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Modular functions and transcendence questions

Yu. V. Nesterenko

Abstract. We prove results on the transcendence degree of a field generated by numbers connected with the modular function $j(\tau)$. In particular, we show that π and e^{π} are algebraically independent and we prove Bertrand's conjecture on algebraic independence over \mathbb{Q} of the values at algebraic points of a modular function and its derivatives.

Bibliography: 19 items.

§ 1. Statement of results

It is well known that if $\tau \in \mathbb{C}$, Im $\tau > 0$, is an imaginary quadratic irrational number, then the value of the modular function $j(\tau)$ is an algebraic number (see [1], Chapter 5). For example,

$$j(i) = 1728$$
, $j'(i) = 0$, $j(\zeta) = j'(\zeta) = j''(\zeta) = 0$,

where $i^2 = -1$ and $\zeta = e^{2\pi i/3}$.

In 1937, Schneider [2] proved that $j(\tau)$ is transcendental if τ is algebraic but not imaginary quadratic. The purpose of the present article is to study algebraic independence of numbers connected with $j(\tau)$.

In 1916, Ramanujan [3] defined the three functions

$$P(z) = 1 - 24 \sum_{n=1}^{\infty} \sigma_1(n) z^n, \qquad Q(z) = 1 + 240 \sum_{n=1}^{\infty} \sigma_3(n) z^n,$$

$$R(z) = 1 - 504 \sum_{n=1}^{\infty} \sigma_5(n) z^n,$$

where $\sigma_k(n) = \sum_{d|n} d^k$, and proved several identities for these functions. In particular, he showed that they satisfy the differential equations

$$\theta P = \frac{1}{12}(P^2 - Q), \qquad \theta Q = \frac{1}{3}(PQ - R), \qquad \theta R = \frac{1}{2}(PR - Q^2),$$
 (1)

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where $\theta = z \frac{d}{dz}$. In 1969, Mahler proved that the functions P(z), Q(z), and R(z)are algebraically independent over $\mathbb{C}(z)$. (This is an easy consequence of a theorem

The above properties of Ramanujan's functions, together with the fact that their Taylor coefficients are integers that do not increase very rapidly, play a crucial role in the proof of our results.

The functions $E_4(\tau) = Q(e^{2\pi i \tau})$ and $E_6(\tau) = R(e^{2\pi i \tau})$, Im $\tau > 0$, are modular forms of weight 4 and 6, respectively. The function $E_2(\tau) = P(e^{2\pi i \tau})$ also has some modular properties. All of this is essential in deriving the corollaries of the theorem below, which is the basic result of the paper.

Theorem 1. For each $q \in \mathbb{C}$ with 0 < |q| < 1 at least three of the numbers q, P(q), Q(q), R(q) are algebraically independent over \mathbb{Q} .

$$\Delta = \frac{1}{1728}(Q^3 - R^2) = z + \cdots, \qquad J(z) = \frac{Q(z)^3}{\Delta(z)} = \frac{1}{z} + 744 + \sum_{n=1}^{\infty} c(n)z^n.$$

The last series gives the Fourier expansion of the modular function $j(\tau) = J(e^{2\pi i \tau})$. From the differential equations (1) and the definition of J(z) and $\Delta(z)$ it follows

$$\theta \Delta = P \Delta, \qquad J = \frac{Q^3}{\Delta} \,, \qquad \theta J = -\frac{Q^2 R}{\Delta} \,, \qquad \theta^2 J = \frac{-P Q^2 R + 4 Q R^2 + 3 Q^4}{6 \Delta} \,. \label{eq:delta_$$

The last three formulae can be inverted to give

$$P = 6\frac{\theta^2 J}{\theta J} - 4\frac{\theta J}{J} - 3\frac{J}{J - 1728} \,, \qquad Q = \frac{(\theta J)^2}{J(J - 1728)} \,, \qquad R = -\frac{(\theta J)^2}{J^2(J - 1728)} \,.$$

If we take into account that one can have $j(\tau) = 0$, $j(\tau) = 1728$, or $j'(\tau) = 0$ only for values of τ that are equivalent to i or ζ under the modular group, then we obtain the following consequences of the above equalities.

Corollary 1. If $\tau \in \mathbb{C}$ with $\text{Im } \tau > 0$ is not equivalent to i or ζ under the modular group, and we set $q = e^{2\pi i \tau}$, then each of the quadruples

- (1) q, J(q), J'(q), J''(q), (2) q, $j(\tau)$, $\pi^{-1}j'(\tau)$, $\pi^{-2}j''(\tau)$

contains at least three numbers that are algebraically independent over \mathbb{Q} .

Corollary 2. Suppose that q is an algebraic number with 0 < |q| < 1. Then each set

- (1) P(q), Q(q), R(q),
- (2) J(q), $\theta J(q)$, $\theta^2 J(q)$

is algebraically independent over Q. In particular, all of these numbers are transcendental.

The first corollary follows immediately from Theorem 1. To prove Corollary 2, we note that if $q = e^{2\pi i \tau}$, then, by the Gel'fond-Schneider theorem, τ cannot be an algebraic irrationality; hence, it is not equivalent to i or ζ under the modular group. The corollary now follows from Corollary 1.

Corollary 3. Let $\wp(z)$ be the Weierstrass \wp -function with algebraic invariants g_2, g_3 ; let ω_1, ω_2 be its periods, $\operatorname{Im}(\omega_2/\omega_1) \neq 0$, and let η_1 be the quasi-period corresponding to ω_1 . Then the numbers

$$e^{2\pi i(\omega_2/\omega_1)}$$
, $\frac{\omega_1}{\pi}$, $\frac{\eta_1}{\pi}$

are algebraically independent over \mathbb{Q} .

In proving Corollary 3, without loss of generality we may assume that ω_1 and ω_2 are fundamental periods and $\text{Im}(\omega_2/\omega_1) > 0$. Setting $q = e^{2\pi i(\omega_2/\omega_1)}$, $\omega = \omega_1$, and $\eta = \eta_1$, we have (see [1], Chapter 4)

$$P(q) = 3 \cdot \frac{\omega}{\pi} \cdot \frac{\eta}{\pi}, \qquad Q(q) = \frac{3}{4} \left(\frac{\omega}{\pi}\right)^4 g_2, \qquad R(q) = \frac{27}{8} \left(\frac{\omega}{\pi}\right)^6 g_3.$$

These formulae imply that all three numbers P(q), Q(q), R(q) are algebraic over the field $\mathbb{Q}(\omega/\pi, \eta/\pi)$. By Theorem 1, $\mathbb{Q}(q, \omega/\pi, \eta/\pi)$ has transcendence degree 3.

Corollary 4. Let $\wp(z)$ be the Weierstrass \wp -function with algebraic invariants g_2, g_3 and with complex multiplication by the field k. If ω is any period of $\wp(z)$, η is the corresponding quasi-period, and τ is any element of k with $\operatorname{Im} \tau \neq 0$, then each of the sets

$$\{\pi, \ \omega, \ e^{2\pi i \tau}\}, \qquad \{\omega, \ \eta, \ e^{2\pi i \tau}\}$$

is algebraically independent over \mathbb{Q} .

In fact, without loss of generality we may assume that ω is a primitive period, that is, it is not a multiple of any other period of $\wp(z)$. We may further assume that $\tau = \omega_2/\omega$, where $\omega_1 = \omega$ and ω_2 form a basis of the period lattice of $\wp(z)$. In the complex multiplication case the numbers ω_2 and η_2 are algebraic over $\mathbb{Q}(\omega,\eta)$ (see [5], Chapter 3). Using the Legendre relation, we then find that η is algebraic over $\mathbb{Q}(\omega,\eta)$ and π is algebraic over $\mathbb{Q}(\omega,\eta)$. Corollary 4 now follows from Corollary 3.

Corollary 5. Each of the sets

$$\{\pi, e^{\pi}, \Gamma(1/4)\}, \{\pi, e^{\pi\sqrt{3}}, \Gamma(1/3)\}$$

is algebraically independent over \mathbb{Q} . In particular, π and e^{π} are algebraically independent over \mathbb{Q} .

To prove this, we apply Corollary 4 to the elliptic curves with complex multiplication given by the equations

$$y^2 = 4x^3 - 4x, \qquad y^2 = 4x^3 - 4.$$

The first has complex multiplication field $\mathbb{Q}(i)$, while the second has complex multiplication field $\mathbb{Q}(\zeta)$, where $\zeta = e^{2\pi i/3}$. Their periods ω are, respectively,

$$\omega = 2 \int_1^\infty \frac{dx}{(4x^3 - 4x)^{1/2}} = \frac{\Gamma(1/4)^2}{(8\pi)^{1/2}} \quad \text{and} \quad \omega = 2 \int_1^\infty \frac{dx}{(4x^3 - 4)^{1/2}} = \frac{\Gamma(1/3)^3}{\pi \cdot 2^{8/3}}.$$

For any natural number D there exists a Weierstrass \wp -function with algebraic invariants and with complex multiplication field $\mathbb{Q}(\sqrt{-D})$. Thus, we obtain the following corollary.

Corollary 6. For any natural number D the numbers

$$\pi$$
, $e^{\pi\sqrt{D}}$

are algebraically independent over Q.

We now make some remarks of a historical nature.

The transcendence of J(z) for algebraic q, 0 < |q| < 1, was conjectured in 1969 by Mahler [6] and proved in 1995 by Barré-Sirieix, Diaz, Gramain, and Philibert [7]. Their paper provided the impetus for the research that led to the proof of our results.

The algebraic independence over $\mathbb Q$ of the numbers ω_1/π and η_1/π in Corollary 3 and the numbers π and ω in Corollary 4 was first proved by Chudnovskii in 1976 (see [8] and [9]). As a consequence, in the same papers Chudnovskii proved algebraic independence over $\mathbb Q$ of the numbers π and $\Gamma(1/4)$ and the numbers π and $\Gamma(1/3)$ (see Corollary 5). In 1977, Bertrand [10] showed that the first of these results of Chudnovskii is equivalent to the statement that $\theta J(q)$ and $\theta^2 J(q)$ are algebraically independent over $\mathbb Q$, where J(q) is an algebraic number not equal to 0 or 1728. In [10] he conjectured that $J(q), \theta J(q)$, and $\theta^2 J(q)$ are algebraically independent for any algebraic q with 0 < |q| < 1. This conjecture is proved in Corollary 2.

The statement about algebraic independence of π and e^{π} (Corollary 5) is one of the oldest conjectures in the folklore of transcendental number theory. A reference to this conjecture in print can be found in [11].

The method of proof of Theorem 1 can also be used to obtain quantitative results.

Theorem 2. Suppose that $q \in \mathbb{C}$, 0 < |q| < 1, and $\theta_1, \theta_2, \theta_3 \in \mathbb{C}$ are such that the numbers q, P(q), Q(q) and R(q) are all algebraic over the field $\mathbb{Q}(\theta_1, \theta_2, \theta_3)$. Then there exists a constant γ_1 depending only on q and the θ_i such that the following inequality holds for any polynomial $A \in \mathbb{Z}[x_1, x_2, x_3]$, $A \not\equiv 0$:

$$|A(\theta_1, \theta_2, \theta_3)| > \exp(-\gamma_1 t(A)^4 \ln^{24} t(A)),$$

where $t(A) = \ln H(A) + \deg A$ and H(A) is the maximum modulus of the coefficients of A. In particular, this estimate holds for each of the sets

$$\pi, e^{\pi}, \Gamma(1/4)$$
 and $\pi, e^{\pi\sqrt{3}}, \Gamma(1/3)$.

As we already noted, the functions P(z), Q(z) and R(z) are algebraically independent over $\mathbb{C}(z)$. The next theorem, which gives a bound for the measure of their algebraic independence, plays an important role in the proof of Theorem 1.

Theorem 3. Let L_1 and L_2 be integers with $L_1 \ge 1$, $L_2 \ge 1$. Any polynomial $A(z, x_1, x_2, x_3) \in \mathbb{C}[z, x_1, x_2, x_3]$ with $A \not\equiv 0$, $\deg_z A \leqslant L_1$ and $\deg_{x_i} A \leqslant L_2$ satisfies the inequality

ord
$$A(z, P(\dot{z}), Q(z), R(z)) \leq cL_1L_2^3$$
,

where $c = 2 \cdot 10^{45}$.

Here and in the sequel the symbol ord $\varphi(z)$ denotes the multiplicity of the zero of $\varphi(z)$ at the point z=0.

We conclude this section by mentioning that analogues of Theorems 1 and 2 can be proved in the p-adic domain.

§ 2. Reduction of Theorems 1 and 2 to Theorem 3

Lemma 2.1. For all sufficiently large integers N there exists a polynomial A in $\mathbb{Z}[z, x_1, x_2, x_3]$, $A \not\equiv 0$, such that

$$\deg_x A \leqslant N$$
, $\deg_{x_i} A \leqslant N$, $i = 1, 2, 3$, $\ln H(A) \leqslant 85N \ln N$,

where H(A) is the maximum modulus of the coefficients of A, and the function

$$F(z) = A(z, P(z), Q(z), R(z))$$

satisfies the equations

$$F^{(k)}(0) = 0, \qquad k = 0, 1, \dots, \left\lceil \frac{(N+1)^4}{2} \right\rceil - 1.$$

Proof. For $k \ge 1$ we have

$$\sigma_k(n) = \sum_{d|n} d^k \leqslant n^k \sum_{d|n} 1 \leqslant n^{k+1}.$$

Actually, there is a more precise bound for $\sigma_k(n)$ (see [12]), but the trivial estimate is sufficient for our proof. We further have

$$1 + \sum_{n=1}^{\infty} n^k z^n \ll \sum_{n=0}^{\infty} (n+1) \cdots (n+k) z^n = \frac{k!}{(1-z)^{k+1}},$$
 (2)

where \ll is the symbol for one function majorizing another; consequently,

$$P(z) \ll \frac{24 \cdot 2!}{(1-z)^3}, \qquad Q(z) \ll \frac{240 \cdot 4!}{(1-z)^5}, \qquad R(z) \ll \frac{504 \cdot 6!}{(1-z)^7}, \qquad z \ll \frac{1}{1-z}.$$

For any vector $\overline{k} = (k_0, k_1, k_2, k_3), k_i \in \mathbb{Z}, 0 \leqslant k_i \leqslant N, i = 0, 1, 2, 3$, we have

$$z^{k_0} P(z)^{k_1} Q(z)^{k_2} R(z)^{k_3} = \sum_{n=0}^{\infty} d(\overline{k}, n) z^n \ll \frac{c_1^{3N}}{(1-z)^{16N}},$$
 (3)

where $c_1 = 504 \cdot 6!$ and $d(\overline{k}, n) \in \mathbb{Z}$. From (3) and (2) it follows that

$$|d(\overline{k}, n)| \le c_1^{3N} (n + 16N)^{16N} \le (c_2 nN)^{16N} \le (nN)^{17N}, \quad n \ge 1,$$
 (4)

if N is sufficiently large; and $|d(\overline{k},0)| \leq 1$.

Now let

$$A = \sum_{0 \le k_i \le N} a(\overline{k}) z^{k_0} x_1^{k_1} x_2^{k_2} x_3^{k_3},$$

where the set of integers $a(\overline{k})$ is chosen to be a non-trivial solution of the system of homogeneous linear equations

$$\sum_{0\leqslant k_i\leqslant N}d(\overline{k},n)a(\overline{k})=0, \qquad n=0,1,\ldots,\left[\frac{(N+1)^4}{2}\right]-1.$$

Here the number of variables $u = (N+1)^4$ and the number of equations $v = [(N+1)^4/2]$ satisfy the inequality $2v \le u$. Hence, if we use (4) and Siegel's lemma on homogeneous systems of linear equations (see, for example, [11]), we see that the system has a non-trivial integer solution satisfying the inequality

$$\max_{\overline{k}} |a(\overline{k})| \leq (N+1)^4 \left(\frac{(N+1)^4}{2}N\right)^{17N} \leq N^{85N}.$$

Lemma 2.1 is proved.

We introduce the notation

$$M = \operatorname{ord} F(z)$$
.

From Lemma 2.1 and Theorem 3 it follows that

$$\frac{1}{2}N^4 \leqslant M \leqslant cN^4. \tag{5}$$

We further set $r = \min((1 + |q|)/2, 2|q|)$. Then |q| < r < 1.

Lemma 2.2. If N is sufficiently large, then for all $z \in \mathbb{C}$ with $|z| \leq r$ we have

$$|F(z)| \leqslant |z|^M N^{189N}.$$

Proof. Let F(z) have the following Taylor expansion at the origin:

$$F(z) = \sum_{n=M}^{\infty} b_n z^n.$$

Then

$$b_n = \sum_{0 \leqslant k_i \leqslant N} d(\overline{k}, n) a(\overline{k}) \in \mathbb{Z}$$

and, by (4),

$$|b_n| \leqslant \sum_{0 \leqslant k_i \leqslant N} (nN)^{17N} N^{85N} \leqslant n^{17N} N^{103N}, \qquad n \geqslant M.$$

For $|z| \leq r$ we have

$$|F(z)| \leq \sum_{n=M}^{\infty} |b_n| \cdot |z|^n = \sum_{n=0}^{\infty} |b_{n+M}| \cdot |z|^{n+M}$$

$$\leq |z|^M N^{103N} \sum_{n=0}^{\infty} (n+M)^{17N} |z|^n$$

$$\leq |z|^M N^{103N} (M+1)^{17N} \left(1 + \sum_{n=1}^{\infty} n^{17N} |z|^n\right)$$

$$\leq |z|^M N^{103N} (cN^4 + 1)^{17N} \frac{(17N)!}{(1-r)^{17N+1}} \leq |z|^M N^{189N}$$

if N is sufficiently large. Lemma 2.2 is proved.

Lemma 2.3. There exists an integer T, $0 \le T \le \gamma N \ln N$, for which

$$|F^{(T)}(q)| > \left(\frac{1}{2}|q|\right)^{2M},$$

where $\gamma = 190(\ln |r/q|)^{-1}$.

Proof. Suppose that the following inequalities hold:

$$4^{L+1}|F^{(k)}(q)| \le \left(\frac{1}{2}|q|\right)^M, \qquad 0 \le k \le L = [\gamma N \ln N].$$
 (6)

On the circle C: |z| = r we have $z\overline{z} = r^2$ and

$$|r^2 - \overline{q}z| = |r^2 - \overline{q}z| \cdot \left| \frac{\overline{z}}{r} \right| = |r\overline{z} - r\overline{q}| = r|z - q|.$$

Using these relations and Lemma 2.2, we find that the integral

$$I = \frac{1}{2\pi i} \int_C \frac{F(z)}{z^{M+1}} \cdot \left(\frac{r^2 - \overline{q}z}{r(z-q)}\right)^{L+1} dz$$

can be bounded as follows:

$$|I| \leqslant N^{189N}.\tag{7}$$

Meanwhile,

$$I = \mathop{\rm Res}_{z=0} G(z) + \mathop{\rm Res}_{z=q} G(z), \tag{8}$$

where

$$G(z) = \frac{F(z)}{z^{M+1}} \cdot \left(\frac{r^2 - \overline{q}z}{r(z-q)}\right)^{L+1}.$$

The residues at 0 and q are computed to be

$$\operatorname{Res}_{z=0} G(z) = b_M \left(-\frac{r}{q} \right)^{L+1}$$

and

$$\operatorname{Res}_{z=q} G(z) = \frac{1}{L!} \left(\frac{d}{dz} \right)^{L} \left(F(z) \cdot \frac{(r^2 - \overline{q}z)^{L+1}}{z^{M+1}r^{L+1}} \right) \bigg|_{z=q} = \sum_{k=0}^{L} \frac{F^{(L-k)}(q)}{(L-k)!r^{L+1}} \cdot c_k, \quad (9)$$

where

$$c_k = \frac{1}{k!} \left(\frac{d}{dz} \right)^k \left(\frac{(r^2 - \overline{q}z)^{L+1}}{z^{M+1}} \right) \bigg|_{z=q} = \frac{1}{2\pi i} \int_{C_1} \frac{(r^2 - \overline{q}z)^{L+1}}{z^{M+1}(z-q)^{k+1}} \, dz$$

and C_1 is the circle |z-q|=|q|/2. Since on C_1 we have

$$\begin{split} |z|\geqslant |q|-|z-q|\geqslant \frac{1}{2}|q|,\\ |r^2-\overline{q}z|=|r^2-q\overline{q}+\overline{q}(q-z)|\leqslant r^2-|q|^2+\frac{1}{2}|q|^2\leqslant r^2\leqslant 2r|q| \end{split}$$

and

$$\left| \frac{(r^2 - \overline{q}z)^{L+1}}{(z - q)^{k+1}} \right| \leqslant (2r|q|)^{L+1} \left(\frac{1}{2} |q| \right)^{-k-1} \leqslant (4r)^{L+1},$$

it follows that

$$|c_k| \leqslant (4r)^{L+1} \left(\frac{1}{2}|q|\right)^{-M}, \qquad 0 \leqslant k \leqslant L.$$

The last inequality together with (6) and (9) gives us

$$\left| \underset{z=q}{\operatorname{Res}} G(z) \right| \leqslant \sum_{k=0}^{L} \frac{1}{(L-k)!} \leqslant e.$$

Now since $b_M \in \mathbb{Z}$, from (8) and (7) we obtain

$$\left|\frac{r}{q}\right|^{L+1} \leqslant \left|b_M\left(-\frac{r}{q}\right)^{L+1}\right| = \left|\underset{z=0}{\operatorname{Res}}\,G(z)\right| \leqslant |I| + \left|\underset{z=q}{\operatorname{Res}}\,G(z)\right| \leqslant N^{189N} + e.$$

But this inequality is false for $N \geqslant 2$. This contradiction shows that there exists an integer $T \leqslant L$ for which

$$|F^{(T)}(q)| > 4^{-L-1} \left(\frac{1}{2}|q|\right)^M > \left(\frac{1}{2}|q|\right)^{2M}$$

if N is sufficiently large. Lemma 2.3 is proved.

Lemma 2.4. There exists a sequence of polynomials

$$A_N \in \mathbb{Z}[z, x_1, x_2, x_3], \qquad N \geqslant N_0,$$

such that

$$\deg A_N \leqslant 2\gamma N \ln N, \qquad \ln H(A_N) \leqslant 2\gamma N \ln^2 N, \tag{10}$$

$$\exp(-\varkappa_2 N^4) \leqslant \left| A_N(q, P(q), Q(q), R(q)) \right| \leqslant \exp(-\varkappa_1 N^4), \tag{11}$$

where

$$\varkappa_1 = \frac{1}{4} \ln \frac{1}{r}, \qquad \varkappa_2 = 3c \ln \frac{2}{|q|}$$

and c is the constant in Theorem 3.

Proof. Any polynomial $B \in \mathbb{C}[z, x_1, x_2, x_3]$ satisfies the identity

$$\frac{d}{dz}B(z,P(z),Q(z),R(z)) = z^{-1}DB(z,P(z),Q(z),R(z)), \tag{12}$$

where

$$D = z \frac{\partial}{\partial z} + \frac{1}{12} (x_1^2 - x_2) \frac{\partial}{\partial x_1} + \frac{1}{3} (x_1 x_2 - x_3) \frac{\partial}{\partial x_2} + \frac{1}{2} (x_1 x_3 - x_2^2) \frac{\partial}{\partial x_3}$$
(13)

is the operator corresponding to the system of differential equations (1). We now define the sequence of polynomial A_N (for N sufficiently large) by setting

$$A_N(z, x_1, x_2, x_3) = (12z)^T (z^{-1}D)^T A(z, x_1, x_2, x_3).$$

Here A is the polynomial that was constructed in Lemma 2.1 and T is the integer in Lemma 2.3. Using induction, it is easy to prove the identity

$$(z^{-1}D)^T = z^{-T} \prod_{k=0}^{T-1} (D-k), \qquad T \geqslant 1, \tag{14}$$

so that $A_N \in \mathbb{Z}[z, x_1, x_2, x_3]$. From (12) it follows that

$$A_N(z, P(z), Q(z), R(z)) = (12z)^T F^{(T)}(z).$$

Hence, Lemma 2.3 and (5) give the lower bound

$$|A_N(q, P(q), Q(q), R(q))| \geqslant \left(\frac{1}{2}|q|\right)^{3M} \geqslant \exp(-\varkappa_2 N^4).$$

To obtain an upper bound we use the formula

$$F^{(T)}(q) = \frac{T!}{2\pi i} \int_{C_2} \frac{F(z)}{(z-q)^{T+1}} dz,$$

where C_2 is the circle |z-q|=r-|q|. Using the inequality

$$|z| \leqslant |z - q| + |q| = r,$$

along with (5), Lemma 2.2 and Lemma 2.3, we conclude that

$$\left|A_N\left(q,P(q),Q(q),R(q)\right)\right|\leqslant 12^T\cdot T!\cdot (r-|q|)^{-T}r^M\cdot N^{189N}\leqslant \exp(-\varkappa_1N^4).$$

This proves (11).

To prove (10) we use (14). If $B \in \mathbb{C}[z, x_1, x_2, x_3]$ and $B \ll H(1+z+x_1+x_2+x_3)^S$, then for any integer k we have

$$\begin{split} (D+k)B \ll |k|H(1+z+x_1+x_2+x_3)^S + HS(1+z+x_1+x_2+x_3)^{S-1} \\ \times \left(z+(x_1^2+x_2)+(x_1x_2+x_3)+(x_1x_3+x_2^2)\right) \\ \ll |k|H(1+z+x_1+x_2+x_3)^S + HS(1+z+x_1+x_2+x_3)^{S+1} \\ \ll H(S+|k|)(1+z+x_1+x_2+x_3)^{S+1}. \end{split}$$

Hence,

$$12^{T} \prod_{k=0}^{T-1} (D-k)B(z,x_1,x_2,x_3) \ll 12^{T} H(S+2T)^{T} (1+z+x_1+x_2+x_3)^{S+T}$$

and, since

$$A \ll N^{85N} (1 + z + x_1 + x_2 + x_3)^{4N}$$

we find that

$$A_N(z, x_1, x_2, x_3) \ll N^{85N} (48N + 24T)^T (1 + z + x_1 + x_2 + x_3)^{4N+T}$$

But this means that

$$\deg A_N \leqslant 4N + T \leqslant 2\gamma N \ln N,$$

$$H(A_N) \leqslant N^{85N} (48N + 24T)^T \cdot 5^{4N} \leqslant \exp(2\gamma N \ln^2 N),$$

and we have proved (10). Lemma 2.4 is proved.

Lemma 2.5. Let $\overline{\omega} = (\omega_1, \dots, \omega_m) \in \mathbb{C}^m$. Suppose that there exists a sequence of polynomials $A_N \in \mathbb{Z}[x_1, \dots, x_m]$ such that

$$\deg A_N \leqslant \sigma(N), \qquad \ln H(A_N) \leqslant \sigma(N)$$

and

$$\exp(-\varkappa_2\lambda(N)) \leqslant |A_N(\overline{\omega})| \leqslant \exp(-\varkappa_1\lambda(N)),$$

where $\kappa_2 > \kappa_1 > 0$ are constants, and $\sigma(N)$ and $\lambda(N)$ are functions that approach infinity as N increases and satisfy the conditions

$$\lim_{N \to \infty} \frac{\sigma(N+1)}{\sigma(N)} = 1 \quad \text{and} \quad \lim_{N \to \infty} \frac{\lambda(N)}{\sigma(N)^{k+1}} = \infty$$

for some integer k. Then

$$\operatorname{tr} \operatorname{deg}_{\mathbb{Q}} \mathbb{Q}(\omega_1, \ldots, \omega_m) \geqslant k+1.$$

Proof. This lemma is an easy consequence of Theorem 2.11 of Philippon [13], in which we take $\delta(N) = \sigma(N)$, $S(N) = \varkappa_1 \lambda(N)$, $R(N) = 2\varkappa_2 \lambda(N)$ and the sequence of principal ideals I_N is taken to be (A_N) .

Theorem 1 now follows from Lemmas 2.4 and 2.5 with $m=4, k=2, \sigma(N)=2\gamma N \ln^2 N, \lambda(N)=N^4, \varkappa_1$ and \varkappa_2 as in Lemma 2.4, and $\overline{\omega}=\left(q,P(q),Q(q),R(q)\right)$.

To prove Theorem 2, instead of Lemma 2.5 we use a criterion of Ably [14]. In the notation of his paper we take n=4, k=2, $u(x)=\gamma_2x\ln^{-8}x$ with a sufficiently small constant γ_2 and c_0 equal to a sufficiently large constant. For example, we can set $\gamma_2=\varkappa_1(3\gamma)^{-4}$ and $c_0=6^4\varkappa_2/\varkappa_1$. We further take a sequence of principal ideals $I_N=(G_N)$, where $G_N=A_S$ and S is the largest integer for which $2\gamma S\ln^2 S \leqslant N$. Then the criterion in [14] implies the following result.

Theorem 4. Let $q \in \mathbb{C}$, 0 < |q| < 1, and $\overline{\omega} = (1, q, P(q), Q(q), R(q)) \in \mathbb{C}^5$. There exists a positive constant γ_3 depending only on q such that the following inequality holds for any unmixed homogeneous ideal $I \subset \mathbb{Z}[x_0, \ldots, x_4]$ with h(I) = 2 and $I \cap \mathbb{Z} = (0)$:

$$|I(\overline{\omega})| \geqslant \exp(-\gamma_3 t(I)^4 \ln^{24} t(I)).$$

See [15] for the definitions of h(I), t(I), and $|I(\overline{\omega})|$. In § 3 below we give analogous definitions for the polynomial ring $\mathbb{C}[z, x_0, \dots, x_m]$.

Remark. In an entirely similar way one can use Lemma 2.4 and Ably's criterion to prove a bound for $|I(\overline{\omega})|$ for ideals with h(I) = 3. The resulting inequality has the form

$$|I(\overline{\omega})| \geqslant \exp(-\gamma_4 t(I)^2 \ln^8 t(I)).$$

We now show how to derive Theorem 2 from Theorem 4. Under the conditions in Theorem 2 we have

$$\operatorname{tr} \operatorname{deg}_{\mathbb{Q}} \mathbb{Q}(q, P(q), Q(q), R(q)) = 3.$$

We first consider the case when the numbers $\theta_1, \theta_2, \theta_3$ are a subset of q, P(q), Q(q), R(q). Let $E \in \mathbb{Z}[x_0, \ldots, x_4]$ be an irreducible homogeneous polynomial with $E(\overline{\omega}) = 0$, and let B be the homogeneous polynomial defined by setting $B = x_0^{\deg A} A(x_1/x_0, \ldots, x_4/x_0)$. Applying Proposition 1 of [15] to the prime ideal $\mathfrak{p} = (E)$, we find that $|\mathfrak{p}(\overline{\omega})| = 0$. If we now use Theorem 4 and Proposition 3 of [15], then we obtain

$$|\overline{\omega}|^{-\deg B}|B(\overline{\omega})|\geqslant \exp(-96t(E)t(B)-\gamma_3T_1^4\ln^{24}T_1),$$

where $T_1 \leq 32t(\mathfrak{p})t(B)$. Since t(B) = t(A), this gives us the inequality in Theorem 2.

We now consider the general case when the numbers q, P(q), Q(q) and R(q) are algebraic over $\mathbb{Q}(\theta_1, \theta_2, \theta_3)$. Among these four numbers exactly three are algebraically independent over \mathbb{Q} ; let us denote them $\omega_1, \omega_2, \omega_3$. Then all of the θ_i are algebraic over $\mathbb{Q}(\omega_1, \omega_2, \omega_3)$. If $C \in \mathbb{Z}[x_1, x_2, x_3]$ is such that $\xi = C(\omega_1, \omega_2, \omega_3)$ has the property that all of the $\xi \theta_i$ are integral over the ring $\mathbb{Z}[\omega_1, \omega_2, \omega_3]$, then, computing the norm of $\xi^{3 \deg A}(\theta_1, \theta_2, \theta_3)$ over the field $\mathbb{Q}(\omega_1, \omega_2, \omega_3)$, we obtain a polynomial A_1 in $\omega_1, \omega_2, \omega_3$ with integer coefficients. If the inequality proved above is applied to this polynomial, then we find that

$$|\operatorname{Norm}(\xi^{3 \operatorname{deg} A} A(\theta_1, \theta_2, \theta_3))| = |A_1(\omega_1, \omega_2, \omega_3)| \ge \exp(-\gamma_4 t(A_1)^4 \ln^{24} t(A_1)),$$

from which (just as in the proof of Liouville's theorem) we easily derive the required lower bound for $|A(\theta_1, \theta_2, \theta_3)|$. This completes the proof of Theorem 4.

§ 3. Theorem 5 and the derivation of Theorem 3 from Theorem 5

The proof of Theorem 3 is based on a method that uses ideas of commutative algebra and elimination theory. In [16] this method was used to prove a similar result for $m \ge 1$ functions that form a solution of an algebraic system of differential

equations with constant coefficients; and in [17] it was used for $m \ge 1$ functions that satisfy linear differential equations with coefficients in $\mathbb{C}(z)$. We shall apply the general lemmas proved in those papers to the case of the polynomial ring $\mathcal{A} = \mathbb{C}[z, x_0, x_1, x_2, x_3]$ (m = 3). We shall also use lemmas from [18] and [19]. Our notation and definitions will be compatible with those in [15] and [19]; in each case we shall give the appropriate reference.

We consider an unmixed ideal I of \mathcal{A} that is homogeneous in the x_i and has the property that $r = 4 - h(I) \ge 1$. As in [16] and [17], h(I) denotes the height of the ideal I. In § 1 of [18] the ideal I was associated with a certain principal ideal $\overline{I}(r)$ of the ring $\mathbb{C}[z, \overline{u}_1, \ldots, \overline{u}_r]$, where \overline{u}_i is the set of variables u_{i0}, \ldots, u_{i3} . Let F be a generator of $\overline{I}(r)$ (the associated form of the ideal I); it is defined up to a factor in \mathbb{C} . We set

$$N(I) = \deg_{\overline{u}} F, \qquad B(I) = \deg_z F.$$

For each element φ of the field $\mathbb{C}((z))$ of formal power series with complex coefficients, we let $\operatorname{ord} \varphi$ denote the exponent of the first power of z with non-zero coefficient. We let \mathcal{K} denote the algebraic closure of $\mathbb{C}((z))$. The function ord extends uniquely to \mathcal{K} and maps that field to the set $\mathbb{Q} \cup \{\infty\}$.

For
$$\overline{\varphi} = (\varphi_0, \dots, \varphi_3) \in \mathcal{K}^4$$
 we set $|\overline{\varphi}| = \min_j (\operatorname{ord} \varphi_j)$.

We now define ord $I(\overline{\varphi})$; it will play a role similar to ord $P(\varphi)$ for $P \in \mathcal{A}$. As in § 1 of [18], we consider r skew-symmetric matrices $S^{(i)} = \|s_{jk}^{(i)}\|$, where i, j, k vary in the range $1 \leq i \leq r, 0 \leq j \leq 3, 0 \leq k \leq 3$. We suppose that, except for the skew-symmetry relation $s_{jk}^{(i)} + s_{kj}^{(i)} = 0$, there is no other algebraic relation over \mathcal{A} among the $s_{jk}^{(i)}$. Given a polynomial $E \in \mathcal{K}[\overline{u}_1, \dots, \overline{u}_r]$ and a vector $\overline{\varphi} = (\varphi_0, \dots, \varphi_3) \in \mathcal{K}^4$ we let $\varkappa(E)$ denote the polynomial in $s_{jk}^{(i)}, 0 \leq j < k \leq 3, 1 \leq i \leq r$, with coefficients in \mathcal{K} that is obtained by substituting the vector $S^{(i)}\overline{\varphi}$ in place of the variable $\overline{u}_i = (u_{i0}, \dots, u_{i3})$ in E for $i = 1, \dots, r$. The map \varkappa is clearly a homomorphism from the ring $\mathcal{K}[\overline{u}_1, \dots, \overline{u}_r]$ to the polynomial ring $\mathcal{K}[s_{jk}^{(i)}, 0 \leq j < k \leq 3, 1 \leq i \leq r]$. We now define ord $\varkappa(E)$ to be the smallest value that ord takes on the coefficients of the polynomial $\varkappa(E)$.

We set

$$\operatorname{ord} I(\overline{\varphi}) = \operatorname{ord} \varkappa(F) - rN(I)|\overline{\varphi}|,$$

where F is the associated form of the ideal I.

Theorem 5. Let I be an unmixed ideal of A that is homogeneous in the variables x_0, \ldots, x_3 and satisfies $I \cap \mathbb{C}[z] = (0)$ and $r = 4 - h(I) \ge 1$. Then

ord
$$I(\overline{f}) \leq \rho^{2r-1} (B(I)N(I)^{r/(4-r)} + N(I)^{3/(4-r)}),$$

where $\overline{f} = (1, P(z), Q(z), R(z))$ and $\rho = 10^9$.

Theorem 5 is proved by induction on r as r goes from 1 to 3.

Theorem 3 can easily be derived from Theorem 5 using the following lemma.

Lemma 3.1. Suppose that the polynomial $C \in A$ is homogeneous in the variables x_i and I = (C) is the corresponding principal ideal of A. Then $N(I) = \deg_{\overline{x}} C$, $B(I) = \deg_z C$, and for every $\overline{\varphi} \in \mathcal{K}^4$ one has

$$\operatorname{ord} C(\overline{\varphi}) \leqslant \operatorname{ord} I(\overline{\varphi}) + |\overline{\varphi}| \operatorname{deg}_{\overline{x}} C.$$

Proof. See [17], Proposition 1.

Now let A be the polynomial in Theorem 3, and set

$$C = x_0^{\deg A} A\left(z, \frac{x_1}{x_0}, \frac{x_2}{x_0}, \frac{x_3}{x_0}\right),$$

where deg A is the total degree of A in the x_i and I = (C) is the corresponding principal ideal of A. Since h(I) = 1, we see that Theorem 3 follows easily from Lemma 3.1 and Theorem 5 with r = 3. According to Lemma 3.1, we have $B(I) \leq L_1$, $N(I) \leq L_2$, and

$$\operatorname{ord} A(z, P(z), Q(z), R(z)) \leq \operatorname{ord} I(\overline{f}) \leq \rho^5 (L_1 + 1) L_2^3 \leq c L_1 L_2^3$$

This completes the derivation of Theorem 3 from Theorem 5.

\S 4. Ideals invariant with respect to the operator D

In § 2 (see (13)) we defined the differential operator D acting on the ring $\mathfrak{R} = \mathbb{C}[z, x_1, x_2, x_3]$. Recall that for any $G(z, x_1, x_2, x_3)$ in this ring we have the identity

$$\theta G(z, P(z), Q(z), R(z)) = DG(z, P(z), Q(z), R(z)). \tag{15}$$

If an ideal I of the ring \mathfrak{R} or $\mathbb{C}[x_1, x_2, x_3]$ has the property that $DI \subset I$, then we say that I is D-invariant.

Lemma 4.1. There exist only two D-invariant principal prime ideals of \Re , namely, the ideals generated by z and by $\Delta = x_2^3 - x_3^2$.

Proof. We have Dz = z and

$$D\Delta = x_2^2(x_1x_2 - x_3) - x_3(x_1x_3 - x_2^2) = x_1\Delta.$$

Thus, the principal ideal generated by z or by Δ is D-invariant.

Now suppose that $A \in \mathfrak{R}$ is any irreducible polynomial with the property that $A \mid DA$. In other words, we have DA = BA for some $B \in \mathfrak{R}$. Given a polynomial $F(z, x_1, x_2, x_3) \in \mathfrak{R}$, we define its weight $\omega(F)$ as follows:

$$\omega(F) = \deg_t F(z, tx_1, t^2x_2, t^3x_3).$$

Then

- (1) if all monomials of F have the same weight, then either DF = 0 or else $\omega(DF) \geqslant \omega(F)$;
- (2) $\omega(DF) \leq \omega(F) + 1$ for any $F \in \Re$;
- (3) $\omega(FG) = \omega(F) + \omega(G)$ for any two polynomials $F, G \in \mathfrak{R}$.

These properties follow trivially from the definition of the weight and the definition of D. Using (2) and (3) we find that the relation DA = BA implies that

$$\omega(B) + \omega(A) = \omega(DA) \le \omega(A) + 1,$$

so that $\omega(B) \leq 1$. This inequality means that

$$DA = (ax_1 + b)A, (16)$$

where $a, b \in \mathbb{C}[z]$. If we compare the degrees in z of the polynomials in (16), we conclude that $a, b \in \mathbb{C}$.

Let C denote the sum of the monomials of A that have minimal weight. If we compare the sum of the monomials of weight $\omega(C)$ on both sides of (16) and use property (1), then we find that

$$z\frac{\partial C}{\partial z} = bC.$$

In this equality we compare the coefficients of the highest power of z in $C \in \mathfrak{R}$; we conclude that $b \in \mathbb{Z}$.

Now let

$$A(z, P(z), Q(z), R(z)) = c_m z^m + c_{m+1} z^{m+1} + \cdots, \quad c_m \neq 0, \qquad m \geqslant 0.$$

If we apply the operator θ to both sides of this equation and use (15), (16) and the definition of P(z), we find that

$$mc_mz^m+\cdots=(a(1-24z+\cdots)+b)(c_mz^m+\cdots),$$

and hence a + b = m. We have thereby shown that a and b are integers.

The differential operator D can be extended in a natural way to the field $\mathbb{C}(z,x_1,x_2,x_3)$. It is easy to verify that the identify (15) remains valid for any element F of this field. By (15), the relation DF=0 for $F\in\mathbb{C}(z,x_1,x_2,x_3)$ means that

$$\theta F(z, P(z), Q(z), R(z)) = 0,$$

and hence $F(z, P(z), Q(z), R(z)) = c \in \mathbb{C}$. But since the functions P, Q and R are algebraically independent over $\mathbb{C}(z)$, this implies that $F = c \in \mathbb{C}$; thus, the field of constants of the operator D in $\mathbb{C}(z, x_1, x_2, x_3)$ is simply \mathbb{C} .

We easily verify that

$$D(\Delta^{-a}z^{-b}) = -(ax_1 + b)\Delta^{-a}z^{-b},$$

and so, using (16), we obtain

$$D(A \cdot \Delta^{-a} z^{-b}) = 0,$$

that is, by what was just proved, we have $A \cdot \Delta^{-a} z^{-b} = c$ and $A = c \Delta^a z^b$ for some constant $c \in \mathbb{C}$. Since A is an irreducible polynomial, only two cases are possible: a = 1, b = 0, and a = 0, b = 1. This completes the proof of Lemma 4.1.

Remark. What we have actually proved is that every D-invariant principal ideal of \mathfrak{R} has the form $(\Delta^a z^b)$, where a and b are non-negative integers.

We next determine all *D*-invariant prime ideals of the ring $\mathbb{C}[x_1, x_2, x_3]$.

Corollary. There is one and only one D-invariant principal prime ideal in $\mathbb{C}[x_1, x_2, x_3]$. It is generated by Δ .

In fact, if A is an irreducible polynomial that generates a D-invariant principal prime ideal in $\mathbb{C}[x_1, x_2, x_3]$, then the principal ideal of \mathfrak{R} generated by A is also a D-invariant prime ideal. Lemma 4.1 then implies that $A = \Delta$.

Lemma 4.2. All of the D-invariant prime ideals of $\mathbb{C}[x_1, x_2, x_3]$ of dimension 0 have the form $(x_1 - c, x_2 - c^2, x_3 - c^3)$, where $c \in \mathbb{C}$.

Proof. Every prime ideal \mathfrak{p} of $\mathbb{C}[x_1, x_2, x_3]$ of dimension 0 has the form $(x_1 - c_1, x_2 - c_2, x_3 - c_3)$ with $c_i \in \mathbb{C}$. Suppose that $D\mathfrak{p} \subset \mathfrak{p}$. Applying the operator D to the polynomials $x_1 - c_1$, we find that the three polynomials

$$x_1^2 - x_2, \quad x_1 x_2 - x_3, \quad x_1 x_3 - x_2^2$$
 (17)

belong to \mathfrak{p} . Since they all must vanish at the point (c_1, c_2, c_3) , we conclude that $c_2 = c_1^2$ and $c_3 = c_1c_2 = c_1^3$, that is, $\mathfrak{p} = (x_1 - c, x_2 - c^2, x_3 - c^3)$ with $c = c_1$. Conversely, every ideal of the form $\mathfrak{p} = (x_1 - c, x_2 - c^2, x_3 - c^3)$ contains the polynomials (17) and hence is *D*-invariant.

Lemma 4.3. The system of differential equations

$$(x^{2} - f)f' = 4(xf - g),$$

$$(x^{2} - f)g' = 6(xg - f^{2})$$
(18)

has a unique solution in algebraic functions f(x), g(x) with f(1) = 1, g(1) = 1, namely, $f = x^2$, $g = x^3$.

Proof. The functions $f=x^2$, $g=x^3$ form a solution of the system (18). Hence, in what follows we may suppose that $f \not\equiv x^2$. We set $u=x^2-f\not\equiv 0,\ v=xf-g$. Then

$$xg - f^2 = x^2u - u^2 - xv$$

and

$$uu' = 2ux - uf' = 2ux - 4v,$$

$$uv' = uf + 4xv - 6x^2u + 6u^2 + 6xv = 5u^2 - 5x^2u + 10xv = 5u^2 - \frac{5}{2}xuu'.$$

We obtain the system of differential equations

$$uu' = 2xu - 4v, 2v' = 10u - 5xu'.$$
 (19)

Since u(x) and v(x) are algebraic functions, there exists a natural number e and a parametrization (here we are considering parametrizations of branches at x = 1; recall that u(1) = v(1) = 0)

$$x = 1 + t^e$$
, $u = \sum_{k=\lambda}^{\infty} a_k t^k$, $v = \sum_{k=\mu}^{\infty} b_k t^k$,

where $\lambda \geqslant 1$, $\mu \geqslant 1$, and $a_{\lambda}b_{\mu} \neq 0$. We may assume that e has been chosen to be as small as possible.

The expansions of the functions in (19) have the following initial terms:

$$xu = a_{\lambda}t^{\lambda} + \cdots, \qquad xu' = \frac{\lambda}{e}a_{\lambda}t^{\lambda-e} + \cdots,$$
 $uu' = \frac{\lambda}{e}a_{\lambda}^{2}t^{2\lambda-e} + \cdots, \qquad v' = \frac{\mu}{e}b_{\mu}t^{\mu-e} + \cdots.$

If we substitute these expansions into the second equation in (19), then we obtain

$$2\frac{\mu}{e}b_{\mu}t^{\mu-e} + \dots = 10a_{\lambda}t^{\lambda} + \dots - 5\frac{\lambda}{e}a_{\lambda}t^{\lambda-e} + \dots$$
 (20)

Comparing exponents of the smallest powers of t on the left and right, we conclude that $\lambda = \mu$. If we now substitute these expansions into the first equation in (19), then we obtain

$$\frac{\lambda}{e}a_{\lambda}^{2}t^{2\lambda-e} + \dots = 2a_{\lambda}t^{\lambda} + \dots - 4b_{\lambda}t^{\lambda} - \dots$$
 (21)

This implies that the inequality $2\lambda - e \geqslant \lambda$ must hold, and hence $\lambda \geqslant e$.

We next compare the coefficients of $t^{\lambda-e}$ on the left and right in (20). We obtain:

$$2b_{\lambda} = -5a_{\lambda}.\tag{22}$$

If $\lambda > e$, then a comparison of the coefficients of t^{λ} in (21) gives

$$2a_{\lambda}-4b_{\lambda}=0.$$

which, combined with (22), contradicts the assumption that $a_{\lambda} \neq 0$. Hence $\lambda = e$, and from (21) we find that

$$a_e^2 = 2a_e - 4b_e.$$

By (22), this means that $a_e^2 = 12a_e$. Since $a_e \neq 0$, we conclude that $a_e = 12$ and $b_e = -30$.

Now suppose that $e \ge 2$. Let r be the smallest number for which the conditions $a_r \ne 0$ and $r \ne 0 \pmod{e}$ hold. This number exists because e was chosen to be minimal (otherwise, the first equation in (19) would imply that all of the non-zero coefficients b_k have index divisible by e). If we now compare the coefficients of t^r on the left and right of the first equation in (19), then we find that

$$\frac{r+e}{e}a_ra_e=2a_r-4b_r$$

or, since $a_e = 12$,

$$12\frac{r+e}{e}a_r = 2a_r - 4b_r. (23)$$

By making the same substitution in the second equation in (19) and comparing the coefficients of t^{r-e} , we find that

$$2\frac{r}{e}b_r = -5\frac{r}{e}a_r$$

or

$$2b_r + 5a_r = 0.$$

The last equation along with (23) give us $ra_r = 0$; since r > e, this means that $a_r = 0$, which is impossible. Thus, e = 1, the functions u(x) and v(x) are single-valued in a neighbourhood of the point x = 1, and

$$u(1) = v(1) = 0,$$
 $u'(1) = 12,$ $v'(1) = -30.$ (24)

We now prove that all the derivatives of u(x) and v(x) are uniquely determined, that is, there exists a unique solution of the system of differential equations (19) that is analytic in a neighbourhood of x = 1 and satisfies the conditions in (24). Let $k \ge 2$. If we differentiate the first equation in (19) k times, then we find that the function

$$uu^{(k+1)} + (k+1)u'u^{(k)} - 2xu^{(k)} + 4v^{(k)}$$

can be expressed as a polynomial in x and $u, u', \ldots, u^{(k-1)}$. Taking (24) into account, we then find that the quantity

$$(6k+5)u^{(k)}(1) + 2v^{(k)}(1)$$

is uniquely determined by $u(1), u'(1), \ldots, u^{(k-1)}(1)$. In exactly the same way, if we differentiate the second equation in (19) k-1 times, we find that the quantity

$$5u^{(k)}(1) + 2v^{(k)}(1)$$

can be expressed uniquely in terms of $u^{(j)}(1)$ and $v^{(j)}(1)$, $0 \le j \le k-1$. But then the derivatives $u^{(k)}(1)$ and $v^{(k)}(1)$ are also uniquely determined. This proves uniqueness of a solution to (19) that is analytic in a neighbourhood of x = 1 and satisfies initial conditions (24).

In some neighbourhood of the point x = 1 the equation x = P(z) determines z uniquely as an analytic function of x that vanishes at x = 1. We set

$$F(x) = Q(P^{-1}(x)), \qquad G(x) = R(P^{-1}(x)).$$

Then F(x) and G(x) are analytic functions in a neighbourhood of x = 1 that, as we easily verify, satisfy the system of differential equations (18) and the initial conditions F(1) = 1, G(1) = 1. In addition,

$$z = -\frac{1}{24}(x-1) + \cdots$$
, $F(x) = 1 - 10(x-1) + \cdots$, $G(x) = 1 + 21(x-1) + \cdots$.

But then it is easy to see that the functions $U(x) = x^2 - F(x)$ and V(x) = xF(x) - G(x) satisfy the system of differential equations (19) and the initial conditions (24). By the uniqueness proved above, we conclude that u(x) = U(x), v(x) = V(x), and hence f(x) = F(x), g(x) = G(x). Thus, F(x) and G(x) are algebraic functions. But if we substitute x = P(z) into the identity A(x, F(x)) = 0, $A(x,y) \in \mathbb{C}[x,y]$, then it becomes an algebraic relation A(P(z),Q(z)) = 0 between the functions P(z) and Q(z). Since P(z), Q(z) and R(z) are algebraically independent over $\mathbb{C}(z)$, such an algebraic relation is impossible. This contradiction completes the proof of Lemma 4.3.

Proposition 1. All non-trivial D-invariant prime ideals of $\mathbb{C}[x_1, x_2, x_3]$ having a zero at the point (1,1,1) are given by the following list:

- $\begin{array}{ll} (1) \ \ (x_1-1,x_2-1,x_3-1); \\ (2) \ \ (x_1^2-x_2,x_1^3-x_3); \\ (3) \ \ (x_2^3-x_3^2). \end{array}$

They all contain the polynomial $\Delta = x_2^3 - x_3^2$.

Proof. Let \mathfrak{p} be a *D*-invariant prime ideal of $\mathbb{C}[x_1, x_2, x_3]$. If dim $\mathfrak{p} = 0$, then \mathfrak{p} has the form in (1), by Lemma 4.2 (in this case c = 1, by assumption).

If dim p = 2, that is, if p is a principal ideal, then its extension p^e to \Re is also a D-invariant principal prime ideal. Since $z \notin \mathfrak{p}^e$, it follows from Lemma 4.1 that $\mathfrak{p}^e = (\Delta)$. But then $\mathfrak{p} = (\Delta)$, that is, we have case (3).

Now suppose that dim $\mathfrak{p}=1$. If $\mathfrak{p}\cap\mathbb{C}[x_1]\neq (0)$, then, since the ideal \mathfrak{p} is prime, we have $x_1 - c \in \mathfrak{p}$ for some constant c. But then $x_1^2 - x_2 = 12D(x_1 - c) \in \mathfrak{p}$, and so $x_2 - c^2 \in \mathfrak{p}$. Also, $x_1x_2 - x_3 = 3D(x_2 - c^2) \in \mathfrak{p}$ and $x_3 - c^3 \in \mathfrak{p}$. However, the inclusion $(x_1 - c, x_2 - c^2, x_3 - c^3) \subset \mathfrak{p}$ is impossible, since dim $\mathfrak{p} = 1$. Thus, $\mathfrak{p} \cap \mathbb{C}[x_1] = (0)$; hence, $\mathfrak{p} \cap \mathbb{C}[x_1, x_2] \neq (0)$ and $\mathfrak{p} \cap \mathbb{C}[x_1, x_3] \neq (0)$. Let $A(x_1, x_2)$ and $B(x_1, x_3)$ denote irreducible polynomials of $\mathbb{C}[x_1, x_2, x_3]$ that lie in \mathfrak{p} . By assumption, we have A(1,1) = B(1,1) = 0; hence, there exist algebraic functions y = y(x), z = z(x) such that

$$A(x,y(x)) = 0, \quad B(x,z(x)) = 0, \quad y(1) = z(1) = 1,$$
 (25)

and the triple of functions (x, y(x), z(x)) is a zero of the ideal \mathfrak{p} . If we differentiate the first equation in (25) with respect to x, then we obtain

$$\frac{\partial A}{\partial x}(x,y(x)) + \frac{\partial A}{\partial y}(x,y(x))y'(x) = 0.$$
 (26)

If we then take into account that the triple (x, y(x), z(x)) is a zero of the ideal p and, in particular, of the polynomial $DA \in \mathfrak{p}$, then we conclude that

$$\frac{1}{12}(x^2 - y(x))\frac{\partial A}{\partial x}(x, y(x)) + \frac{1}{3}(xy(x) - z(x))\frac{\partial A}{\partial y}(x, y(x)) = 0.$$
 (27)

Multiplying (26) by $\frac{1}{12}(x^2 - y(x))$, subtracting the resulting expression from (27), and using the fact that

$$\frac{\partial A}{\partial y}\big(x,y(x)\big)\not\equiv 0,$$

we find that the functions y(x), z(x) satisfy the first differential equation in (18). Similarly, if we differentiate the second equation in (25), then we find that y(x), z(x)also satisfy the second differential equation in (18). By Lemma 4.3, we now know that $y(x) = x^2$ and $z(x) = x^3$, that is, $A(x_1, x_2) = x_1^2 - x_2$, $B(x_1, x_3) = x_1^3 - x_3$. This means that p contains the prime ideal in part (2) of Proposition 1. Since these two ideals have the same dimension, we conclude that they are the same ideal. Proposition 1 is proved.

Proposition 2. Every D-invariant prime ideal of $\mathbb{C}[z, x_1, x_2, x_3]$ with a zero at (0, 1, 1, 1) contains either the polynomial z, or else the polynomial $\Delta = x_2^3 - x_3^2$.

Proof. If \mathfrak{p} is a *D*-invariant prime ideal of \mathfrak{R} , then $\mathfrak{q} = \mathfrak{p} \cap \mathbb{C}[x_1, x_2, x_3]$ is a *D*-invariant prime ideal of $\mathbb{C}[x_1, x_2, x_3]$. If $\mathfrak{q} \neq (0)$, then, by Proposition 1, we have $\Delta \in \mathfrak{q} \subset \mathfrak{p}$. If $\mathfrak{q} = (0)$, then \mathfrak{p} is a principal ideal, in which case, by Lemma 4.1, $z \in \mathfrak{p}$. Proposition 2 is proved.

§ 5. Auxiliary lemmas

The proof of the next lemma can be found in [17]; it is part of the proof of Theorem 2 of that paper (see pp. 44–45, and also [16], pp. 267–268). For convenience we shall give the proof of Lemma 5.1 in the general case when $\mathcal{A} = \mathbb{C}[z, x_0, \ldots, x_m]$ with $m \geq 1$. We shall say that an ideal of \mathcal{A} is homogeneous if it is homogeneous in the variables x_0, \ldots, x_m .

Lemma 5.1. Let \mathfrak{p} be a homogeneous prime ideal in A with $\mathfrak{p} \cap \mathbb{C}[z] = (0)$. Suppose that the vector $\overline{\varphi} \in \mathcal{K}^{m+1}$ is not a zero of the ideal \mathfrak{p} , $|\overline{\varphi}| = 0$, C and A are homogeneous polynomials in A with $A \in \mathfrak{p}$ and $C \notin \mathfrak{p}$, and the following inequality holds for some $s \geqslant 0$:

$$\operatorname{ord} C(\overline{\varphi}) \geqslant \operatorname{ord} A(\overline{\varphi}) - s.$$

If $r = m + 1 - h(I) \ge 2$, then there exists an unmixed homogeneous ideal $J \subset A$ such that its set of zeros in projective space over the algebraic closure of $\mathbb{C}(z)$ is the same as the set of zeros of the ideal (\mathfrak{p}, C) and, in addition,

- (1) $N(J) \leqslant N(\mathfrak{p}) \deg_{\overline{x}} C$,
- (2) $B(J) \leq B(\mathfrak{p}) \deg_{\overline{x}} C + N(\mathfrak{p}) \deg_z C$,
- (3) $\operatorname{ord} \mathfrak{p}(\overline{\varphi}) \leq \operatorname{ord} J(\overline{\varphi}) + B(\mathfrak{p}) \operatorname{deg}_{\overline{\varphi}} C + N(\mathfrak{p}) (\operatorname{deg}_{z} C + s).$

If r=1, then (3) still holds if we formally set ord $J(\overline{\varphi})=0$ in this case.

Proof. Without loss of generality we may assume that $x_0 \notin \mathfrak{p}$. From Lemmas 2 and 3 of [17] we know that there exists a finite normal extension \mathbb{K} of the field $\mathbb{C}(z,\overline{u}_1,\ldots,\overline{u}_{r-1})$ such that the associated form F of \mathfrak{p} splits over \mathbb{K} into a product of linear forms

$$F = a \prod_{j=1}^{N(\mathfrak{p})} (u_{r0} + \alpha_1^{(j)} u_{r1} + \dots + \alpha_m^{(j)} u_{rm}),$$

where $a \in \mathbb{C}[z, \overline{u}_1, \dots, \overline{u}_{r-1}]$ and $\alpha_i^{(j)} \in \mathbb{K}$, and where each of the points $(1 : \alpha_1^{(j)} : \dots : a_m^{(j)}) \in \mathbb{P}_{\mathbb{K}}^m$ is a common zero of \mathfrak{p} . Moreover, if

$$G = a^{\deg_{\overline{x}} Q} \prod_{j=1}^{N(p)} C(1, \alpha_1^{(j)}, \dots, \alpha_m^{(j)}),$$
 (28)

then $G \in \mathbb{C}[z, \overline{u}_1, \dots, \overline{u}_{r-1}], G \not\equiv 0$, and

$$\deg_z G \leqslant B(\mathfrak{p}) \deg_{\overline{x}} C + N(\mathfrak{p}) \deg_z C. \tag{29}$$

Furthermore, if $r \geqslant 2$, then

$$\deg_{\overline{u}}, G \leqslant N(\mathfrak{p}) \deg_{\overline{x}} C, \tag{30}$$

and there exists an unmixed homogeneous ideal $J \subset A$ with h(J) = m - r + 2, with associated form $w^{-1}G$ for some $w \in \mathbb{C}[z]$, and with the same set of zeros as (\mathfrak{p}, C) in projective space over the algebraic closure of $\mathbb{C}(z)$. For $r \geq 2$ the inequalities (29) and (30) imply (1) and (2) of Lemma 5.1.

It remains to prove (3). According to Lemma 5 of [17] and the remark following the proof of Lemma 6 of [17], there exists a ring homomorphism

$$\tau \colon \mathbb{C}\big[z,\overline{u}_1,\ldots,\overline{u}_{r-1},a^{-1},\alpha_1^{(1)},\ldots,\alpha_m^{(N(\mathfrak{p}))}\big] \to \mathcal{K}$$

over $\mathbb{C}[z]$ such that

$$\operatorname{ord} \tau(a) = \operatorname{ord} \varkappa(a), \qquad \operatorname{ord} \tau(G) = \operatorname{ord} \varkappa(G)$$
 (31)

(the map \varkappa was defined in §3) and, in addition, if we set $\beta_i^{(j)} = \tau(\alpha_i^{(j)})$, then the vectors $\overline{\beta}_j = (1, \beta_1^{(j)}, \dots, \beta_m^{(j)})$, $j = 1, \dots, N(\mathfrak{p})$, are zeros of \mathfrak{p} and satisfy the inequalities

$$\operatorname{ord} \varkappa(a) + \sum_{j=1}^{N(\mathfrak{p})} |\overline{\beta}_{j}| \geqslant (r-1)N(\mathfrak{p})|\overline{\varphi}| \geqslant 0, \tag{32}$$

$$\operatorname{ord} \varkappa(a) + \sum_{j=1}^{N(\mathfrak{p})} \left(\|\overline{\varphi} - \overline{\beta}_j\| + |\overline{\beta}_j| \right) \geqslant \operatorname{ord} \mathfrak{p}(\overline{\varphi}) + (r-1)N(\mathfrak{p})|\overline{\varphi}| \geqslant \operatorname{ord} \mathfrak{p}(\overline{\varphi}).$$
(33)

Here and in the sequel, for two vectors $\overline{\varphi} = (\varphi_0, \dots, \varphi_m) \in \mathcal{K}^{m+1}$ and $\overline{\psi} = (\psi_0, \dots, \psi_m) \in \mathcal{K}^{m+1}$ we use the notation

$$||\overline{\varphi} - \overline{\psi}|| = \min_{0 \le i < j \le m} \operatorname{ord}(\varphi_i \psi_j - \varphi_j \psi_i) - |\overline{\varphi}| - |\overline{\psi}|.$$

We note that always $\|\overline{\varphi} - \overline{\psi}\| \geqslant 0$.

Let j be a fixed index, $1 \le j \le N(\mathfrak{p})$. We choose n so that $|\overline{\beta}_j| = \operatorname{ord} \beta_n^{(j)}$. Since $A(\overline{\beta}_j) = 0$, it follows from Lemma 1 of [17] applied with V = A and $W = x_n^{\deg_{\overline{x}} R}$ that

ord
$$A(\overline{\varphi}) \geqslant ||\overline{\varphi} - \overline{\beta}_i||$$
.

By assumption, we now know that

$$\operatorname{ord} C(\overline{\varphi}) \geqslant \|\overline{\varphi} - \overline{\beta}_j\| - s. \tag{34}$$

If we again apply Lemma 1 of [17], this time with V=C and $W=x_i^{\deg_{\overline{x}}Q}$, where i is chosen so that ord $\varphi_i=|\overline{\varphi}|=0$, then we find that

$$\operatorname{ord}\left(C(\overline{\varphi})(\beta_{i}^{(j)})^{\deg_{\overline{x}}C} - C(\overline{\beta}_{j})\varphi_{i}^{\deg_{\overline{x}}C}\right) \geqslant \|\overline{\varphi} - \overline{\beta}_{j}\| + |\overline{\beta}_{j}| \deg_{\overline{x}}C. \tag{35}$$

Furthermore, from the identity

$$C(\overline{\beta}_j)\varphi_i^{\deg_{\overline{x}}C} = C(\overline{\varphi})(\beta_i^{(j)})^{\deg_{\overline{x}}C} - \left(C(\overline{\varphi})(\beta_i^{(j)})^{\deg_{\overline{x}}C} - C(\overline{\beta}_j)\varphi_i^{\deg_{\overline{x}}C}\right)$$

and the inequalities (34) and (35) we conclude that

$$\operatorname{ord} C(\overline{\beta}_{i}) \geqslant \|\overline{\varphi} - \overline{\beta}_{i}\| + |\overline{\beta}_{i}| \operatorname{deg}_{\overline{x}} C - s. \tag{36}$$

Since, by (28), we have

$$\tau(G) = \tau(a)^{\deg_{\overline{x}}C} \prod_{j=1}^{N(\mathfrak{p})} C(1, \beta_1^{(j)}, \dots, \beta_m^{(j)}),$$

it follows from (31) and (36), along with (32) and (33), that

$$\operatorname{ord} \varkappa(G) = \operatorname{ord} \tau(G) = \operatorname{deg}_{\overline{x}} C \cdot \operatorname{ord} \tau(a) + \sum_{j=1}^{N(\mathfrak{p})} \operatorname{ord} C(\overline{\beta}_{j})$$

$$\geqslant \operatorname{deg}_{\overline{x}} C \cdot \operatorname{ord} \varkappa(a) + \sum_{j=1}^{N(\mathfrak{p})} \left(\|\overline{\varphi} - \overline{\beta}_{j}\| + |\overline{\beta}_{j}| \operatorname{deg}_{\overline{x}} C - s \right)$$

$$\geqslant \operatorname{ord} \varkappa(a) + \sum_{j=1}^{N(\mathfrak{p})} \left(\|\overline{\varphi} - \overline{\beta}_{j}\| + |\overline{\beta}_{j}| \right) - sN(\mathfrak{p})$$

$$\geqslant \operatorname{ord} \mathfrak{p}(\overline{\varphi}) - sN(\mathfrak{p}). \tag{37}$$

Suppose that r=1. Then $G \in \mathbb{C}[z]$, and from (37) and (29) we obtain

$$\operatorname{ord}\mathfrak{p}(\overline{\varphi})\leqslant\operatorname{ord}G+sN(\mathfrak{p})\leqslant\deg_zG+sN(\mathfrak{p})\leqslant B(\mathfrak{p})\deg_{\overline{x}}C+N(\mathfrak{p})(\deg_zC+s).$$

We have thus proved the required inequality in the case r = 1.

Now suppose that $r \ge 2$. Using (37) and (29), we obtain

$$\begin{split} \operatorname{ord} \mathfrak{p}(\overline{\varphi}) &\leqslant \operatorname{ord} \varkappa(G) + sN(\mathfrak{p}) = \operatorname{ord} w + \operatorname{ord} J(\overline{\varphi}) + sN(\mathfrak{p}) \\ &\leqslant \operatorname{ord} J(\overline{\varphi}) + sN(\mathfrak{p}) + \operatorname{deg}_z G \\ &\leqslant \operatorname{ord} J(\overline{\varphi}) + B(\mathfrak{p}) \operatorname{deg}_{\overline{x}} C + N(\mathfrak{p}) (\operatorname{deg}_z C + s). \end{split}$$

This proves the required inequality for $r \ge 2$.

We shall want to apply this lemma with m=3, s=0 and r=1 or 2. We shall also refer to the lemmas in [17] in the case m=3. We now return to our earlier notation $\mathcal{A}=\mathbb{C}[z,x_0,x_1,x_2,x_3]$.

Lemma 5.2. Let $\mathfrak{p} \subset \mathcal{A}$ be a homogeneous prime ideal with $\mathfrak{p} \cap \mathbb{C}[z] = (0)$ and $r = 4 - h(\mathfrak{p}) \geqslant 1$ and let ν and μ be non-negative integers satisfying the inequalities

$$\nu^{4-r} \geqslant \lambda N(\mathfrak{p}), \qquad (\mu+1)\nu^{3-r} \geqslant \lambda B(\mathfrak{p}),$$
 (38)

where $\lambda = 18^3 = 5832$. There exists a polynomial $E \in \mathfrak{p}$ that is homogeneous in \overline{x} and satisfies the inequalities

$$\deg_z E \leqslant \mu$$
, $\deg_{\overline{x}} E \leqslant \nu$.

Proof. See [19], Corollary 2 of Theorem 2 with m = 3.

We note that the inequalities in (38) hold if, for example,

$$\nu = 1 + \left[\lambda N(\mathfrak{p})^{1/(4-r)}\right], \qquad \mu = \left[\lambda B(\mathfrak{p})N(\mathfrak{p})^{-(3-r)/(4-r)}\right], \tag{39}$$

where $[\cdot]$ denotes the integer part.

We now define a homogeneous analogue of the operator D in (13) for the ring A:

$$T = z \frac{\partial}{\partial z} + \frac{1}{12} (x_1^2 - x_0 x_2) \frac{\partial}{\partial x_1} + \frac{1}{3} (x_1 x_2 - x_0 x_3) \frac{\partial}{\partial x_2} + \frac{1}{2} (x_1 x_3 - x_2^2) \frac{\partial}{\partial x_3}.$$

Then $Tx_0 = 0$.

Lemma 5.3. If p is a homogeneous prime ideal of A with

$$\operatorname{ord}\mathfrak{p}(\overline{f}) > B(\mathfrak{p}) + 3N(\mathfrak{p}),$$

then there do not exist homogeneous prime ideals $q \subset p$, $q \neq (0)$, with $Tq \subset q$.

In this lemma, as in Theorem 3, we use the notation $\overline{f} = (1, P(z), Q(z), R(z))$.

Proof. We shall follow the proof of Lemma 7 in [17] and use the facts about D-invariant ideals that were proved in § 3. Suppose that Lemma 5.3 is false and there exists a homogeneous prime ideal $\mathfrak{q} \subset \mathfrak{p}$, $\mathfrak{q} \neq (0)$, that satisfies the condition that $T\mathfrak{q} \subset \mathfrak{q}$.

According to Lemma 6 of [17], there exists a zero $\overline{\beta} \in \mathcal{K}^4$ of the ideal \mathfrak{p} for which

$$\operatorname{ord} \mathfrak{p}(\overline{f}) \leqslant rN(\mathfrak{p}) \cdot \|\overline{f} - \overline{\beta}\| + B(\mathfrak{p}).$$

Taking into account the inequality in the lemma, we conclude that $\|\overline{f} - \overline{\beta}\| > 1$.

Let n be chosen so that ord $\beta_n = |\overline{\beta}|$. Let V be any homogeneous polynomial in \mathfrak{p} , and let $W = x_n^d$, where $d = \deg_{\overline{x}} V$. If we apply Lemma 1 of [17] to the polynomials V and W, then we obtain the inequality

$$\operatorname{ord} \left(V(\overline{f}) W(\overline{\beta}) - V(\overline{\beta}) W(\overline{f}) \right) \geqslant \|\overline{f} - \overline{\beta}\| + \left(|\overline{f}| + |\overline{\beta}| \right) d$$

or, since $V(\overline{\beta})=0,$ $|\overline{f}|=0,$ and $||\overline{f}-\overline{\beta}||>1,$ we have the inequality

$$\operatorname{ord} V(\overline{f}) + d \cdot \operatorname{ord} \beta_n > 1 + |\overline{\beta}| \cdot d.$$

Recalling the definition of n, we finally obtain

$$V(\overline{f}) > 1. (40)$$

The inequality (40) does not hold for the polynomial $V = x_0$. Hence, $x_0 \notin \mathfrak{p}$ and so $x_0 \notin \mathfrak{q}$.

Let \mathfrak{q}_0 denote the ideal of $\mathfrak{R} = \mathbb{C}[z, x_1, x_2, x_3]$ consisting of all polynomials C for which there exist polynomials $C_0 \in \mathfrak{q}$ and $B \in \mathcal{A}$ such that

$$C = C_0 + (x_0 - 1)B, (41)$$

where C_0 is homogeneous. We claim that \mathfrak{q}_0 is a prime ideal. In fact, if $U, V \in \mathfrak{R}$ are such that $UV \in \mathfrak{q}_0$, then there exist homogeneous polynomials U_0 , V_0 and W_0 with

$$U = U_0 + (x_0 - 1)U_1, \qquad V = V_0 + (x_0 - 1)V_1, \qquad UV = W_0 + (x_0 - 1)W_1, \quad W_0 \in \mathfrak{q},$$

where $U_1, V_1, W_1 \in \mathcal{A}$. We may suppose that $\deg_{\overline{x}} U_0 + \deg_{\overline{x}} V_0 = \deg_{\overline{x}} W_0$, since otherwise we could multiply any of the polynomials U_0 , V_0 or W_0 by a suitable power of x_0 to obtain the equality. From the above relations it follows that

$$W_0 - U_0 V_0 = (x_0 - 1)G, \qquad G \in \mathcal{A},$$

and, since $W_0 - U_0 V_0$ is a homogeneous polynomial, we conclude that $W_0 - U_0 V_0 = 0$ and $U_0 V_0 = W_0 \in \mathfrak{q}$. Since \mathfrak{q} is a prime ideal, either U_0 or V_0 lies in \mathfrak{q} . But then, by the definition of \mathfrak{q}_0 , either U or V must belong to \mathfrak{q}_0 . This proves that \mathfrak{q}_0 is a prime ideal.

Applying the operator T to (41), we obtain

$$TC = TC_0 + (x_0 - 1)TB$$

and for some $B_1 \in \mathcal{A}$ we have

$$DC = TC_0 + (x_0 - 1)B_1$$

Since $T\mathfrak{q} \subset \mathfrak{q}$, it follows that $TC_0 \in \mathfrak{q}$, and so $DC \in \mathfrak{q}_0$. We have thus proved that \mathfrak{q}_0 is a *D*-invariant ideal. From (41) and (40) it follows that for any $C \in \mathfrak{q}_0$

$$\operatorname{ord} C(\overline{f}) = \operatorname{ord} C_0(\overline{f}) > 1. \tag{42}$$

In particular, this means that the point (0,1,1,1) is a zero of the ideal \mathfrak{q}_0 . By Proposition 4.2, either $z \in \mathfrak{q}_0$, or else $\Delta = x_2^3 - x_3^2 \in \mathfrak{q}_0$. However, neither of these two polynomials satisfies (42). This contradiction completes the proof of Lemma 5.3.

Lemma 5.4. Let $\mathfrak{p} \subset \mathcal{A}$ be a homogeneous prime ideal with $\mathfrak{p} \cap \mathbb{C}[z] = (0)$, $r = 4 - h(\mathfrak{p}) \geqslant 1$, and

$$\operatorname{ord}\mathfrak{p}(\overline{f}) > B(\mathfrak{p}) + 3N(\mathfrak{p}).$$

Then there exists a homogeneous polynomial $A \in \mathfrak{p}$ such that $C = TA \notin \mathfrak{p}$ and

$$\deg_{\overline{x}} C \leqslant 9\lambda^2 N(\mathfrak{p})^{1/(4-r)}, \quad \deg_z C \leqslant 3\lambda \big(B(\mathfrak{p})+1\big) N(\mathfrak{p})^{-(3-r)/(4-r)}, \tag{43}$$

where, as in Lemma 5.2, $\lambda = 18^3 = 5832$.

Proof. Let E be a homogeneous polynomial in \mathfrak{p} , for which the following expression attains its minimum value:

$$N(\mathfrak{p}) \deg_z E + (B(\mathfrak{p}) + 1) \deg_{\overline{x}} E.$$

The polynomial E is obviously irreducible. For brevity we set $L = \deg_{\overline{x}} E$ and $M = \deg_z E$. We define μ and ν using (39). By Lemma 5.2, we find that

$$N(\mathfrak{p})M + (B(\mathfrak{p}) + 1)L \leqslant N(\mathfrak{p})\mu + (B(\mathfrak{p}) + 1)\nu \leqslant 3\lambda(B(\mathfrak{p}) + 1)N(\mathfrak{p})^{1/(4-r)},$$

from which it follows that

$$L \leqslant 3\lambda N(\mathfrak{p})^{1/(4-r)}, \qquad M \leqslant 3\lambda (B(\mathfrak{p}) + 1)N(\mathfrak{p})^{-(3-r)/(4-r)}. \tag{44}$$

Since $\mathfrak{p} \cap \mathbb{C}[z] = (0)$, we have $L \geqslant 1$.

We set $E_0 = E$ and $E_1 = TE$. We have

$$\deg_z E_1 \leqslant M$$
, $\deg_{\overline{x}} E_1 \leqslant L + 1$.

If $E_1 \notin \mathfrak{p}$, then we set A = E and $C = E_1$, and Lemma 5.4 holds. Hence, from now on we suppose that $E_1 \in \mathfrak{p}$.

We set $J_0 = (E_0)$ and $J_1 = (E_0, E_1)$. We have $h(J_0) = 1$. The polynomial E_0 is irreducible, and $E_0 \nmid E_1$, since otherwise the principal prime ideal J_0 would be T-invariant, and this is impossible by Lemma 5.3. But then J_1 is an unmixed ideal with $h(J_1) = 2$. According to Lemma 4 of [17] applied to the ideals J_0 and J_1 and the polynomial E_1 we have

$$N(J_1) \leqslant N(J_0) \deg_{\overline{x}} E_1 \leqslant (L+1)^2,$$
 $B(J_1) \leqslant B(J_0) \deg_{\overline{x}} E_1 + N(J_0) \deg_{\overline{x}} E_1 \leqslant M(2L+1) \leqslant 2M(L+1).$

We have used the relations in Lemma 3.1 to bound $N(J_0)$ and $B(J_0)$.

We set

$$\chi_0 = 0$$
, $\chi_1 = a_1 = c_1 = 1$, $b_1 = 2$

and define χ_i, a_i, b_i, c_i as follows for i = 2, 3:

$$\chi_{i+1} = \chi_i + 2\lambda b_i, \qquad i = 1, 2, \tag{45}$$

$$c_{i+1} = c_i + \lambda b_i, \qquad i = 1, 2,$$
 (46)

$$a_{i+1} = a_i c_{i+1}, i = 1, 2,$$
 (47)

$$b_{i+1} = b_i c_{i+1} + a_i, \qquad i = 1, 2.$$
 (48)

For each k = 0, 1, 2, 3 we now let J_k denote the ideal of A that is generated by the polynomials T^iE , $0 \le i \le \chi_k$.

Let k be the largest integer such that $J_k \subset \mathfrak{p}$ and there exist homogeneous polynomials E_0, \ldots, E_k satisfying the conditions

- (1) $\deg_{\overline{x}} E_j \leqslant c_j(L+1)$ and $\deg_z E_j \leqslant M$ for $j=1,\ldots,k$;
- (2) the ideal $a_k = (E_0, \dots, E_k)$ is contained in J_k ;
- (3) all the primary components of a_k contained in \mathfrak{p} have height k+1, and if u_k is the unmixed ideal that is the intersection of these components, then

$$N(u_k) \le a_k (L+1)^{k+1}, \qquad B(u_k) \le b_k M (L+1)^k.$$
 (49)

For k=1 it is easy to see that these conditions hold with the polynomials E_0 and E_1 that we constructed before; hence, $k \ge 1$. On the other hand, the inclusion $\mathfrak{u}_k \subset \mathfrak{p}$ implies that $k+1=h(\mathfrak{u}_k) \leqslant h(\mathfrak{p})=4-r$, so that $k \leqslant 3-r$. We suppose that $J_{k+1} \subset \mathfrak{p}$ and show how this would lead to a contradiction.

Let \mathfrak{b} be a primary component of the ideal \mathfrak{u}_k , set $\mathfrak{q} = \sqrt{\mathfrak{b}}$, and let l be the exponent of \mathfrak{q} , that is, the smallest natural number such that $\mathfrak{q}^l \subset \mathfrak{b}$. We now prove that

$$l \leqslant 2\lambda b_k. \tag{50}$$

Suppose that, on the contrary, $l > 2\lambda b_k$. Then, according to (49) and the relation (1) in Proposition 2 of [17], we have

$$a_k(L+1)^{k+1} \geqslant N(\mathfrak{u}_k) \geqslant lN(\mathfrak{q}) \geqslant l > 2\lambda b_k \geqslant 2\lambda a_k.$$

This implies that $(L+1)^3 \ge (L+1)^{k+1} > 2\lambda$ and $L \ge 22$. Furthermore,

$$\left(\frac{L-1}{L+1}\right)^3 \geqslant \left(\frac{21}{23}\right)^3 > \frac{1}{2}$$

and $(L+1)^3 < 2(L-1)^3$. Using Proposition 2 of [17] and (49), from this we see that

$$\lambda B(\mathfrak{q}) \leqslant \frac{\lambda B(\mathfrak{u}_k)}{l} \leqslant \frac{B(\mathfrak{u}_k)}{2b_k} \leqslant \frac{1}{2}(M+1)(L+1)^k < (M+1)(L-1)^k,$$
$$\lambda N(\mathfrak{q}) \leqslant \frac{\lambda N(\mathfrak{u}_k)}{l} \leqslant \frac{N(\mathfrak{u}_k)}{2b_k} \leqslant \frac{1}{2}(L+1)^{k+1} < (L-1)^{k+1}.$$

We have also used the fact that $a_k \leq b_k$. Thus, the ideal q satisfies (38) with $\nu = L - 1$ and $\mu = M$. By Lemma 5.2, there exists a homogeneous polynomial $G \in \mathfrak{q}$ that satisfies the inequalities

$$\deg_{\overline{x}} G \leqslant L - 1, \qquad \deg_z G \leqslant M.$$

Since $G \in \mathfrak{q} \subset \mathfrak{p}$, these inequalities contradict the definition of the polynomial E. This proves (50).

We next prove that there exist i,j with $0 \le i \le k$ and $0 \le j \le 2\lambda b_k$ such that $T^jE_i \notin \mathfrak{q}$. The ideal \mathfrak{q} is isolated in the set of associated prime ideals of \mathfrak{a}_k . Hence, there exists a polynomial $H \notin \mathfrak{q}$ such that $G^lH \in \mathfrak{a}_k$ for any $G \in \mathfrak{q}$. Suppose that there did not exist i,j with the desired properties. Then, since $l \le 2\lambda b_k$, we should have $T^l(G^lH) \in \mathfrak{q}$. Since $G \in \mathfrak{q}$, this implies that $(TG)^lH \in \mathfrak{q}$; and, since \mathfrak{q} is a prime ideal and $H \notin \mathfrak{q}$, we have $TG \in \mathfrak{q}$. Thus, $T\mathfrak{q} \subset \mathfrak{q}$, which contradicts Lemma 5.3. This proves the existence of indices i,j with the desired properties.

Let $\mathfrak{q}_1,\ldots,\mathfrak{q}_s$ be all of the associated prime ideals of \mathfrak{u}_k . By what was proved above for every v with $1\leqslant v\leqslant s$ there exist i_v,j_v with $0\leqslant i_v\leqslant k$ and $0\leqslant j_v\leqslant 2\lambda b_k$ such that $T^{j_v}E_{i_v}\notin\mathfrak{q}_v$. We set

$$E_{k+1} = \sum_{v=1}^{s} \eta_v T^{j_v} E_{i_v},$$

where the $\eta_v \in \mathbb{C}$ are chosen so that $E_{k+1} \notin \mathfrak{q}_v$ for $1 \leq v \leq s$, and we prove that the polynomials E_0, \ldots, E_{k+1} satisfy the conditions (1)-(3) with k replaced by k+1.

Since the operator T does not increase the degree in z of a polynomial, it follows that $\deg_z E_{k+1} \leq M$. Using (46), we also obtain

$$\deg_{\overline{x}} E_{k+1} \leqslant c_k(L+1) + 2\lambda b_k \leqslant c_{k+1}(L+1).$$

Thus, (1) holds with k replaced by k + 1. Condition (2) follows from (45).

Let \mathfrak{r} be an associated prime ideal of \mathfrak{a}_{k+1} that is contained in \mathfrak{p} (such ideals exist because $\mathfrak{a}_{k+1} \subset J_{k+1} \subset \mathfrak{p}$). Since $\mathfrak{r} \supset \mathfrak{a}_{k+1} = (\mathfrak{a}_k, E_{k+1})$, we have $\mathfrak{p} \supset \mathfrak{r} \supset \mathfrak{a}_k$. If we take into account the fact that the set of all associated prime ideals of \mathfrak{a}_k contained in \mathfrak{p} coincides with the set $\mathfrak{q}_1, \ldots, \mathfrak{q}_s$, then we conclude that there exists an ideal $\mathfrak{q}_j \subset \mathfrak{r}$. Hence, $h(\mathfrak{r}) \geqslant h(\mathfrak{q}_j) = k+1$. If $h(\mathfrak{r}) = k+1$, then $\mathfrak{r} = \mathfrak{q}_j$. But this is impossible, since $E_{k+1} \in \mathfrak{r}$ and $E_{k+1} \notin \mathfrak{q}_j$. Thus, $h(\mathfrak{r}) \geqslant k+2$; since \mathfrak{a}_{k+1} is generated by k+2 polynomials, this means that $h(\mathfrak{r}) = k+2$. We have proved that all the primary components of \mathfrak{a}_{k+1} that are contained in \mathfrak{p} have height k+2.

Let $\mathfrak{a}_k = \mathfrak{u}_k \cap \mathfrak{a}'$, where \mathfrak{a}' is the intersection of the primary components of \mathfrak{a}_k that do not occur in \mathfrak{u}_k . If I were a primary component of \mathfrak{u}_{k+1} and $\mathfrak{r} = \sqrt{I}$, then the inclusion $\mathfrak{r} \supset \mathfrak{a}'$ would imply that $\mathfrak{p} \supset \mathfrak{r} \supset \mathfrak{a}'$, which is impossible. Hence, \mathfrak{a}' is not contained in \mathfrak{r} ; so from the inclusion $\mathfrak{u}_k \cap \mathfrak{a}' = \mathfrak{a}_k \subset \mathfrak{a}_{k+1} \subset I$ we obtain $\mathfrak{u}_k \subset I$. Thus, $\mathfrak{u}_k \subset \mathfrak{u}_{k+1}$, and since $E_{k+1} \in \mathfrak{a}_{k+1} \subset \mathfrak{u}_{k+1}$, we have $(\mathfrak{u}_k, E_{k+1}) \subset \mathfrak{u}_{k+1}$. Using Lemma 4 of [17], (49), (47) and (48), we see that this inclusion implies that

$$N(\mathfrak{u}_{k+1}) \leqslant N(\mathfrak{u}_k) \deg_{\overline{x}} E_{k+1} \leqslant a_{k+1} (L+1)^{k+2},$$

 $B(\mathfrak{u}_{k+1}) \leqslant B(\mathfrak{u}_k) \deg_{\overline{x}} E_{k+1} + N(\mathfrak{u}_k) \deg_z E_{k+1} \leqslant b_{k+1} M(L+1)^{k+1}.$

We have proved that the polynomials E_0,\ldots,E_{k+1} satisfy all the conditions (1)–(3) with k replaced by k+1. In other words, the assumption that $J_{k+1} \subset \mathfrak{p}$ led us to a contradiction with the definition of k. Thus, J_{k+1} is not contained in \mathfrak{p} and, since $k \leq 3-r \leq 2$, the inclusion $J_3 \subset \mathfrak{p}$ is impossible. This means that there exists an index i $(0 \leq i \leq \chi_3 - 1)$ such that $A = T^iE \in \mathfrak{p}$ but $C = TA \notin \mathfrak{p}$. From (45)–(48) it follows that $\chi_3 = 8\lambda^2 + 10\lambda + 1$. Finally, using (44) we find that

$$\begin{split} \deg_z C \leqslant M \leqslant 3\lambda \big(B(\mathfrak{p}) + 1\big) N(\mathfrak{p})^{-(3-r)/(4-r)}, \\ \deg_{\overline{\pi}} C \leqslant \deg_{\overline{\pi}} E + \chi_3 &= L + \chi_3 \leqslant 9\lambda^2 N(\mathfrak{p})^{1/(4-r)}. \end{split}$$

Lemma 5.4 is proved.

§ 6. End of the proof of Theorem 5

We now proceed directly to the proof of Theorem 5. Suppose that there exist ideals that satisfy the hypothesis of the theorem but not its conclusion. Let I be such an ideal having maximal height h(I) and let r = 4 - h(I). Then

ord
$$I(\overline{f}) > \rho^{2r-1} (B(I)N(I)^{r/(r-4)} + N(I)^{3/(4-r)}).$$
 (51)

Let I_1, \ldots, I_s be all the primary components of I that have trivial intersection with $\mathbb{C}[z]$, let k_j be their exponents and let $\mathfrak{p}_j = \sqrt{I_j}$ be their radicals. By assumption,

 $s \geqslant 1$. The following equalities hold for some $b \in \mathbb{C}[z]$ by Proposition 2 of [17]:

$$\sum_{j=1}^{s} k_j N(\mathfrak{p}_j) = N(I), \tag{52}$$

$$\deg_z b + \sum_{j=1}^s k_j B(\mathfrak{p}_j) = B(I), \tag{53}$$

$$\operatorname{ord} b + \sum_{j=1}^{s} k_{j} \operatorname{ord} \mathfrak{p}_{j}(\overline{f}) = \operatorname{ord} I(\overline{f}).$$
 (54)

If for all j one has $\operatorname{ord} \mathfrak{p}_{j}(\overline{f}) \leqslant \rho^{2r-1} (B(\mathfrak{p}_{j})N(\mathfrak{p}_{j})^{r/(4-r)} + N(\mathfrak{p}_{j})^{3/(4-r)})$, then from (52)–(54) we obtain

$$\operatorname{ord} I(\overline{f}) = \operatorname{ord} b + \sum_{j=1}^{s} k_{j} \operatorname{ord} \mathfrak{p}_{j}(\overline{f})$$

$$\leq \operatorname{deg}_{z} b + \sum_{j=1}^{s} k_{j} \rho^{2r-1} B(\mathfrak{p}_{j}) N(\mathfrak{p}_{j})^{r/(4-r)} + \rho^{2r-1} \sum_{j=1}^{s} k_{j} N(\mathfrak{p}_{j})^{3/(4-r)}$$

$$\leq \rho^{2r-1} N(I)^{r/(4-r)} \left(\operatorname{deg}_{z} b + \sum_{j=1}^{s} k_{j} B(\mathfrak{p}_{j}) \right) + \rho^{2r-1} \left(\sum_{j=1}^{s} k_{j} N(\mathfrak{p}_{j}) \right)^{3/(4-r)}$$

$$\leq \rho^{2r-1} (B(I) N(I)^{r/(4-r)} + N(I)^{3/(4-r)}).$$

which contradicts (51). This contradiction means that there exists a prime ideal \mathfrak{p} with $\mathfrak{p} \cap \mathbb{C}[z] = (0)$ and $h(\mathfrak{p}) = 4 - r$ such that

$$\operatorname{ord}\mathfrak{p}(\overline{f}) > \rho^{2r-1}(B(\mathfrak{p})N(\mathfrak{p})^{r/(4-r)} + N(\mathfrak{p})^{3/(4-r)}). \tag{55}$$

Let A and C be the polynomials whose existence was proved in Lemma 5.4.

We first consider the case r=1. If we apply Lemma 5.1 to the ideal \mathfrak{p} , the polynomials A and C, and the vector \overline{f} (in our case s=0), and if we use (43) with r=1, then we find that

$$\operatorname{ord}\mathfrak{p}(\overline{\varphi})\leqslant B(\mathfrak{p})\operatorname{deg}_{\overline{x}}C+N(\mathfrak{p})\operatorname{deg}_{z}C\leqslant (9\lambda^{2}+3\lambda)B(\mathfrak{p})N(\mathfrak{p})^{1/3}+3\lambda N(\mathfrak{p}),$$

which contradicts (55) for r=1. Thus, this case is impossible.

In the case $r \ge 2$ we consider the ideal J whose existence is ensured by Lemma 5.1. We apply the lemma to the ideal $\mathfrak p$ and the polynomials A and C that are constructed in Lemma 5.4. Using (43), according to Lemma 5.1 we have the bounds

$$N(J) \leqslant N(\mathfrak{p}) \deg_{\overline{x}} C \leqslant 9\lambda^2 N(\mathfrak{p})^{(5-r)/(4-r)}, \tag{56}$$

$$B(J) \leqslant B(\mathfrak{p}) \deg_{\overline{x}} C + N(\mathfrak{p}) \deg_{z} C$$

$$\leqslant (9\lambda^{2} + 3\lambda)B(\mathfrak{p})N(\mathfrak{p})^{1/(4-r)} + 3\lambda N(\mathfrak{p})^{1/(4-r)}.$$
 (57)

The ideal J satisfies the relation: $h(J) = h(\mathfrak{p}) + 1 > h(I)$. From the definition of I it follows that

ord
$$J(\overline{f}) \leq \rho^{2r-3} (B(J)N(J)^{(r-1)/(5-r)} + N(I)^{3/(5-r)}).$$

Using this inequality, the inequality (3) of Lemma 5.1, and also (56), (57) and (43), we find that

$$\begin{split} \operatorname{ord} \mathfrak{p}(\overline{f}) \leqslant \operatorname{ord} J(\overline{f}) + B(\mathfrak{p}) \operatorname{deg}_{\overline{x}} C + N(\mathfrak{p}) \operatorname{deg}_{z} C \\ \leqslant \rho^{2r-3} \Big((9\lambda^{2} + 3\lambda) B(\mathfrak{p}) N(\mathfrak{p})^{\frac{1}{4-r}} + 3\lambda N(\mathfrak{p})^{\frac{1}{4-r}} \Big) \Big(9\lambda^{2} N(\mathfrak{p})^{\frac{5-r}{4-r}} \Big)^{\frac{r-1}{5-r}} \\ + \rho^{2r-3} \Big(9\lambda^{2} N(\mathfrak{p})^{\frac{5-r}{4-r}} \Big)^{\frac{3}{5-r}} + (9\lambda^{2} + 3\lambda) B(\mathfrak{p}) N(\mathfrak{p})^{\frac{1}{4-r}} + 3\lambda N(\mathfrak{p})^{\frac{1}{4-r}} \\ < \rho^{2r-1} \Big(B(\mathfrak{p}) N(\mathfrak{p})^{r/(4-r)} + N(\mathfrak{p})^{3/(4-r)} \Big). \end{split}$$

This inequality contradicts (55) and thereby completes the proof of Theorem 5.

This proof of Theorem 5 was presented in such a way as to open up the possibility of obtaining a more general result. We consider an arbitrary system of differential equations

$$t(z)y'_{j} = F_{j}(z, y_{1}, \dots, y_{m}), \qquad j = 1, \dots, m,$$
 (58)

where t(z) and $F_i(z, y_1, \ldots, y_m)$ are polynomials in all the indicated variables.

Definition. We say that a solution $f_1(z), \ldots, f_m(z)$ of the system (58) has the D-property at the point $\xi \in \mathbb{C}$ if these functions are analytic at ξ and the set of all prime ideals of $\mathbb{C}[z, y_1, \ldots, y_m]$ having the following two properties has non-zero intersection: (1) the ideal is invariant relative to the operator

$$D = t(z)\frac{\partial}{\partial z} + \sum_{i=1}^{m} F_j(z, y_1, \dots, y_m) \frac{\partial}{\partial y_j},$$

and (2) its variety of zeros in \mathbb{C}^{m+1} does not contain the analytic curve $(z, f_1(z), \ldots, f_m(z))$, but contains the point of this curve corresponding to $z = \xi$.

If the functions $f_i(z)$ are algebraically independent over $\mathbb{C}(z)$, then this condition means that either the set of *D*-invariant prime ideals having zero at $(\xi, f_1(\xi), \ldots, f_m(\xi))$ has non-trivial intersection, or else there are no such ideals at all.

For example, in Proposition 2 of §4 we saw that the set of functions (P(z), Q(z), R(z)) has the D-property at the point $\xi = 0$. The polynomial $z(y_2^3 - y_3^2)$ is contained in the intersection of all D-invariant prime ideals with zero at (0, 1, 1, 1). In [17] we used the fact that if a set of functions $f_1(z), \ldots, f_m(z)$ is algebraically independent over $\mathbb{C}(z)$ and satisfies a homogeneous (in y_j) linear system of differential equations (58), then it has the D-property at any point $\xi \in \mathbb{C}$.

It is easy to prove that if ξ is not a zero of the polynomial t(z) (that is, ξ is not a singular point), then any solution to (58) that is analytic at ξ has the D-property. In fact, suppose that $(\xi, f_1(\xi), \ldots, f_m(\xi))$ is one of the zeros of the D-invariant prime ideal \mathfrak{p} . Since we have the identity $t(z)^n \varphi^{(n)}(z) = B(z, f_1(z), \ldots, f_m(z))$ for any polynomial $A \in \mathfrak{p}$ and any $n \geq 0$, where $\varphi(z) = A(z, f_1(z), \ldots, f_m(z))$ and $B = t(z)^n (t(z)^{-1}D)^n A \in \mathfrak{p}$, it follows that, in view of the condition $t(\xi) \neq 0$, we have $\varphi^{(n)}(\xi) = 0$ for $n \geq 0$; that is, $\varphi(z) \equiv 0$, and so the entire curve $(z, f_1(z), \ldots, f_m(z))$ belongs to the variety of zeros of \mathfrak{p} . Thus, in this case there do not exist any prime ideals with the properties in the definition. A similar situation occurred

in [16], where we considered a set of functions that satisfies a system of differential equations (58) with constant coefficients $(t(z) \equiv 1 \text{ and the } F_j \text{ do not depend on } z)$, and hence has the D-property at any point ξ .

The next result is a generalization of Theorem 3.

Theorem 6. Suppose that the solution $f_1(z), \ldots, f_m(z)$ of the system (58) consists of functions that are algebraically independent over $\mathbb{C}(z)$ and have the D-property at a point $\xi \in \mathbb{C}$. Then for any $A \in \mathbb{C}[z, x_1, \ldots, x_m]$ with $A \not\equiv 0$ one has

$$\operatorname{ord}_{z=\xi} A(z, f_1(z), \dots, f_m(z)) \leq \gamma_5 (\deg_z A + 1) (\deg_{\overline{x}} A)^m,$$

where γ_5 is a constant that depends only on the point ξ and the functions $f_i(z)$.

Proof. The proof of this theorem is almost identical to the proof of Theorem 3. Instead of the inequality in Theorem 5, one uses induction on r, $1 \le r \le m$, to prove that the inequality

ord
$$I(\overline{f}) \le \rho^{2r-1} (B(I)N(I)^{r/(m+1-r)} + N(I)^{m/(m+1-r)})$$
 (59)

(with ρ a sufficiently large constant) holds for any unmixed homogeneous ideal $I \subset \mathcal{A} = \mathbb{C}[z,x_0,\ldots,x_m]$ for which $I \cap \mathbb{C}[z] = (0)$ and h(I) = m+1-r. To do this it suffices to make some natural modifications in the proof of Lemma 5.3. One makes a new choice of sequences a_i,b_i,c_i,χ_i ; and the bound on the degrees in z of the E_j in the proof of Lemma 5.4 is written in the form $\deg_z E_j \leqslant c_j(M+1)$. No significant changes in the proof are needed.

The algebraic independence condition for the functions $f_j(z)$ in Theorem 6 can be omitted, and one can also prove a bound for the sum of the multiplicities of the zeros at different points. The next result is a generalization of Theorem 6 and the main theorems of [16] and [17].

Theorem 7. Suppose that the solution $f_1(z), \ldots, f_m(z)$ of the system (58) has the D-property at each of the distinct complex points ξ_1, \ldots, ξ_q . Then for any $A \in \mathbb{C}[z, x_1, \ldots, x_m]$ with $A(z, f_1(z), \ldots, f_m(z)) \not\equiv 0$ one has

$$\sum_{j=1}^{q} \operatorname{ord}_{z=\xi_{j}} A(z, f_{1}(z), \dots, f_{m}(z)) \leqslant \gamma_{6} (\deg_{z} A + q) (\deg_{\overline{x}} A)^{k},$$

where γ_6 is a constant that depends only on the points ξ_j , the functions $f_i(z)$ and the system (58), and k is the maximum number of functions $f_i(z)$ that are algebraically independent over $\mathbb{C}(z)$.

This theorem is proved by the method used above. In Lemma 5.3 the condition that $\mathfrak{q} \neq (0)$ must be replaced by the condition that \mathfrak{q} not be contained in the ideal \mathfrak{E} consisting of all the algebraic relations among the functions $f_j(z)$ over $\mathbb{C}[z]$. The analogue of the inequality in Theorem 5 has the form

$$\sum_{j=1}^{q} \operatorname{ord}_{z=\xi_{j}} I(\overline{f}) \leqslant \rho^{2r-1} (B(I)N(I)^{r/(k+1-r)} + qN(I)^{k/(k+1-r)}),$$

where, as before, r = m + 1 - h(I); this is proved by induction on r for $1 \le r \le k$. In [16] and [17] one can find all of the modifications in the proof that are needed because one has a large number of points.

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Moscow State University

E-mail address: nest@nw.math.msu.su

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