ for some $c \in\{1,2\}$. If $\sqrt{-1} \in F_{0}$, then $\langle 1,1,1,1, c\rangle \cong\langle 1,-1,1,-1, c\rangle$, which has Schur index 2 and so has nontrivial intersection with the line defined by $2^{k_{1}} y_{1}+2^{k_{2}} y_{2}+2^{k_{3}} y_{3}+2^{k_{4}} y_{4}+2^{k_{5}} y_{5}=0$. Thus there exists a solution to (4.1).

Conversely, suppose (4.1) has a solution in $\mathbb{P}^{4}\left(F_{0}\right)$. Then one of $\langle 1,1,1,2,2\rangle,\langle 1,1,1,1,2\rangle$ or $\langle 1,1,1,1,1\rangle$ is isotropic. In all cases we must have -1 is a sum of 6 , and hence 4 squares in $F_{0}$.

## $4.5 m \geq 6$

Proposition 4.5.1. Suppose $m \geq 6$. If $F_{0}$ is a nonreal number field, then (4.1) has a solution in $\mathbb{P}^{4}\left(F_{0}\right)$.

Proof. Suppose $m \geq 6$ and $F_{0}$ is a nonreal number field. After dividing out by $2^{k_{m}}$ we can substitute the first equation of (4.1) into the second to obtain a quadratic form of degree $m-1$ representing the system (4.1). Over a nonreal number field, any quadratic form of dimension $\geq 5$ is isotropic (see [9, XI, Corollary 1.5, 379]), and thus we have a solution to (4.1).

### 4.6 Proof of Theorem 4.2

The proof of Theorem 4.2 has mostly been summarized in propositions of the earlier subsections. We need to extend Propositions 4.3.1 and 4.4.1 to arbitrary $m \geq 4$ and $m \geq 5$ respectively. We will need the following:

For a fixed field $F_{0}$, if (4.1) has a solution in $\mathbb{P}^{m-1}\left(F_{0}\right)$ for all $n$ with $m$ digits in its base 2 expansion, then (4.1) will have a solution for all $n$ with $\geq m$ digits in its base 2 expansion.

Suppose $n=2^{\ell_{1}}+\cdots+2^{\ell_{m^{\prime}}}+2^{k_{1}}+\cdots+2^{k_{m}}$ with $\ell_{1}>\cdots>\ell_{m^{\prime}}>k_{1}>\cdots>k_{m}$ and let $\left(a_{1}\right.$ : $\cdots: a_{m}$ ) be a solution to (4.1) in $\mathbb{P}^{m-1}\left(F_{0}\right)$ for $2^{k_{1}}+\cdots+2^{k_{m}}$. Then we obtain a solution to (4.1) in $\mathbb{P}^{\ell+m-1}\left(F_{0}\right)$ for $2^{\ell_{1}}+\cdots+2^{\ell_{m^{\prime}}}+2^{k_{1}}+\cdots+2^{k_{m}}$ being $\left(0: \cdots: 0: a_{1}: \cdots: a_{m}\right)$.

Lastly we need to show that if (4.1) has a solution in $\mathbb{P}^{m-1}\left(F_{0}\right)$ then $F_{0}$ must be nonreal. Suppose (4.1) has a solution in $\mathbb{P}^{m-1}\left(F_{0}\right)$. Then as each of the $2^{k_{i}}$ are positive integers, the second equation gives us 0 as the sum of $n$ squares, not all zero:

$$
y_{1}^{2}+\cdots+y_{1}^{2}+\cdots+y_{m}^{2}+\cdots+y_{m}^{2}=0 .
$$

By dividing out by a nonzero $y_{i}$ we get -1 as a sum of squares. This completes the Theorem.

## Chapter 5

## The Hermite-Joubert problem for $p=3$

Theorem 5.1. Let $n=3^{k_{1}}+\cdots+3^{k_{m}} \geq 0$ be the base 3 representation of $n$, with $k_{1} \geq \cdots \geq k_{m}$, and $F_{0}$ be a field with characteristic $\neq 3$. Consider the system of equations:

$$
\begin{align*}
& 3^{k_{1}} y_{1}+\cdots+3^{k_{m}} y_{m}=0 \\
& 3^{k_{1}} y_{1}^{3}+\cdots+3^{k_{m}} y_{m}^{3}=0 \tag{5.1}
\end{align*}
$$

The following statements hold:

- If $m=1$ then (5.1) has no solution in $\mathbb{P}^{0}\left(F_{0}\right)$.
- If $m=2$ then (5.1) has a solution in $P^{1}\left(F_{0}\right)$ if and only if char $F_{0}$ divides $3^{2\left(k_{1}-k_{2}\right)}-1$ or $k_{1}=k_{2}$.
- If $m=3, F_{0}=\mathbb{Q}$, the $k_{i}$ are distinct and $k_{1}+k_{2}+k_{3} \not \equiv 2(\bmod 3)$, then $(5.1)$ has no solution in $P^{2}\left(F_{0}\right)$.
- If $m=3, F_{0}=\mathbb{Q}$, the $k_{i}$ are distinct and $k_{1} \not \equiv k_{2}(\bmod 3)$, then $(5.1)$ has no solution in $P^{2}\left(F_{0}\right)$.
- If $m \geq 5$ and char $F_{0} \neq 0$ then (5.1) has a solution in $P^{m-1}\left(F_{0}\right)$.
- If $m \geq 6$ then (5.1) has a solution in $P^{1}\left(F_{0}\right)$.
- If the $k_{i}$ are not distinct then (5.1) has a solution in $P^{1}\left(F_{0}\right)$.

We will first prove the statements in the theorem for when $m=3$. Define an equivalence relation on $m$-tuples of integers by $\left(k_{1}, \ldots, k_{m}\right) \sim\left(\ell_{1}, \ldots \ell_{m}\right)$ if and only if there exist $\sigma \in \mathrm{S}_{m}$ and $\lambda \in \mathbb{Z}$ such that $k_{i} \equiv \ell_{\sigma(i)}+\lambda(\bmod 3)$ for each $i$.

Remark 5.2. If $\left(k_{1}, \ldots, k_{m}\right) \sim\left(\ell_{1}, \ldots \ell_{m}\right)$ then the solutions to $3^{k_{1}} y_{1}^{3}+\cdots+3^{k_{m}} y_{m}^{3}=0$ in $\mathbb{P}^{m-1}\left(F_{0}\right)$ are in bijective correspondence with the solutions to $3^{\ell_{1}} z_{1}^{3}+\cdots+3^{\ell_{m}} z_{m}^{3}=0$ in $\mathbb{P}^{m-1}\left(F_{0}\right)$. If $\ell_{i}=$ $k_{\sigma(i)}+\lambda+3 s_{i}$, then this correspondence is given by $z_{i}=3^{-s_{i}} y_{\sigma(i)}$. In particular, if all solutions to $3^{k_{1}} y_{1}^{3}+\cdots+3^{k_{m}} y_{m}^{3}=0$ have a zero coordinate $y_{i}$, then all solutions to $3^{\ell_{1}} z_{1}^{3}+\cdots+3^{\ell_{m}} z_{m}^{3}=0$ have a zero coordinate $z_{j}$.

Lemma 5.3. Let $k_{1}>k_{2}>k_{3} \geq 0$ be integers such that $k_{1}+k_{2}+k_{3} \not \equiv 2(\bmod 3)$. Then the system of equations

$$
\begin{align*}
& 3^{k_{1}} y_{1}+3^{k_{2}} y_{2}+3^{k_{3}} y_{3}=0 \\
& 3^{k_{1}} y_{1}^{3}+3^{k_{2}} y_{2}^{3}+3^{k_{3}} y_{3}^{3}=0 \tag{5.2}
\end{align*}
$$

has no solution in $\mathbb{P}^{2}(\mathbb{Q})$.
Proof. Suppose we have a solution $\left(y_{1}: y_{2}: y_{3}\right)$ to (5.2) with a zero coordinate, say $y_{i}=0$. Then this gives rise to a solution in $\mathbb{P}^{1}(\mathbb{Q})$ to $(5.1)$ for $n^{\prime}=n-3^{k_{i}}$, which has only two terms in its base 3 expansion. By Lemma 3.2, this is impossible in characteristic 0 . Thus it suffices to show that all solutions $3^{k_{1}} y_{1}^{3}+3^{k_{2}} y_{2}^{3}+3^{k_{3}} y_{3}^{3}=0$ have a zero coordinate.

By Remark 5.2, it suffices to show that for some $\left(\ell_{1}, \ldots \ell_{m}\right) \sim\left(k_{1}, k_{2}, k_{3}\right)$ the only solutions to $3^{\ell_{1}} y_{1}^{3}+3^{\ell_{2}} y_{2}^{3}+3^{\ell_{3}} y_{3}^{3}=0$ have a zero coordinate.

Under the equivalence relation described above, $\left(k_{1}, k_{2}, k_{3}\right)$ is equivalent to one of $(0,0,0)$, $(0,0,1),(0,1,2)$, or $(0,0,2)$. For this lemma we need only consider the first three, noting that if $\left(k_{1}, k_{2}, k_{3}\right) \sim\left(\ell_{1}, \ldots \ell_{m}\right)$, then $k_{1}+k_{2}+k_{3} \equiv \ell_{1}+\ell_{2}+\ell_{3}(\bmod 3)$.

Thus to prove the lemma it suffices to show that the only solutions in $\mathbb{P}^{2}(\mathbb{Q})$ to $3^{\ell_{1}} y_{1}^{3}+3^{\ell_{2}} y_{2}^{3}+$ $3^{\ell_{3}} y_{3}^{3}=0$ have a zero coordinate, for each case of $\left(\ell_{1}, \ell_{2}, \ell_{3}\right) \in\{(0,0,0),,(0,0,1),(0,1,2)\}$. We consider each case in turn.

- The only solutions to $y_{1}^{3}+y_{2}^{3}+y_{3}^{3}=0$ have a zero coordinate, as a result of Fermat's Last Theorem.
- By [13], the only solution to $y_{1}^{3}+y_{2}^{3}+3 y_{3}^{3}=0$ is $(1:-1: 0)$.
- Suppose $\left(y_{1}: y_{2}: y_{3}\right)$ is a solution to $y_{1}^{3}+3 y_{2}^{3}+9 y_{3}^{3}=0$. We may choose a representation satisfying $\operatorname{gcd}\left(y_{1}, y_{2}, y_{3}\right)=1$. Then we have $3 \mid y_{1}$, and hence $3 \mid y_{2}$ and $3 \mid y_{3}$. This contradicts the minimality of $\left(y_{1}, y_{2}, y_{3}\right)$.

Lemma 5.4. Let $k_{1}>k_{2}>k_{3} \geq 0$ be integers such that $k_{1} \not \equiv k_{2}(\bmod 3)$. Then the system of equations

$$
\begin{align*}
& 3^{k_{1}} y_{1}+3^{k_{2}} y_{2}+3^{k_{3}} y_{3}=0 \\
& 3^{k_{1}} y_{1}^{3}+3^{k_{2}} y_{2}^{3}+3^{k_{3}} y_{3}^{3}=0 \tag{5.3}
\end{align*}
$$

has no solution in $\mathbb{P}^{2}(\mathbb{Q})$.
Proof. Assume the contrary and suppose we have a nontrivial solution $\left(y_{1}, y_{2}, y_{3}\right) \neq(0,0,0)$ to (5.3) with $y_{1}, y_{2}, y_{3} \in \mathbb{Q}$. By dividing each equation by $3^{k_{3}}$ and replacing $k_{1}$ and $k_{2}$ by $k_{1}-k_{3}$ and $k_{2}-k_{3}$ respectively, we may assume without loss of generality that $k_{3}=0$. Now the second equation is $y_{3}=-3^{k_{1}} y_{1}-3^{k_{2}} y_{2}$, and by substituting this into the first, we gain a nontrivial solution to

$$
\begin{equation*}
3^{k_{1}} y_{1}^{3}+3^{k_{2}} y_{2}^{3}-3^{3 k_{1}} y_{1}^{3}-3^{2 k_{1}+k_{2}+1} y_{1}^{2} y_{2}-3^{k_{1}+2 k_{2}+1} y_{1} y_{2}^{2}-3^{3 k_{2}} y_{2}^{3}=0 . \tag{5.4}
\end{equation*}
$$

As $\left(y_{1}, y_{2}, y_{3}\right)$ is nontrivial we must have one of $y_{1}$ and $y_{2}$ are nonzero, and hence by (5.4) both are nonzero. Define

$$
\begin{aligned}
M_{1} & :=v_{3}\left(3^{k_{1}} y_{1}^{3}\right)=k_{1}+3 v 3\left(y_{1}\right), \\
M_{2} & :=v_{3}\left(3^{k_{2}} y_{2}^{3}\right)=k_{2}+3 v_{3}\left(y_{2}\right), \quad \text { and } \\
M & :=\min \left(M_{1}, M_{2}\right) .
\end{aligned}
$$

Here $v_{3}$ denotes the 3 -adic valuation. Since $k_{1} \not \equiv k_{2}(\bmod 3)$, we have $M_{1} \neq M_{2}$, so that $v_{3}\left(3^{k_{1}} y_{1}^{3}+\right.$ $\left.3^{k_{2}} y_{2}^{3}\right)=M$. We now show that the 3-adic valuation of the rest of the terms in the left hand side of (5.4) is larger than $M$, which would imply that the 3 -adic valuation of the entire equation is $M$, contradicting it being 0 . We consider each term separately:

- $v_{3}\left(3^{3 k_{1}} y_{1}^{3}\right)=3 k_{1}+3 v_{3}\left(y_{1}\right)>M_{1} \geq M$
- $v_{3}\left(3^{2 k_{1}+k_{2}+1} y_{1}^{2} y_{2}\right)=2 k_{1}+k_{2}+2 v_{3}\left(y_{1}\right)+v_{3}\left(y_{2}\right)+1>\frac{2}{3} M_{1}+\frac{1}{3} M_{2}>\frac{2}{3} M+\frac{1}{3} M=M$
- $v_{3}\left(3^{k_{1}+2 k_{2}+1} y_{1} y_{2}^{2}\right)=k_{1}+2 k_{2}+v_{3}\left(y_{1}\right)+2 v_{3}\left(y_{2}\right)+1>\frac{1}{3} M_{1}+\frac{2}{3} M_{2}>\frac{2}{3} M+\frac{1}{3} M=M$
- $v_{3}\left(3^{3 k_{2}} y_{2}^{3}\right)=3 k_{2}+3 v_{3}\left(y_{2}\right)>M_{2} \geq M$

Thus we have obtained a contradiction in the 3-adic valuation of the equation and so there are no solutions to (5.3) when $k_{1} \not \equiv k_{2}(\bmod 3)$.

Proof of Theorem 5.1. If $m=1$ then $\mathbb{P}^{0}\left(F_{0}\right)$ consists of a point, which is not a solution to (5.1).
If $m=2$, our result follows from Lemma 3.2.
If $m=3$, our result follows from Lemmas 5.3 and 5.4.
If the $k_{i}$ are not distinct, say $k_{i}=k_{j}$, then we obtain a solution $y=\left(y_{1}: \cdots: y_{m}\right)$ by setting $y_{i}=1, y_{j}=-1$ and $y_{h}=0$ for $h \neq i, j$.

The other cases are covered by Lemma 3.1(2) and (3).

Note that Lemma 3.1(1) shows that in characteristic $\neq 3$ there exist solutions to (5.1) when $m=$ 4 and $\left(k_{1}, k_{2}, k_{3}, k_{4}\right) \sim(0,0,1,1)$ or when $m=5$ and $\left(k_{1}, k_{2}, k_{3}, k_{4}, k_{5}\right) \nsim(0,0,0,1,2)$. Except for the remaining cases of $m=3,4$ or 5 , being when $\left(k_{1}, \ldots, k_{m}\right)$ is equivalent to one of $(0,0,2),(0,0,1,2)$ or $(0,0,0,1,2)$, we have a complete description of when there exist solutions to (5.1) in $\mathbb{P}^{m-1}(\mathbb{Q})$.

Proof of Theorem 1.4. Theorem 1.4 readily follows from Theorem 1.3 and Theorem 5.1

## Chapter 6

## Numerical evidence in support of Conjecture 1.5

We now consider Conjecture 1.5. By Theorem 1.3, it is equivalent to the following
Conjecture 6.1. Let $k_{1}>k_{2}>0$. Then the system of equations

$$
\begin{align*}
& 3^{k_{1}} y_{1}+3^{k_{2}} y_{2}+y_{3}=0 \\
& 3^{k_{1}} y_{1}^{3}+3^{k_{2}} y_{2}^{3}+y_{3}^{3}=0 \tag{6.1}
\end{align*}
$$

has no solution in $\mathbb{P}^{2}(\mathbb{Q})$.
Note in this statement we have assumed $k_{3}=0$. This is without loss of generality, as in the proof of Lemma 5.4.

The remaining case left unsolved by Lemmas 5.3 and 5.4 is when $k_{1} \equiv k_{2}(\bmod 3)$ and $k_{1}+$ $k_{2} \equiv 2(\bmod 3)$. That is, when $k_{1} \equiv k_{2} \equiv 1(\bmod 3)$. Indeed in this case we have $\left(k_{1}, k_{2}, 0\right) \sim$ $(0,0,2)$, but there do exist nonzero solutions to

$$
y_{1}^{3}+y_{2}^{3}+9 y_{3}^{3}
$$

as illustrated by $(1: 2:-1)$. However the set of solutions to this equation (after a change of variables) does not necessarily intersect with $3^{k_{1}} y_{1}+3^{k_{2}} y_{2}+y_{3}=0$ to give a solution to (6.1), as we have already seen when $k_{1} \not \equiv k_{2}(\bmod 3)$.

The smallest $n$ of this form, not covered by Theorem 1.4, is $85=3^{4}+3^{1}+3^{0}$. In this case we have shown computationally that (6.1) does not have a solution in $\mathbb{P}^{2}(\mathbb{Q})$. In fact we have covered
a lot more.
Proposition 6.2. Theorem 1.4 remains true for all $k_{1}>k_{2}>k_{3} \geq 0$ such that $k_{2}-k_{3}<k_{1}-k_{3}<$ 30000.

We introduce some notation before giving proofs of our computational evidence. For $a, b$ nonnegative integers, define

$$
f_{a, b}:=\left(3^{6 a+2}-1\right) x^{3}+3^{4 a+2 b+3} x^{2}+3^{2 a+4 b+3} x+\left(3^{6 b+2}-1\right) .
$$

Here $f_{a, b} \in \mathbb{Z}[x]$ is primitive. For abbreviation, we write $l_{a, b}$ to denote the leading coefficient of $f_{a, b}$ (that is $l_{a, b}:=3^{6 a+2}-1$ ).

Lemma 6.3. Let $k_{1} \geq k_{2}>0$ with $k_{1} \equiv k_{2} \equiv 1(\bmod 3)$. Write $k_{1}=3 a+1$ and $k_{2}=3 b+1$. Then the system (6.1) has a rational solution in $\mathbb{P}^{2}(\mathbb{Q})$ if and only if $f_{a, b}$ has a rational root.

Moreover, suppose there exists an integer $m$ such that $\left(l_{a, b}, m\right)=1$ and $f_{a, b}(z) \not \equiv 0(\bmod m)$ for all integers $0 \leq z<m$. Then $f_{a, b}$ does not have a rational root.

Proof. The system (6.1) is equivalent to

$$
3^{k_{1}} y_{1}^{3}+3^{k_{2}} y_{2}^{3}-3^{3 k_{1}} y_{1}^{3}-3^{2 k_{1}+k_{2}+1} y_{1}^{2} y_{2}-3^{k_{1}+2 k_{2}+1} y_{1} y_{2}^{2}-3^{3 k_{2}} y_{2}^{3}=0 .
$$

We may consider this as a cubic polynomial in one variable $y=y_{1} / y_{2}$ :

$$
3^{k_{1}}\left(1-3^{2 k_{1}}\right) y^{3}-3^{2 k_{1}+k_{2}+1} y^{2}-3^{k_{1}+2 k_{2}+1} y+3^{k_{2}}\left(1-3^{2 k_{1}}\right)=0 .
$$

After taking out a factor of $-3^{k_{2}}$ and writing $k_{1}=3 a+1, k_{2}=3 b+1$ for integers $a>b \geq 0$, this becomes

$$
3^{3 a-3 b}\left(3^{6 a+2}-1\right) y^{3}+3^{6 a+3} y^{2}+3^{3 a+3 b+3} y+\left(3^{6 b+2}-1\right)=0 .
$$

After the substitution $y=3^{b-a} x$, we then have that (6.1) has a rational solution in $\mathbb{P}^{2}(\mathbb{Q})$ if and only if

$$
f_{a, b}=\left(3^{6 a+2}-1\right) x^{3}+3^{4 a+2 b+3} x^{2}+3^{2 a+4 b+3} x+\left(3^{6 b+2}-1\right)
$$

has a rational root.
Now suppose there exists an integer $m$ such that $\left(l_{a, b}, m\right)=1$ and $f_{a, b}(z) \not \equiv 0(\bmod m)$ for all integers $0 \leq z<m$.

As $f_{a, b}$ is primitive, by Gauss's lemma it suffices to show that $f_{a, b}$ is irreducible over the integers. By assumption we have $f_{a, b}$ is irreducible modulo $m$, and as $\left(l_{a, b}, m\right)=1$ we thus have $f_{a, b}$ is irreducible over the integers.

The most efficient algorithm we have found to prove that $f_{a, b}$ has no rational root is to find some small prime $q$ such that $q$ does not divide $l_{a, b}$, and $f_{a, b}$ is irreducible modulo $q$, as in Lemma 6.3. For the $b<a<10000$ of which we checked, we found primes $q$ (dependent on $a$ and $b$ ) smaller than 500 for which $q$ does not divide $l_{a, b}$ and $f_{a, b}$ is irreducible modulo $q$.

Remark 6.4. Note that there is no such absolute bound on primes independent of $a$ and $b$ for which this algorithm will work. The equation $f_{a, b}$ does indeed have solutions if $a=b$. Thus if all the primes less than some integer $B$ appear in the prime factorization of both $a$ and $b$, then $f_{a, b}$ will have solutions modulo $q$ for all primes less than $B$.

To aid in computation, we first establish some preliminary results.
For tuples of integers we say $(a, b) \equiv(x, y)(\bmod m)$ if both $a \equiv x(\bmod m)$ and $b \equiv y(\bmod m)$, so there are $m^{2}$ classes of tuples of integers modulo $m$.

Lemma 6.5. Let $m$ be a positive integer with $(3, m)=1$.
(a) Suppose $\left(a^{\prime}, b^{\prime}\right) \equiv(a, b)\left(\bmod \operatorname{ord}_{m}(3)\right), f_{a, b}$ is irreducible modulo $m$, and $\left(l_{a, b}, m\right)=1$. Then $f_{a^{\prime}, b^{\prime}}$ has no rational root.
(b) Let o be a positive integer multiple of $\operatorname{ord}_{m}(3)$. Suppose $f_{a, b}$ is irreducible modulo $m$ and $\left(l_{a, b}, m\right)=1$. Then $f_{a^{\prime}, b^{\prime}}$ has no rational solutions if $\left(a^{\prime}, b^{\prime}\right) \cong\left(a+i \operatorname{ord}_{m}(3), b+j \operatorname{ord}_{m}(3)\right)$ for any $i, j \in \mathbb{Z}$.
(c) Let $r$ be a positive integer divisor of $\operatorname{ord}_{m}(3)$ and write $\operatorname{ord}_{m}(3)=r s$. Suppose $f_{a, b}$ is irreducible modulo $m$ and $\left(l_{a, b}, m\right)=1$. If $f_{r i+a, r j+b}$ is irreducible modulo $m$ and $\left(l_{r i+a, r j+b}, m\right)=1$ for all tuples $(r i+a, r j+b)$ with $0 \leq i, j<s$, then $f_{a^{\prime}, b^{\prime}}$ has no rational solutions for any $\left(a^{\prime}, b^{\prime}\right) \cong$ $(a, b)(\bmod r)$.

Proof. (a) As $a$ and $b$ only appear in $f_{a, b}$ in exponents of 3 , and $\left(a^{\prime}, b^{\prime}\right) \equiv(a, b)\left(\bmod \operatorname{ord}_{m}(3)\right)$, we have $f_{a^{\prime}, b^{\prime}} \equiv f_{a, b}(\bmod m)$. In particular we have $l_{a^{\prime}, b^{\prime}} \equiv l_{a, b}(\bmod m)$ and $f_{a^{\prime}, b^{\prime}}$ is irreducible modulo $m$. Thus $f_{a^{\prime}, b^{\prime}}$ is irreducible over the integers and hence as $f_{a^{\prime}, b^{\prime}}$ is primitive, it is irreducible over the rationals.

Part (b) follows from (a) by the projection $\mathbb{Z} /(o \mathbb{Z}) \rightarrow \mathbb{Z} /\left(\operatorname{ord}_{m}(3) \mathbb{Z}\right)$. Part (c) follows from (a) by the natural isomorphism $\mathbb{Z} /(s \mathbb{Z}) \rightarrow(r \mathbb{Z}) /\left(\operatorname{ord}_{m}(3) \mathbb{Z}\right)$ of additive groups.

For a fixed modulus $o$, parts (b) and (c) of Lemma 6.5 give us a way of searching for classes of tuples $(a, b)$ modulo $o$ for which $f_{a^{\prime}, b^{\prime}}$ has no rational root, for any $\left(a^{\prime}, b^{\prime}\right) \equiv(a, b)(\bmod o)$.

Proposition 6.6. Suppose $(a, b)$ is not congruent to one of

$$
(0,0),(1,1),(2,2),(3,3),(4,4),(1,3),(2,3),(3,1),(3,2),(3,4),(4,3)
$$

modulo 5. Then $f_{a, b}$ has no rational solution.
Proof. By Lemma 6.5 (c), it suffices to find an integer $m$ where $\operatorname{ord}_{m}(3)=5 s$ for some integer $s$, such that $f_{5 i+a, 5 j+b}$ is irreducible modulo $m$ and $\left(l_{r i+a, r j+b, m}\right)=1$ for all $0 \leq i, j<s$.

Suppose $(a, b)$ is congruent to one of

$$
(0,1),(0,3),(0,4),(1,0),(1,2),(2,1),(2,4),(3,0),(4,0),(4,2)
$$

modulo 5 . Let $m=31$, where we have $\operatorname{ord}_{31}(3)=30$. By checking computationally, we prove that the assumptions of Lemma 6.5 (c) are satisfied. For more detail, see Appendix A.

Otherwise, $(a, b)$ is congruent to one of

$$
(0,1),(0,2),(1,0),(1,4),(2,0),(4,1)
$$

modulo 5. Now let $m=61$, where we have $\operatorname{ord}_{61}(3)=10$. Again, by checking computationally, we prove that the assumptions of Lemma 6.5 (c) are satisfied.

This proposition allows the algorithm based on Lemma 6.3 to skip many values of $a$ and $b$.
We return to Proposition 6.2. Without loss of generality we may assume $k_{3}=0$ and $k_{2}<$ $k_{1}<30000$. By Theorem 1.4 we may then assume $k_{1} \equiv k_{2} \equiv 1(\bmod 3)$. Let $a=\left(k_{1}-1\right) / 3$ and $b=\left(k_{2}-1\right) / 3$. Then $b<a<10000$. By Lemma 6.3, it suffices to prove that $f_{a, b}$ has no rational roots. By Proposition 6.6, we may assume $(a, b)$ is congruent to one of

$$
(0,0),(1,1),(2,2),(3,3),(4,4),(1,3),(2,3),(3,1),(3,2),(3,4),(4,3)
$$

modulo 5. We ensured computationally (see Appendix A) that there exists a prime $p$ such that $f_{a, b}$ has no roots modulo $p$ and $\left(l_{a, b}, p\right)=1$. Thus by Lemma $6.3 f_{a, b}$ has no rational root. This completes the proof of the proposition.

Remark 6.7. Testing with other orders $o$ of 3 (as opposed to 5) has ruled out a larger fraction of tuples in $(\mathbb{Z} / o \mathbb{Z})^{2}$. By checking moduli up to 300 , we have ruled out $14 / 25$ of the tuples $(a, b)$ modulo 5. The same approach has ruled out 42548/44100 tuples modulo 210. Unfortunately, due to the large scale, this fact can not be easily used to optimize the algorithm proving no rational solutions to $f_{a, b}$.

Remark 6.8. An alternate approach to proving the conjecture would be to consider $x^{3}+y^{3}+9 z^{3}=$ 0 as an elliptic curve. When $k_{1}+k_{2}+k_{3} \equiv 2(\bmod 3)$, after substitution the system of equations becomes equivalent to

$$
\begin{aligned}
& x^{3}+y^{3}+9 z^{3}=0 \\
& 3^{a} x+3^{b} y+z=0
\end{aligned}
$$

for some integers $a>b>0$. The group of rational points lying on the elliptic curve $x^{3}+y^{3}+9 z^{3}=0$ is cyclic, generated by $(1: 2:-1)$, see [13].

## Chapter 7

## Conclusion

In this thesis we addressed a generalization of the classical Hermite-Joubert problem, as in Question 1.1. By restricting to the case of $p$-fields, we gave a necessary and sufficient condition in Theorem 1.3 to answer the analogous Question 1.2. This condition is phrased in terms of the existence of an $F$-point on a particular variety, depending on $n$ and $p$. In Theorems 4.1 and 4.2 we obtain an almost complete description of the Hermite-Joubert problem in the case of $p=2$.

Theorem 1.4 yields us new examples in the case of $p=3$ where the Hermite-Joubert problem has a negative solution, the smallest being where $n=13$. In Chapter 6 we give numerical evidence in support of Conjecture 1.5, showing more examples where the Hermite-Joubert problem has a negative solution.

A further potential research direction would be to give a proof of Conjecture 1.5. Theorem 1.3 considers the existence of an element $a \in E$ such that $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{p}\right)=0$. In the case of $p=2$ or $p=3$ this is analagous to the classical Hermite-Joubert problem, where this condition is equivalent to the $1^{\text {st }}$ and $p^{\text {th }}$ coefficients of the characteristic polynomial of $a$ being zero. Another possible direction would be to study the existence of an element $a$ such that the $1^{\text {st }}$ and $p^{\text {th }}$ coefficients of the characteristic polynomial of $a$ are zero, for a more general $p$.

In the case of $p=2$, for particular separable field extensions $E / F$, we showed that there exists an element $a \in E$ such that $\operatorname{tr}_{E / F}(a)=\operatorname{tr}_{E / F}\left(a^{2}\right)=0$. It would be of interest to see how this element might be constructed.

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## Appendix A

## Maple code

## with (ListTools)

\# The following function tests for solutions for specific $k 1=3 a+1, k 2=3 b+1$ and a specific modulus $m$. It returns an empty list if a solution is found, and $[(a, b)]$ otherwise.

```
find_rational_root_test_modulus := proc (a, b, m) local i, g, lc,
    f, d, n, found;
    1c := 3^(6*a+2)-1;
    g := 1c* *^ 3+3^(4*a+2*b+3)*x^2+3^(2*a+4*b+3)*x+3^(6*b+2)-1;
    if igcd(lc, m) = 1 then
        f := modp1(ConvertIn(g, x), m);
        found := false;
        for i from 0 to m-1 while not found do
            if modp1(Eval(f, i), m) = 0 then
                found := true
            end if
        end do;
        if not found then
            return [[a, b]]
        end if
    end if;
    return []
end proc
```

\# The following function returns tuples of integers $(t 1, t 2)$ such that for all $(a, b)$ satisfying $a \equiv t 2\left(\bmod \operatorname{ord}_{m}(3)\right)$ and $b \equiv t 2\left(\bmod \operatorname{ord}_{m}(3)\right)$ there are provably no solutions to our system. Here $o$ is the order of 3 modulo $m$

```
find_rational_root_test_modulus_complete := proc (m, o) local i,
    j, list, revlist, ord, n;
    n := 3;
    list := [];
    for i from o to 2*o-1 do
        for j from 0 to o-1 do
            list := [op(list), op(find_rational_root_test_modulus(i, j,
                    m) )]
        end do
    end do;
    list := map(x }->\mathrm{ map(y }->\mathrm{ y mod o, x), list);
    return sort(list)
end proc
    # Assumes (3,m)=1.
ord_3_modulo_m := proc (m) local ord, n;
    n := 3;
    ord := 1;
    while n mod m <> 1 do n := 3*n;
        ord := ord+1
    end do;
    return ord
end proc
```

\# The following function returns integers for which 3 has order a multiple of $o$, and their orders.

```
ord_3_multiple_o := proc (l, h, o) local m, ord, n, list;
    list := [];
    m := max(1, 4);
    while m <= h do
        n := 3;
        ord := ord_3_modulo_m(m);
```

```
    if ord mod o = 0 then
        list := [op(list), [m, ord]]
        end if ;
        m := m+1;
        if m}\operatorname{mod}3=0\mathrm{ then
            m := m+1
        end if
    end do;
    return list
end proc
```

\# The following function takes as input a modulus and a list of tuples. It is intended for moduli that are a multiple of $o$. It outputs the set of tuples modulo $o$ which occur completely in the list. For example, if $m=10, o=5$ and all of $(1,4),(1,9),(6,4),(6,9)$ occur in the list, then $(1,4)$ will be included in the output.

```
reduce_to_mod_o := proc (m, list, o) local i, j, i2, j2, 12,
    found;
12 := [];
for i from 0 to o-1 do
    for j from 0 to o-1 do
            found := true;
            for i2 from 0 to iquo(m, o)-1 while found do
                for j2 from 0 to iquo(m, o)-1 while found do
                    if not member([ i+o*i2, j+o*j2], list) then
                        found := false
                end if
                end do
            end do;
            if found = true then
                12:= [op(12 ), [i, j]]
            end if
        end do
    end do;
    return 12
```


## end proc

\# The following function takes as input a modulus and a list of tuples. It is intended for moduli that divide $o$. It outputs the set of tuples modulo $m$ which reduce mod o to something occuring in the list. For example, if $m=5, o=10$ and $(1,4)$ occurs in the list, then $(1,4),(1,9),(6,4),(6,9)$ will be included in the output.

```
increase_to_mod_o := proce(m, list, o) local i, j, d, t, 12;
    12 := [];
    d := o/m;
    for t in list do
            for i from 0 to d-1 do
                for j from 0 to d-1 do
                    12:= [op(12 ), [t[1]+m*i, t[2]+m*j]]
            end do
        end do
    end do;
    return 12
end proc
```

\# The following function returns tuples of integers $(t 1, t 2)$ such that for all $(k 1, k 2)$ satisfying $k 1 \equiv t 2(\bmod o)$ and $k 2 \equiv t 2(\bmod o)$ there are provably no solutions to our system.

```
find_all_mod_o := proc (1, h, o) local x, list;
    list := [];
    for x in ord_3_multiple_o(l, h, o) do
        list := [op(list), op(reduce_to_mod_o(x[2],
            find_rational_root_test_modulus_complete(x[1], x[2]), o))]
    end do;
    for x from 2 to o-1 do
        if }0\operatorname{mod}x=0 the
            list := [op(list), op(increase_to_mod_o(x, find_all_mod_o(l
                , h, x), o))]
        end if
    end do;
    return MakeUnique(sort(list))
end proc
```

\# The result of the following command shows that the only possible values of $(a(\bmod 5), b$ $(\bmod 5))$ are $(0,0),(1,1),(2,2),(3,3),(4,4),(1,3),(2,3),(3,1),(3,2),(3,4),(4,3)$. Note in our system we are proving for $k 1=3 a+1, k 2=3 b+1$.

MakeUnique (sort ([op(find_all_mod_o(31, 31, 5)), op(find_all_mod_o (61, 61, 5))])
\# The above command outputs

$$
[[0,1],[0,2],[0,3],[0,4],[1,0],[1,2],[1,4],[2,0],[2,1],[2,4],[3,0],[4,0],[4,1],[4,2]]
$$

\# The following function attempts to prove that our system has no solution for a fixed $k 1, k 2$. It does this by checking modulo primes up to a bound $B$ until it finds no solution.

```
find_rational_root_test_prime := proc (a, b, B) local i, g, f, d,
    lc, found, p;
1c := 3^(6*a+2)-1;
g := 1c*x^ 3+3^(4*a+2*b+3)*x^2+3^(2*a+4*b+3)*x+3^(6*b+2)-1;
d := diff(g, x);
p := 2;
while p <= B do
    if irem(lc, p) <> 0 and 'mod'(Gcd(g, d), p) = 1 then
        f := modp1(ConvertIn(g, x), p);
        found := false;
        for i from 0 to p-1 while not found do
            if modp1(Eval(f, i), p) = 0 then
                found := true
            end if
        end do;
        if not found then
            return p
        end if
    end if;
    p := nextprime(p)
    end do;
    return false
```


## end proc

\# This is the main function. It checks for $k 1>k 2$ for $a>b$ between $l$ and $h$. Here we say $k 1=3 a+1$ and $k 2=3 b+1 . B$ is an upper bound on the primes used in the subfunction. It is optimized to check only when $(a(\bmod 5), b(\bmod 5))$ is one of $(0,0),(1,1),(2,2),(3,3),(4,4)$, $(1,3),(2,3),(3,1),(3,2),(3,4),(4,3)$, which is proven sufficient above. It returns false after no solutions are found.
find_rational_root_test_prime_main_optimized $:=\operatorname{proc}(1, h, B)$
local i, $j$;
for i from iquo(1, 5) to iquo $(h, 5)$ do
for j from 0 to $\mathrm{i}-1$ do
if find_rational_root_test_prime ( $5 * \mathrm{i}, ~ 5 * \mathrm{j}, ~ B)=$ false then
return true, $5 * \mathrm{i}, 5 * \mathrm{j}$
elif find_rational_root_test_prime ( $5 * i+1,5 * j+1, B)=$ false then
return true, $5 * \mathrm{i}+1,5 * \mathrm{j}+1$
elif find_rational_root_test_prime ( $5 * i+1,5 * j+3, B)=$ false then
return true, $5 * \mathrm{i}+1,5 * \mathrm{j}+3$
elif find_rational_root_test_prime $(5 * i+2,5 * j+2, B)=$ false then
return true, $5 * \mathrm{i}+2,5 * \mathrm{j}+2$
elif find_rational_root_test_prime ( $5 * i+2,5 * j+3, B)=$ false then
return true, $5 * \mathrm{i}+2,5 * \mathrm{j}+3$
elif find_rational_root_test_prime ( $5 * \mathrm{i}+3,5 * \mathrm{j}+1, \mathrm{~B})=\mathrm{false}$ then
return true, $5 * \mathrm{i}+3,5 * \mathrm{j}+1$
elif find_rational_root_test_prime $(5 * i+3,5 * j+2, B)=$ false then
return true, $5 * \mathrm{i}+3,5 * \mathrm{j}+2$
elif find_rational_root_test_prime $(5 * i+3,5 * j+3, B)=$ false then
return true, $5 * \mathrm{i}+3,5 * \mathrm{j}+3$

```
            elif find_rational_root_test_prime( 5*i+3, 5*j+4, B) = false
                then
                return true, 5*i+3, 5*j+4
            elif find_rational_root_test_prime( }5*\textrm{i}+4,5*j+3,B)=fals
                then
            return true, 5*i+4, 5*j+3
            elif find_rational_root_test_prime( }5*\textrm{i}+4,5*j+4,B)=fals
                then
                return true, 5*i +4, 5*j+4
            end if
    end do;
    if find_rational_root_test_prime( }5*\textrm{i}+3,5*\textrm{i}+1,\textrm{B})=\mathrm{ false
            then
            return true, 5*i+3, 5*i+1
    elif find_rational_root_test_prime(5*i+3, 5*i+2, B) = false
            then
            return true, 5*i+3, 5*i+2
    elif find_rational_root_test_prime(5*i+4, 5*i+3, B) = false
        then
        return true, 5*i+4, 5*i+3
        end if
    end do;
    return false
end proc
\# The following command returns false.
find_rational_root_test_prime_main_optimized (0, 10000, 500)
\# The following command returns 42548.
numelems(find_all_mod_o (0, 300, 210))
```

