p-adic logarithmic functions and applications

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Abstract We explain how we define *p*-adic logarithmic functions to provide a new lower bound for linear forms in two *p*-adic elliptic logarithms proven in [11]. We adapt the argument that relies on the interpolation method on the variable change introduced by G. Chudnovsky, and on Faà-di-Bruno's formula adapted to matrices whose elements are *p*-adic elliptic logarithmic functions.

1 Introduction

Let K be a number field of finite degree D over \mathbb{Q} . Denote the ring of integers by \mathfrak{O} . Let $A, B \in K$, $\Delta := 4A^3 - 27B^2 \neq 0$ and \mathcal{E} be an elliptic curve defined by

$$Y^2 = X^3 - AX - B.$$

We may assume $A, B \in \mathfrak{O}$ (for; if A or $B \notin \mathfrak{O}$, then there exists a suitable $c \in \mathfrak{O}$ such that the elliptic curve $Y^2 = X^3 - A'X - B'$ with $A' = c^4A \in \mathfrak{O}$, $B' = c^6B \in \mathfrak{O}$ and with the discriminant $\Delta' = c^{12}\Delta$, is isomorphic to \mathcal{E} since the j-invariant remains equal under these multiplications).

Let us denote by $\overline{\mathbb{Q}}$ the algebraic closure of \mathbb{Q} in \mathbb{C} . Let p be a rational prime $\in \mathbb{Q}$ and $|\cdot|_{\infty}$ be an Archimedean valuation on K. For a place v of K over p, we write the valuation $|\cdot|_v$ normalized such that $|x|_v = p^{-ord_p(x)}$ for $x \in \mathbb{Q}$. Denote K_v the completion of K by v, and write \mathbb{Q}_p the completion of \mathbb{Q} by p. The field K_v is a finite extension of \mathbb{Q}_p of local degree $n_v = [K_v : \mathbb{Q}_p]$ with $\sum_{v|p} n_v = D$. Put \mathbb{C}_p the completion of the

algebraic closure of K_v . We note that the algebraic closure of K_v is not complete itself. It is well-known that \mathbb{C}_p is algebraically closed complete field of characteristic 0, in which the algebraic closure of K_v is dense and that there are D distinct embeddings of K into \mathbb{C}_p . Denote again by $|x|_v$ the extension of $|x|_v$ on \mathbb{C}_p .

For $\underline{x} \in \mathbb{P}_N(\overline{\mathbb{Q}})$ having coordinates $\underline{x} = (x_0, \dots, x_N) \in \mathbb{P}_N(K)$, define the absolute logarithmic height of \underline{x} by

$$h(\underline{x}) = rac{1}{[K:\mathbb{Q}]} \sum_v n_v \log(\max\{|x_0|_v,\ldots,|x_N|_v\})$$

where the sum runs over all the normalized places of K. This definition is independent of the choice of the projective coordinates and the choice of the field containing x_0, \ldots, x_N .

Let $a \in \overline{\mathbb{Q}}$ and put h(a) := h(1:a), the absolute logarithmic height of the algebraic number a. We may write $h(a) = h_{\infty}(a) + h_{f}(a)$ where the sum in $h_{\infty}(a)$ runs over all the infinite places and the sum in $h_{f}(a)$ runs over all the finite places:

$$h_{\infty}(a) = rac{1}{[K:\mathbb{Q}]} \sum_{v ext{ infinite}} n_v \log(\max\{1,|a|_v\}),$$

$$h_f(a) = \frac{1}{[K:\mathbb{Q}]} \sum_{v \text{ finite}} n_v \log(\max\{1,|a|_v\}).$$

Now we fix a place v over p and denote $|\cdot| = |\cdot|_v$. For a formal power series $f(z) = \sum_{k=0}^{\infty} a_k z^k \in \mathbb{C}_p[[z]]$, f(z) converges at $z \in \mathbb{C}_p$ if and only if $|a_k z^k| \to 0$. It is known that the radius of convergence is also given by Hadamard's formula.

Let us recall the Lutz-Weil p-adic elliptic function which corresponds to the p-adic version of the Weierstraß elliptic function \wp . Consider \mathcal{E} be an elliptic curve $\subset \mathbb{P}^2(\mathbb{C}_p)$:

$$ZY^2 = X^3 - AXZ^2 - BZ^3 \ (A, B \in \mathfrak{O}, \ 4A^3 \neq 27B^2).$$

Write
$$\lambda_p = \frac{1}{p-1}$$
 if $p \neq 2$, $\lambda_2 = 3$, $\mathcal{C}_p := \{z \in \mathbb{C}_p : |z| < p^{-\lambda_p}\}$ and $\mathcal{C}_v := \mathcal{C}_p \cap K_v$.

It is known that there exist two solutions φ and $-\varphi$ to the differential equation $(\varphi')^2 = 1 - A\varphi^4 - B\varphi^6$ with $\varphi(0) = 0$, defined over $\mathcal{C}_v \to K_v$, analytic in \mathcal{C}_v , after [18] [26]. In fact putting $\varphi^2 = \frac{1}{\wp_{\varphi}}$ we have $\left(\frac{\wp_{\varphi}'}{2}\right)^2 = \wp_{\varphi}^3 - A\wp_{\varphi} - B$ and $\varphi'(0) = 1$. The function $\varphi(z)$ is called the Lutz-Weil p-adic elliptic function. The elliptic curve can be given the structure of the p-adic Lie-group $\mathcal{E}(\mathbb{C}_p) \subset \mathbb{P}^2(\mathbb{C}_p)$ as follows. We may enlarge the domain of the definition of the function φ to \mathcal{C}_p (see e. g. the page 151 of [1] and [2], [23]).

Definition 1.1 For the p-adic Lie-group $\mathcal{E}(\mathbb{C}_p) \subset \mathbb{P}^2(\mathbb{C}_p)$ we have the exponential map:

$$\exp = \exp_{\mathcal{E}} : \mathcal{C}_p \to \mathcal{E}(\mathbb{C}_p) \subset \mathbb{P}^2(\mathbb{C}_p)$$

$$z \mapsto (\varphi(z), -\varphi'(z), \varphi^3(z))$$

The *p*-adic exponential map is locally analytic only. The function φ is odd and injective; indeed, $|\varphi(z)| = |z|, |\varphi'(z)| = 1$ for any $z \in \mathcal{C}_p$, hence $\exp_{\mathcal{E}}$ has no period [3]. There are corresponding addition formula and derivation formula, similar to those of φ .

Let $\beta \in K$. Take u_1 and u_2 in \mathcal{C}_p . We assume $\varphi^2(u_i)$ and $\frac{\varphi}{\varphi'}(u_i) \in K$ (i = 1, 2) i. e. $\exp(u_i) \in \mathcal{E}(K)$ (i = 1, 2). Put $\Lambda = \beta u_1 - u_2$ which is a linear form in two p-adic elliptic logarithms u_1 and u_2 . Write $\hat{h}(P) := \frac{1}{2} \lim_{n \to \infty} \frac{1}{4^n} h(2^n P)$ the Néron-Tate height defined on \mathcal{E} for a rational point $P \in \mathcal{E}(K)$.

We may suppose that none of these 3 numbers β, u_1, u_2 equals to 0, for, otherwise our statement trivially follows thanks to the Liouville inequality: $|\alpha| \geq e^{-[K:\mathbb{Q}]h(\alpha)}$ where $\alpha \in K, \alpha \neq 0$.

Denote non-negative real numbers $h_1, h_2, h_3, \rho, E, a_1, a_2, b$ and d by $h_i = \hat{h}(\exp(u_i))$ $(i = 1, 2), h_3 = \max(1, h(\beta)), \rho = p^{-\lambda_p}, E = \rho/\max(|u_1|, |u_2|), a_1 = \max(1, h_1), a_2 = \max(1, h_2), d = \max(1, \frac{[K:\mathbb{Q}]}{\log E}), g = \max(1, h_4, \log(h_1), \log(h_2), \log(d)).$ We denote further by $h = h_4 = h(\mathcal{E}) := \max\{1, h(1, A, B)\}$ the height of the elliptic curve \mathcal{E} .

Our proncipal result is as follows.

Theorem 1.1 [with R. Takada] Under the assumptions above, if we have

$$|\Lambda| \le \exp(-1.16 \times 10^{35} \times a_1 \cdot a_2 \cdot h_3 \cdot g^3 \cdot d^6 \cdot \log E),$$

then we obtain

$$\Lambda = 0$$

and $\beta = \frac{u_2}{u_1}$ is an algebraic number of degree at most 2 over $\mathbb Q$ with

$$h(\beta) \le \log \left(5.89 \times 10^{17} \times g^2 d^3 \times \max(a_1, \sqrt{a_1 a_2})\right).$$

Corollary 1.1 Whenever we have $\Lambda \neq 0$, then we obtain

$$|\Lambda| > \exp(-1.16 \times 10^{35} \times a_1 \cdot a_2 \cdot h_3 \cdot g^3 \cdot d^6 \cdot \log E).$$

We compare our result with that of G. Rémond and F. Urfels.

Put $b = \max(h_3, h_4, h_1, h_2, d)$ and $c = \max(1, h_4, \log b)$. The result in [20] shows, if

$$|\Lambda| \le \exp(-5.7 \times 10^{26} \times a_1 \cdot a_2 \cdot b \cdot c^3 \cdot d^6 \cdot \log E),$$

then

$$\Lambda = 0$$

and $\beta = \frac{u_2}{u_1}$ is an algebraic number of degree at most 2 over $\mathbb Q$ of height

$$\log(4.29 \times 10^{14} \times c^2 d^3 \times \max(a_1, \sqrt{a_1 a_2})).$$

We refine this result so as to obtain the best possible approximation concening with the height of algebraic coefficients of the linear forms since our bound does not contain $\log h_3$. Our constant is expressed in an explicit manner and the part of h_3 is separately written from other data. However, our numerical constant is larger than that of the statement of [20].

2 p-adic elliptic logarithmic function

We define the p-adic logarithmic function in elliptic case as a reversed function of the $\exp_{\mathcal{E}}$ with an expression of Formal group over \mathfrak{O} , following [15] [24] (see also [6][7]).

Let $P = (X,Y,1) \in \mathcal{E}(K)$. Put t = t(P) = -X/Y, $\omega(t) = -1/Y$. We have $P = (X,Y,1) = (t,-1,\omega(t)) = \left(\frac{\varphi(z)}{\varphi'(z)},-1,\frac{\varphi^3(z)}{\varphi'(z)}\right)$. Let r be a positive real number. We put $\mathcal{E}(r)$ the set of points P in $\mathcal{E}(K)$ with $|t(P)| \leq p^{-r}$. We include the origin in $\mathcal{E}(r)$ by convention, and then $\mathcal{E}(r)$ is a subgroup of $\mathcal{E}(K)$. Denote by \mathfrak{p}_r the set of elements $t \in K$ with $|t| \leq p^{-r}$. The map $P \to t(P)$ establishes a bijection between $\mathcal{E}(r)$ and \mathfrak{p}_r (Theorem 3.2, Chapter III, [15]). There is a power series expansion of $\omega(t)$ in t where the coefficients are polynomials in A, B with coefficients in \mathbb{Z} (Theorem 3.1, Chapter III, [15]). This power series expansion is studied in [7]. Below we rewrite estimates obtained in [7].

Lemma 2.1 Under the notations above, we have $\omega(t) = \sum_{n\geq 3} A_n t^n$ where $A_n \in \mathbb{Z}[A, B]$

is homogeneous of degree n-3 (of weights 4,6 on A,B) of form $A_3=1$ and

$$A_n = \sum_{\substack{4\lambda + 6\mu = n - 3, \\ \lambda, \mu > 0}} a_{\lambda, \mu}^{(n)} A^{\lambda} B^{\mu} ,$$

where $a_{\lambda,\mu}^{(n)} \in \mathbb{Z}$ with

$$|a_{\lambda,\mu}^{(n)}|_{\infty} \le \frac{3^3 \cdot 8^{n-3}}{n^3(\lambda+1)^3(\mu+1)^3} \quad (n \ge 3, \lambda \ge 0, \mu \ge 0) .$$

Moreover, we have

$$h(A_n) \le 3n + (n-3)h .$$

This lemma yields the estimate of the height of Taylor coefficients for the functions

$$\varphi^2(z) = \frac{\omega(t)}{t} = \sum_{n \geq 3} A_n t^{n-1},$$

$$\frac{\omega(t)}{t^2} = \sum_{n>3} A_n t^{n-2}.$$

Since $\exp_{\mathcal{E}}(z) = \left(\frac{1}{\varphi^2(z)}, \frac{-\varphi'(z)}{\varphi^3(z)}, 1\right) = (t, -1, \omega(t))$, the function z = z(t) corresponds to the logarithmic function which is introduced in [15] (see [7] [24]). By writing X, Y in terms of t and $\omega(t)$, the differential form $\Omega(t) = \frac{dX}{2Y}$ is viewed as a formal power series in t, and we define as in [15][24] the formal integral $\log_{\mathcal{E}}(t) = \int \Omega(t)$. With this formal integral we have;

$$\int \Omega(t) = \int rac{dX}{2Y} = \int rac{d\left(rac{1}{arphi^2}
ight)z(t)}{\left(rac{-2arphi'}{arphi^3}
ight)z(t)} dt = \int rac{\left(rac{-2arphi'}{arphi^3}
ight)z(t)}{\left(rac{-2arphi'}{arphi^3}
ight)z(t)} z'(t) dt = z(t)$$

which is indeed the local reversed function around the origin, of the function $t=t(P)=-\frac{X}{Y}=\frac{\varphi(z)}{\varphi'(z)}$.

Definition 2.1 Put $\log_{\mathcal{E}}(t) = z(t)$. We call the function an elliptic p-adic logarithmic function associate to \mathcal{E} .

We rewrite the statement in [7] for convenience in explicit calculations below, by using $h = h(\mathcal{E})$:

Lemma 2.2 The Taylor expansion of $\log_{\mathcal{E}}(t)$ is given by

$$\log_{\mathcal{E}}(t) = z(t) = \sum_{n \ge 1} B_n t^n$$

where
$$B_1 = 1, B_n = \frac{C_n}{2n}, C_n = \sum_{\substack{4\lambda + 6\mu = n - 1, \\ \lambda, \mu > 0}} b_{\lambda, \mu}^{(n)} A^{\lambda} B^{\mu} \quad (n \ge 1) \text{ with } b_{\lambda, \mu}^{(n)} \in \mathbb{Z} \text{ and }$$

$$|b_{\lambda,\mu}^{(n)}|_{\infty} \le \frac{(2^5 \cdot 3 \cdot 5^2)^n}{(n+2)^3(\lambda+1)^3(\mu+1)^3} \quad (n \ge 1, \lambda \ge 0, \mu \ge 0).$$

Concerning the height, we have

$$h(C_n) \le 9n + (n-1)h .$$

Moreover, the domain of convergence of $\log_{\mathcal{E}}(t)$ is $\{z \in \mathbb{C}_p : |z| < 1\}$.

3 Differential operator

Consider a point $u=(0,u_1,u_2)\in\mathbb{C}_p\times\mathcal{C}_p^2$ and the hyperplane W defined by $z_0=\beta z_1-z_2$. To prove our theorem, with respect to the fixed non-Archimedean valuation $|\cdot|=|\cdot|_v$, we note that there is no restriction to suppose $|\beta|\leq 1$, otherwise we may consider $\frac{1}{\beta}u_2-u_1$ instead of Λ .

We are going to look at (Λ, u_1, u_2) . We choose as in [10] a basis of W: $(\beta, 1, 0)$ and (-1, 0, 1). Put $\sigma = (\sigma_1, \sigma_2) \in \mathbb{Z}^2$, $\sigma_1, \sigma_2 \geq 0$, and a differential operator over \mathbb{C}_p^3 along W;

$$D_z^{\sigma} = \left(\beta \frac{\partial}{\partial z_0} + \frac{\partial}{\partial z_1}\right)^{\sigma_1} \circ \left(-\frac{\partial}{\partial z_0} + \frac{\partial}{\partial z_2}\right)^{\sigma_2}.$$

Introduce also a "divided differential operator" along W as in [8];

$$\Delta_z^{\sigma} := \frac{D_z^{\sigma}}{\sigma!} = \frac{D_z^{\sigma}}{\sigma_1! \sigma_2!} = \frac{1}{\sigma_1! \sigma_2!} \left(\beta \frac{\partial}{\partial z_0} + \frac{\partial}{\partial z_1} \right)^{\sigma_1} \circ \left(-\frac{\partial}{\partial z_0} + \frac{\partial}{\partial z_2} \right)^{\sigma_2}.$$

Put $\tau = (\tau_0, \tau_1, \tau_2) \in \mathbb{Z}^3, \tau_0, \tau_1, \tau_2 \geq 0$ and define with $\psi = \varphi^2$;

$$f_{\tau}: \mathbb{C}_p \times \mathcal{C}_p^2 \to \mathbb{C}_p$$

$$(z_0, z_1, z_2) \to z_0^{\tau_0} \psi(z_1)^{\tau_1} \psi(z_2)^{\tau_2}.$$

For T_0, T_1, T_2, S_0, S_1 which are parameters ≥ 0 in \mathbb{Z} with $S_0 \geq 5$, define a matrix

$$\mathcal{M} = (\Delta_z^{\sigma} f_{\tau}(su))_{\tau;(\sigma,s)} = (m_{\tau,\sigma,s}) \tag{1}$$

where the lines are indexed by $\mathcal{T} = \{\tau \in \mathbb{Z}^3 | 0 \le \tau_i \le T_i\}$, the columns by $\mathcal{S} = \{(\sigma, s) = (\sigma_1, \sigma_2, s) \in \mathbb{Z}^3 | \sigma_1 \ge 0, \sigma_2 \ge 0, |\sigma| := \sigma_1 + \sigma_2 < S_0, 0 \le s \le S_1\}$. The number of lines is $L := (T_0 + 1)(T_1 + 1)(T_2 + 1)$. The elements of the matrix are "divided derivatives" instead of the ordinary derivatives in [20].

4 Interpolation matrix

Lemma 4.1 Let $D = [K : \mathbb{Q}]$. For any $L \times L$ minor determinant Δ of \mathcal{M} , we suppose;

$$|\Delta| \le \exp\left(-D(\log(L!) - DL\Big(T_0(h_3 + 6) + 3.4S_0\log(T_0 + 1) + (S_0 + 1)(18 + h_4)\right) + 8S_1^2(T_1h_1 + T_2h_2) + (T_1 + T_2)(16h_4 + 60\log 2 + 12)\right).$$

Then the rank of M is strictly less than L.

We remove $S_0 \log S_0$ in Proposition 6.1 of [20], that is essential for our improvement. For this, we carry out the variable change from z to t.

Now we assume that the rank of \mathcal{M} equals to L. We shall show that there exists an $L \times L$ minor determinant $\Delta \neq 0$ of \mathcal{M} and give a lower bound for $|\Delta|$ which contradicts the assumption of Lemma 4.1.

Recall that our matrix (1) is defined by $\mathcal{M} = (\Delta^{\sigma} f_{\tau}(su))_{\tau;(\sigma,s)}$ where

$$f_{\tau} = z_0^{\tau_0} \psi(z_1)^{\tau_1} \psi(z_2)^{\tau_2}, \quad \psi(z) = \varphi(z)^2.$$

Definition 4.1 We order the set of the indices $(\sigma, s) = (\sigma_1, \sigma_2, s)$ of columns of \mathcal{M} as follows. We order the set of $\sigma := (\sigma_1, \sigma_2)$ by the quantity $|\sigma| = \sigma_1 + \sigma_2$, namely if $|\sigma| < |\sigma'|$ then define $\sigma < \sigma'$. If $|\sigma| = |\sigma'|$ then we order lexicographically (σ_1, σ_2) . We define an order for $(\sigma, s) = (\sigma_1, \sigma_2, s)$ firstly by the order defined above for σ and secondly by the order for s. Since the rank of \mathcal{M} equals to L, then there exist L-tuple of the indices of columns such that the corresponding $L \times L$ minor determinant is nonzero. Choose the minimal L-tuple among such ones by the order defined now. We denote the minimal L-tuple by $(\sigma_{\mu}, s_{\mu})_{1 \leq \mu \leq L}$. We put the corresponding square minimal matrix $\mathcal{N} = (m_{\tau, \sigma_{\mu}, s_{\mu}})$ and denote $\Delta = \det \mathcal{N}$.

We present some properties as follows [11].

Lemma 4.2 For a fixed μ_0 , every column of index $<(\sigma_{\mu_0}, s_{\mu_0})$ is contained in the subspace generated by the columns of index (σ_{μ}, s_{μ}) with $1 \le \mu < \mu_0$.

Lemma 4.3 We have $\det A = \det \mathcal{N} = \Delta \neq 0$.

We are now going to give an upper bound for the height of the number $a_{\tau,\mu}$ by doing variable change of the functions "from z to t".

Lemma 4.4 If s = 0, then we have

$$a_{\tau,\mu} = m_{\tau,\sigma_{\mu},0} = \Delta_z^{\sigma_{\mu}} f_{\tau}(0)$$
$$= \Delta_z^{\sigma_{\mu}} (z_0^{\tau_0} \psi(z_1)^{\tau_1} \psi(z_2)^{\tau_2})(0) = b_{\tau,\sigma_{\mu},0} + c_{\tau,\sigma_{\mu},0}$$

with

$$b_{\tau,\sigma_{\mu},0} = \frac{1}{\sigma_{\mu,1}!\sigma_{\mu,2}!} \left(\frac{\partial}{\partial t_1}\right)^{\sigma_{\mu,1}} \circ \left(\frac{\partial}{\partial t_2}\right)^{\sigma_{\mu,2}} (\beta z(t_1) - z(t_2))^{\tau_0} \left(\frac{\omega(t_1)}{t_1}\right)^{\tau_1} \left(\frac{\omega(t_2)}{t_2}\right)^{\tau_2} (0,0)$$

with exact order $|\sigma_{\mu}| = \sigma_{\mu,1} + \sigma_{\mu,2}$ for $b_{\tau,\sigma_{\mu},0}$. The term $c_{\tau,\sigma_{\mu},0}$ is a sum of the derivatives in (t_1, t_2) of order strictly inferior to $|\sigma_{\mu}|$.

Lemma 4.5 If $s \neq 0$, then we have

$$\begin{split} n_{\tau,\sigma_{\mu},s_{\mu}} &= \frac{1}{\sigma_{\mu,1}!\sigma_{\mu,2}!\psi(s_{\mu}u_{1})^{T_{1}}\psi(s_{\mu}u_{2})^{T_{2}}} \left(\frac{\partial}{\partial t_{1}}\right)^{\sigma_{\mu,1}} \circ \left(\frac{\partial}{\partial t_{2}}\right)^{\sigma_{\mu,2}} \left(F(z(t_{1}),z(t_{2}))\right)|_{t=0} \\ &= d_{\tau,\sigma_{\mu},s_{\mu}} + e_{\tau,\sigma_{\mu},s_{\mu}} \end{split}$$

with

$$F(z(t_1), z(t_2)) = (\beta z(t_1) - z(t_2))^{\tau_0} (\psi(z(t_1) - \psi(s_\mu u_1))^{2\tau_1}$$

$$\times (\psi(z(t_2)) - \psi(s_\mu u_2))^{2\tau_2} T(z(t_1), s_\mu u_1)^{T_1 - \tau_1} T(z(t_2), s_\mu u_2)^{T_2 - \tau_2}$$

and

$$d_{\tau,\sigma_{\boldsymbol{\mu}},s_{\boldsymbol{\mu}}} = \frac{1}{\sigma_{\boldsymbol{\mu},1}!\sigma_{\boldsymbol{\mu},2}!\psi(s_{\boldsymbol{\mu}}u_1)^{T_1}\psi(s_{\boldsymbol{\mu}}u_2)^{T_2}} \left(\frac{\partial}{\partial t_1}\right)^{\sigma_{\boldsymbol{\mu},1}} \circ \left(\frac{\partial}{\partial t_2}\right)^{\sigma_{\boldsymbol{\mu},2}} \left(F(z(t_1),z(t_2))\right)|_{t=0}$$

with exact order $|\sigma_{\mu}| = \sigma_{\mu,1} + \sigma_{\mu,2}$ for $d_{\tau,\sigma_{\mu},s_{\mu}}$. The term $e_{\tau,\sigma_{\mu},s_{\mu}}$ is a sum of the derivatives in (t_1,t_2) of order strictly inferior to $|\sigma_{\mu}|$.

Lemma 4.6 Put further

$$\ell_{\tau,\sigma_{\mu},s_{\mu}} = \begin{cases} b_{\tau,\sigma_{\mu},0} & (if \, s_{\mu} = 0) \\ d_{\tau,\sigma_{\mu},s_{\mu}} & (if \, s_{\mu} \neq 0) \end{cases}$$

$$(2)$$

and the new matrix

$$\mathcal{B} := (\gamma_{ au,\mu}) := (\ell_{ au,\sigma_{\mu},s_{\mu}}).$$

Then we have $\det \mathcal{B} = \det \mathcal{A} = \det \mathcal{N} = \Delta$.

Now we are going to give a lower bound for the height of $\Delta = \det(\gamma_{\tau,\mu}) = \det(a_{\tau,\mu}) \neq 0$. We do not use the differential equation, as is done in [20]. We have then next Lemma to estimate the height of each $\gamma_{\tau,\mu}$.

Lemma 4.7 Consider $\gamma_{\tau,\mu}$ namely either $b_{\tau,\sigma_{\mu},0}$ or $d_{\tau,\sigma_{\mu},s_{\mu}}$. Then we have

$$h(\gamma_{\tau,\mu}) = \begin{cases} h(b_{\tau,\sigma_{\mu},0}) \le 3.4S_0 \log(T_0+1) + T_0h_3 + 6T_0 + (S_0+1)(18+h_4) \\ +3(T_1+T_2), \\ h(d_{\tau,\sigma,s}) \le 3.4S_0 \log(T_0+1) + T_0h_3 + 6T_0 + (S_0+1)(18+h_4) \\ +(T_1+T_2)(16h_4+60\log 2+12) + 8S_1^2(T_1h_1+T_2h_2). \end{cases}$$
(3)

Lemma 4.8 We have

$$h(\Delta) \le \log(L!) + L\Big(T_0(h_3+6) + 3.4S_0\log(T_0+1) + (S_0+1)(18+h_4) + 8S_1^2(T_1h_1 + T_2h_2) + (T_1+T_2)(16h_4+60\log 2+12)\Big).$$

By means of the Liouville inequality, we can complete the proof of Lemma 4.1.

5 Extrapolation

It is possible to prove;

Lemma 5.1 Let Δ be an $L \times L$ minor determinant of \mathcal{M} . Suppose

$$|\Lambda| \le \exp\left(-\frac{L}{S_0}\log E\right) = E^{-L/S_0}.$$

Then we have

$$|\Delta| \le \exp\left(-\frac{L}{2}\left(\frac{L}{S_0} - 2S_0 + 1\right)\log E\right). \tag{4}$$

Lemma 5.2 Assume that there exist $T_0, T_1, T_2, S_0, S_1 \in \mathbb{Z} \geq 0$ with the following conditions. $S_0 \geq 5$, $S_0 - 1 \in 3\mathbb{Z}$, $S_1 \in 3\mathbb{Z}$, $(S_0 + 2)(S_0 + 5)(S_1 + 3) > 2916T_0T_1T_2$, $(S_0 + 2)(S_0 + 5)(T_1 + T_2) > 324T_0T_1T_2$, $(S_0 + 2)(S_1 + 3) > 81 \max\{T_1, T_2\}$, $(S_0 + 2)(T_1 + T_2) > 27T_1T_2$, $S_0 + 2 > 9T_0$. Assume further $\frac{L}{S_0} > \frac{D}{\log E} \times (2Q) + 2S_0 - 1$, $\frac{L}{S_0} > \frac{D}{\log E} \times R$ with $Q = \frac{\log(L!)}{L} + T_0(h_3 + 6) + 3.4S_0 \log(T_0 + 1) + (S_0 + 1)(18 + h_4) + 8S_1^2(T_1h_1 + T_2h_2) + (T_1 + T_2)(16h_4 + 60 \log 2 + 12)$ and with $R = 2\Omega_0^2T_2(\frac{T_1h_1}{4} + T_2h_2) + \frac{13}{2}h_4 + \frac{53}{2}\log 2$. Now suppose

 $|\Lambda| \le \exp\left(-\frac{L}{S_0}\log E\right).$

Then we have $\Lambda = 0$.

We have to choose parameters to achieve the proof of the main theorem. We have $L \leq (T_0 + 1)(T_1 + 1)(T_2 + 1)$.

Put $T_0 = [c_0 a_1 a_2 g^3 d^5]$, $T_1 = [c_1 a_2 b g d^3]$, $T_2 = [c_2 a_1 b g d^3]$, $S_0 = 1 + 3[c_3 a_1 a_2 b g^2 d^5]$, $S_1 = 3[c_4 g d]$, with absolute constants c_0, c_1, c_2, c_3, c_4 .

Since the quantity Q only differs from the assumptions in [20], thanks to the calculations due in [20], it is sufficient to choose;

 $c_0 = 2.13 \times 10^{28}$, $c_1 = c_2 = 9.85 \times 10^{16}$, $c_3 = 6.09 \times 10^{28}$, $c_4 = 5.50 \times 10^6$. Thus we complete the proof of our main theorem since

$$\frac{(1+c_0)(1+c_1)(1+c_2)}{3c_3-2} \le 1.16 \times 10^{35}$$

and

$$\frac{18c_0}{3c_3}\sqrt{c_2 \max\left(\frac{c_1}{4}, c_2\right)} \le 5.89 \times 10^{17}.$$

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