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Helmut MAIER et Michael Th. RASSIAS

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# Prime Avoidance Property of $k$ -th Powers of Piatetski–Shapiro Primes

par HELMUT MAIER et MICHAEL TH. RASSIAS

RÉSUMÉ. Dans des articles antérieurs, les auteurs ont établi la propriété d'évitement des nombres premiers pour les puissances  $k$ -ièmes de nombres premiers et pour les nombres premiers dans les suites de Beatty. Dans cet article, les auteurs considèrent les puissances  $k$ -ièmes des nombres premiers de Piatetski–Shapiro.

ABSTRACT. In previous papers the authors established the prime avoidance property of  $k$ -th powers of prime numbers and of prime numbers within Beatty sequences. In this paper the authors consider  $k$ -th powers of Piatetski–Shapiro primes.

## 1. Introduction

Let  $p_n$  denote the increasing sequence of prime numbers.

The fact that

$$\limsup_{n \rightarrow \infty} \frac{p_{n+1} - p_n}{\log p_n} = \infty$$

was first proved by Westzynthius [15]. Erdős [2] obtained that infinitely often one has:

$$(1.1) \quad p_{n+1} - p_n > C_1 \frac{\log p_n \log_2 p_n}{(\log_3 p_n)^2}$$

with appropriate  $C_1 > 0$ . Here  $\log_k x := \log(\log_{k-1} x)$ .

The order of magnitude of the lower bound (1.1) was slightly improved by Rankin [13] who showed that

$$(1.2) \quad p_{n+1} - p_n > C_2 \frac{\log p_n \log_2 p_n \log_4 p_n}{(\log_3 p_n)^2} = C_2 g_1(p_n).$$

The proof of Rankin differs only slightly from that of Erdős. Rankin is using a better estimate for the number of smooth integers. The approach of Rankin and modifications of it have become known as the Erdős–Rankin method. In [2] and [13] these two authors for sufficiently large  $x$  construct a long interval of integers, which all have greatest common divisor greater than 1 with  $\prod_{p < x} p$ .

The results (1.1) and (1.2) are closely related to the concept of prime-avoidance. The term “prime avoidance” first appears in the paper [5] of K. Ford, D. R. Heath-Brown and S. Konyagin, where they prove the existence of infinitely many “prime-avoiding” perfect  $k$ -th powers for any positive integer  $k$ . Their definition is as follows:

**Definition 1.1.** *Let*

$$(1.3) \quad g_1(m) := \frac{\log m \log_2 m \log_4 m}{(\log_3 m)^2}.$$

*Then an integer  $m$ , for which  $g_1(m) \geq m^{(0)}$  (with  $m^{(0)} > 0$  being a fixed constant) is called prime-avoiding with constant  $c_0 > 0$  if  $m+u$  is composite for all integers  $u$  satisfying  $|u| \leq c_0 g_1(m)$ . We also say,  $m$  has the prime avoidance property with constant  $c_0$ .*

For the sake of larger flexibility we also consider the following variations:

**Definition 1.2.** *Let  $h(m)$  be defined for  $m \geq m^{(1)}$  and let  $h(m) \rightarrow \infty$  for  $m \rightarrow \infty$ . Let  $c_1 > 0$  be a fixed constant.*

(i) *Then an integer  $m \geq m^{(1)}$  (with  $m^{(1)} > 0$  being a fixed constant) is called prime-avoiding with constant  $c_1 > 0$  and function  $h$  if  $h(m) \geq m^{(1)}$  and  $m+u$  is composite for all integers  $u$  satisfying  $|u| \leq c_1 h(m)$ . We also say, that  $m$  has the prime avoidance property with constant  $c_1$  and function  $h$ .*

(ii)  *$m \geq m^{(1)}$  is called prime-avoiding to the right with function  $h$  and constant  $c_2$  if  $m+u$  is composite for all integers  $u \in [0, c_2 h(m)]$ . We say that  $m$  has the prime-avoidance property to the right.*

*An analogous definition applies for “prime-avoiding to the left” and “prime avoidance property to the left”.*

(iii) *The terms “one-sided prime-avoidance property” is used, if prime-avoidance property to the right, to the left or both hold.*

The function  $g_1$  in (1.3) has played a central role for a long time in the study of large gaps between consecutive primes. The connection between these questions can easily be seen as follows:

Let  $m$  be prime-avoiding with function  $h$  and constant  $c_1 > 0$  and let  $p_n$  be the largest prime number  $\leq m$  with  $p_{n+1}$  being the subsequent prime number. Then

$$p_{n+1} - p_n \geq 2c_1 h(m).$$

Also the reverse (large gaps imply the existence of prime-avoiding numbers) can be seen easily.

The results (1.1) and (1.2) also imply prime-avoidance results for most integers of the interval  $(p_n, p_{n+1})$ .

A modified version of this method was applied by the authors of [5] on the problem of prime-avoiding  $k$ -th powers mentioned in the introduction.

In [10] the authors of the present paper extended this result by proving the existence of infinitely many prime-avoiding  $k$ -th powers of prime numbers. For the sake of simplicity they only treat one-sided prime avoidance. They prove the following:

*There are infinitely many  $n$ , such that*

$$p_{n+1} - p_n \geq C_3 g_1(n)$$

*and the interval  $[p_n, p_{n+1}]$  contains the  $k$ -th power of a prime number.*

Their method of proof consists of a combination of the method of [5] with the matrix method of the first author [8]. The matrix  $\mathcal{M}$  employed in this technique is of the form

$$\mathcal{M} := (a_{r,u} : u \in \mathcal{B}) ,$$

where  $P(x) := \prod_{p < x} p$ ,  $D$  is a fixed positive integer, the rows

$$\mathcal{R}(r) := \{a_{r,u} : u \in \mathcal{B}\}$$

of the matrix are *translates* (in closer or wider sense) of the *base-row*  $\mathcal{B}$ .

The essential idea of the construction of the base-row  $\mathcal{B}$  is the Erdős–Rankin method. Here however (as in all combinations with the matrix method) the base-row is partially coprime to  $P(x)$ . The columns

$$\mathcal{C}(u) := (a_{r,u})_{1 \leq r \leq P(x)^{D-1}} , \quad \text{with } (a_{r,u}, P(x)) = 1$$

are called *admissible columns*.

The number of prime numbers in these columns, which are arithmetic progressions (or in the case of [10], shifted powers of elements of arithmetic progressions), can be estimated using theorems on primes in arithmetic progressions.

A famous prize problem of Erdős, being open for more than 70 years, was to replace the function  $g_1$  in (1.2) by a function of higher order of magnitude. This problem could finally be solved in the paper [4] by K. Ford, B. J. Green, S. Konyagin and T. Tao and independently in the paper [12] by J. Maynard. Later all five authors improved on this result in their joint paper [3]. They proved:

$$p_{n+1} - p_n \geq C_4 g_2(p_n)$$

infinitely often, where

$$g_2(m) := \frac{\log m \log_2 m \log_4 m}{\log_3 m} .$$

In the paper [10] the authors of the present paper combined the methods of the papers [3, 5, 8] to obtain the following theorem.

**Theorem 1.1 of [10].** *There is a constant  $C_5 > 0$  and infinitely many  $n$  such that*

$$p_{n+1} - p_n \geq C_5 g_2(n)$$

*and the interval  $[p_n, p_{n+1}]$  contains the  $k$ -th power of a prime.*

Using the Definition 1.2 this can also be phrased as follows:

There is a constant  $C_5 > 0$  such that infinitely many  $k$ -th powers of primes are one-sided prime-avoidant with constant  $\frac{1}{2}C_5$ . In the paper [11] the authors investigated the prime-avoidance of the  $k$ -th powers of prime numbers with Beatty sequences.

**Definition 1.3.** *For two fixed real numbers  $\alpha$  and  $\beta$ , the corresponding non-homogeneous Beatty sequence is the sequence of integers defined by*

$$\mathcal{B}_{\alpha, \beta} := (\lfloor \alpha n + \beta \rfloor)_{n=1}^{\infty}$$

*( $\lfloor u \rfloor$  denotes the largest integer  $\leq u$ ).*

**Definition 1.4.** *For an irrational number  $\gamma$  we define its type  $\tau$  by the relation*

$$\tau := \sup\{\rho \in \mathbb{R} : \liminf n^\rho \|\gamma n\| = 0\}$$

*(see [11]).  $\|\cdot\|$  denotes the distance to the nearest integer.  $\tau$  is called of finite type, if  $\tau < \infty$ .*

In the paper [11] the authors prove:

**Theorem 1.3 of [11].** *Let  $k \geq 2$  be an integer. Let  $\alpha, \beta$  be fixed real numbers with  $\alpha$  being a positive irrational and of finite type. Then there is a constant  $C_6 > 0$ , depending only on  $\alpha$  and  $\beta$ , such that for infinitely many  $n$  we have:*

$$p_{n+1} - p_n \geq C_6 g_2(n)$$

*and the interval  $[p_n, p_{n+1}]$  contains the  $k$ -th power of a prime  $\tilde{p} \in \mathcal{B}_{\alpha, \beta}$ .*

In this paper we deal with the Piatetski-Shapiro primes. We prove:

**Theorem 1.5.** *Let  $c_4 \in (1, 18/17)$  be fixed,  $k \in \mathbb{N}$ ,  $k \geq 2$ . Then there is a constant  $C_6 > 0$ , depending only on  $k$  and  $c_4$ , such that for infinitely many  $n$  we have*

$$p_{n+1} - p_n \geq C_6 g_2(n)$$

*and the interval  $[p_n, p_{n+1}]$  contains the  $k$ -th power of a prime  $\tilde{p} = \lfloor l^{c_4} \rfloor$ .*

Using Definition 1.2 we can phrase Theorem 1.5 as follows:

**Theorem 1.6.** *Let  $c_4 \in (1, 18/17)$  be fixed,  $k \in \mathbb{N}$ ,  $k \geq 2$ . Then there is a constant  $C_7 > 0$  such that infinitely many  $k$ -th powers of primes of the form  $\tilde{p} = \lfloor l^{c_4} \rfloor$  have the one-sided prime avoidance property with constant  $C_7$  and function  $g_2$ .*

One can achieve prime avoidance by a slight modification of the proof of Theorem 1.5.

## 2. Construction of the matrix $\mathcal{M}$

In several papers (e.g. [8, 9, 10, 11]) in which the matrix method was used, it was crucial for the estimate of prime numbers in arithmetic progressions mod  $q$  that  $q$  was a “good modulus”.

The definitions and facts of these papers may be summarised as follows: We recall the definition:

**Lemma 2.1.** *There exists a constant  $C_8 > 0$  such that for every  $C > C_8$  there is a sequence of real numbers  $(x_n(C))$ , depending only on  $C$ , such that for every  $n$  the modulus  $q_n := P(x_n(C))$  is good with respect to  $C$  in the following sense: One has  $L(s, \chi) \neq 0$  for every character  $\chi$  modulo  $q_n$  and every  $s = \sigma + it$  satisfying*

$$(2.1) \quad \sigma > 1 - \frac{C}{\log[q(|t| + 1)]}.$$

**Lemma 2.2.** *Let  $C > 0$  and the sequence  $(x_n) = (x_n(C))$  be given, such that (2.1) is satisfied. Then there exists a constant  $D > 0$  depending only on  $C$ , such that*

$$(2.2) \quad \pi(x_n; q_n, a) \gg \frac{x_n}{\phi(q_n) \log x_n}$$

*uniformly for  $(a, q_n) = 1$  and  $x_n \geq q^{D/2}$ .*

*The constant  $D$  can be made arbitrarily large, if  $C$  is sufficiently large.*

*Proof.* This is Lemma 2 of [8]. It is not due to Gallagher (as stated in [8]), but is a simple consequence of Theorem 7 of Gallagher [6].  $\square$

In the sequel we fix constants  $C, D$ , such that  $D$  is sufficiently large and a sequence  $(x_n) = (x_n(C))$ , such that (2.1) and (2.2) are satisfied.

**Lemma 2.3.** *Let  $x$  be sufficiently large.*

$$(2.3) \quad y = C_9 x \frac{\log x \log_3 x}{\log_2 x},$$

$C_9 > 0$  being a sufficiently small constant. Let  $C_{10} > 0$  be a fixed constant to be specified later. Then there is an integer  $m_0$  satisfying

$$(2.4) \quad 1 \leq m_0 < P(C_{10}x)$$

$$(2.5) \quad (m_0 + 1, P(C_{10}x)) = 1$$

*and an exceptional set  $V$  satisfying*

$$(2.6) \quad \#V \ll x^{1/2+\epsilon}, \quad \epsilon > 0 \text{ arbitrarily small}$$

*such that*

$$(2.7) \quad (m_0 + 1 + rP(C_0x))^k + u - 1$$

*is composite for  $2 \leq u \leq y$ , unless  $u \in V$ .*

*Proof.* By Lemma 3.10 of [10], there is an integer  $m_0$  satisfying all the properties (2.3)–(2.7) with the possible exception of (2.5).

The additional property (2.5) follows from the construction of  $m_0$  in [10]. In Definition 3.9 of [10],  $m_0$  is defined by

$$1 \leq m_0 < P(C_{10}x) \quad (C_{10} \text{ being named } C_0 \text{ in [10]})$$

and the congruences

$$\begin{aligned} m_0 &\equiv c_s \pmod{s} \\ m_0 &\equiv d_p \pmod{p} \\ (2.8) \quad m_0 &\equiv 0 \pmod{q}, \quad q \in (1, x], \quad q \notin S \cup P \\ m_0 &\equiv e_u \pmod{p_u}, \\ m_0 &\equiv g_p \pmod{p}, \text{ for all other primes } p \leq C_{10}x, \quad g_p \text{ being arbitrary.} \end{aligned}$$

Since by the definition of  $\mathcal{A}, \mathcal{B}$  we have

$$c_s \not\equiv -1 \pmod{s}, \quad d_p \not\equiv -1 \pmod{p}$$

we have from the first four congruences in (2.8) that  $m_0 + 1 \not\equiv 0$  for all  $p \in [2, C_{10}x]$  and thus  $(m_0 + 1, P(C_{10}x)) = 1$ .  $\square$

The construction of  $V$  is described in Lemma 3.5 of [10]. We now construct the matrix  $\mathcal{M}$ . We start with the following:

*Remark.* In the proof of Lemma 3.10 in [10], prime numbers from the interval  $(x, C_{10}x]$  are used to construct a number  $m_0$  with  $1 \leq m_0 < P(C_{10}x)$ , such that

$$a_{r,u} = (m_0 + 1 + rP(x))^k + u - 1,$$

$2 \leq u \leq y$ , is composite unless  $u \in V$ . For the investigation of primes in the columns of the matrix  $\mathcal{M}$  the condition needed is that the moduli  $P(C_{10}x)$  are good. This is transformed into the condition that  $P(x)$  is good by renaming  $x$  into  $x/C_{10}$ .

**Definition 2.4.** Let  $n \in \mathbb{N}$  be sufficiently large and let  $x = x_n = x_n(C)$ . The existence has been established in Lemma 2.1. Let  $q = q_n = P(x_n(C))$ . Let  $y$  satisfy (2.3), and  $m_0, V$  satisfy (2.4)–(2.7) with  $C_{10} = 1$ . We let

$$\mathcal{M} := (a_{r,u})_{\substack{1 \leq r \leq P(x)^{D-1} \\ 1 \leq u \leq y}},$$

where  $a_{r,u} := (m_0 + 1 + rP(x))^k + u - 1$

For  $1 \leq r \leq P(x)^{D-1}$  we denote by

$$R(r) := (a_{r,u})_{1 \leq u \leq y}$$

the  $r$ -th row of  $\mathcal{M}$  and for  $1 \leq u \leq y$  we denote by

$$C(u) := (a_{r,u})_{1 \leq r \leq P(x)^{D-1}}$$

the  $u$ -th column of  $\mathcal{M}$ .

### 3. Piatetski-Shapiro primes in arithmetic progressions

**Definition 3.1.** Let  $c_6 > 1$ . Let  $a$  and  $d$  be coprime integers,  $w \geq 1$ . Then we let

$$\pi_{c_6}(w; d, a) = \#\{p \leq w : p \in \mathcal{P}^{(c_6)}, p \equiv a \pmod{d}\}$$

$$\mathcal{P}^{(c_6)} = \{p \text{ prime} : p = \lfloor l^{c_6} \rfloor \text{ for some } l\}.$$

**Lemma 3.2.** Let  $a$  and  $d$  be coprime integers,  $d \geq 1$ . For fixed  $c_6 \in (1, 18/17)$  we have (with  $\gamma = 1/c_6$ ):

$$\pi_{c_6}(w; d, a) = \gamma w^{\gamma-1} \pi(w; d, a) + \gamma(1 - \gamma) \int_2^w u^{\gamma-2} \pi(u; d, a) du$$

$$+ O\left(w^{\frac{17}{39} + \frac{7\gamma}{13} + \epsilon}\right),$$

with the constant being implied in  $O$  not depending on  $d$  or  $a$ .

*Proof.* This is Theorem 8 of [1]. □

We shall use Dirichlet-characters to evaluate  $\pi(u; P(x), m_0 + 1)$  and use the familiar definition

$$\theta(u; g, r) = \sum_{\substack{p \leq u \\ p \equiv r \pmod{g}}} \log p.$$

**Lemma 3.3.** Assume that  $\chi \pmod{r}$  is induced by the primitive character  $\chi^*$ . Then

$$\sum_{p \leq u} |\chi(p) \log p - \chi^*(p) \log p| = O(\log^2(ru)).$$

*Proof.*  $\chi(p) \neq \chi^*(p)$  implies  $\chi(p) = 0$  and thus  $p \mid \text{conductor}(\chi) \mid r$ . Also

$$\sum_{p|r} 1 = O(\log r).$$

Thus

$$\sum_{p \leq u} |\chi(p) \log p - \chi^*(p) \log p| = O(\log u \cdot \log r) = O(\log(ur)^2). \quad \square$$

**Definition 3.4.** We let  $l(r) := m_0 + 1 + rP(x)$ ,

$$\mathcal{R}_1(\mathcal{M}) := \{r : 1 \leq r \leq P(x)^{D-1}, l(r) \in \mathcal{P}^{(c_6)}\},$$

$$\mathcal{R}_2(\mathcal{M}) := \left\{r : \begin{array}{l} 1 \leq r \leq P(x)^{D-1}, \\ R(r) \text{ contains a prime number} \end{array} \right\}.$$

We observe that each row  $R(r)$  with  $r \in \mathcal{R}_1(\mathcal{M})$  has as its first element  $a_{r,1}$ , the  $k$ -th power of the prime  $l(r) \in \mathcal{P}^{(c_6)}$ . We now conclude the proof of Theorem 1.5 by showing that the set  $\mathcal{R}_1(\mathcal{M}) \setminus \mathcal{R}_2(\mathcal{M})$  is non-empty.



**Lemma 3.5.** *We have*

$$\#\mathcal{R}_1(\mathcal{M}) = \frac{P(x)^{D\gamma}}{\phi(P(x))} \left(1 + O\left(e^{-c_7 D}\right)\right).$$

*Proof.* We apply Lemma 3.2 with  $P(x)^D$  instead of  $w$ ,  $d = P(x)$ ,  $a = m_0 + 1$  and obtain

$$\begin{aligned} (3.1) \quad \#\mathcal{R}_1(\mathcal{M}) &= \gamma P(x)^{D(\gamma-1)} \pi(P(x)^D; P(x), m_0 + 1) \\ &\quad + \gamma(1 - \gamma) \int_2^{P(x)^D} u^{\gamma-2} \pi(u; P(x), m_0 + 1) \, du \\ &\quad + O\left(P(x)^{D(\frac{17}{39} + \frac{7\gamma}{13} + \epsilon)}\right). \end{aligned}$$

We now use Dirichlet-characters to evaluate  $\pi(u; P(x), m_0 + 1)$ . We use the familiar definition

$$\theta(u; g, r) = \sum_{\substack{p \leq u \\ p \equiv r \pmod{g}}} \log p$$

and obtain

$$\begin{aligned} (3.2) \quad \theta(u; P(x), m_0 + 1) &= \frac{1}{\phi(P(x))} \sum_{p \leq u} \log p + \frac{1}{\phi(P(x))} \sum_{\substack{\chi \pmod{P(x)} \\ \chi \neq \chi_0}} \sum_{p \leq u} \chi(p) \log p. \end{aligned}$$

From Lemma 3.3 we obtain

$$\begin{aligned} &\left| \theta(u; P(x), m_0 + 1) - \frac{1}{\phi(P(x))} \sum_{p \leq u} \log p \right| \\ &\leq \frac{1}{\phi(P(x))} \sum_{1 < r \leq P(x)} \sum_{\chi^* \pmod{r}}^* \left( \sum_{p \leq u} \chi^*(p) \log p + O((\log P(x))^2) \right), \end{aligned}$$

where  $\sum^*$  denotes the summation over primitive characters.

We first assume that  $u > P(x)^{D/2}$  and apply Lemma 2.2 and obtain:

$$(3.3) \quad \theta(u; P(x), m_0 + 1) = \frac{1}{\phi(P(x))} \left( \sum_{p \leq u} \log p \right) \left( 1 + O\left(e^{-c_7 D}\right) \right).$$

For  $u \leq P(x)^{D/2}$  we use the trivial estimate

$$\theta(u; P(x), m_0 + 1) = O(u \log^2 u).$$

Lemma 3.5 now follows from (3.1), (3.2) and (3.3).  $\square$

We now estimate  $\mathcal{R}_2(\mathcal{M})$ . We have

$$(3.4) \quad |\mathcal{R}_2(\mathcal{M})| \leq \sum_{v \in V} t(v),$$

where

$$(3.5) \quad t(v) = \#\left\{r : l(r) \in \mathcal{P}^{(c_6)}, l(r)^k + v - 1 \text{ prime}\right\}.$$

We majorize the set  $\mathcal{P}^{(c_6)}$  by sieving the simpler set

$$\mathcal{N}^{(c_6)} = \{n \in \mathbb{N} : n = \lfloor l^{c_6} \rfloor \text{ for some } l\}.$$

**Definition 3.6.** Let

$$s_c(x; h, i) = \#\left\{n \in \mathcal{N}^{(c_6)} : n \leq x, n \equiv i \pmod{h}\right\}.$$

**Lemma 3.7.**

$$s_c(x; h, i) = \frac{x^{1/c}}{h} + O_c\left(\left(x^{1+5/c}(\log x)^3\right)^{1/7}\right).$$

*Proof.* This is Proposition 1 of [14]. □

Let

$$Q(x) = \prod_{\substack{x < p \leq P(x) \\ p \text{ prime}}} p.$$

Let

$$(3.6) \quad n(r) = m_0 + 1 + rP(x)$$

$$N(v) := \#\left\{r : \begin{array}{l} n(r) \in \mathcal{N}^{(c_6)}, \quad n(r) \leq 2P(x)^D, \\ (n(r), P(x)) = 1, \quad (n(r)^k + v - 1, Q(x)) = 1 \end{array}\right\}.$$

By the sieve of Eratosthenes we have:

$$(3.7) \quad N(v) \leq \sum_{\substack{r : n(r) \in \mathcal{N}^{(c_6)}, \\ n(r) \leq 2P(x)^D}} \sum_{\tau_1 | (n(r), P(x))} \mu(\tau_1) \sum_{\tau_2 | (n(r)^k + v - 1, Q(x))} \mu(\tau_2).$$

We now apply a Fundamental Lemma in the theory of Combinatorial Sieves.

**Lemma 3.8.** Let  $\kappa > 0$  and  $\Omega > 1$ . There exist two sets of real numbers  $\Lambda^+ = (\lambda_d^+)$  and  $\Lambda^- = (\lambda_d^-)$  depending only on  $\kappa$  and  $\Omega$  with the following properties:

$$(3.8) \quad \lambda_1^\pm = 1,$$

$$(3.9) \quad |\lambda_d^\pm| < 1, \quad \text{if } 1 < d < \Omega,$$

$$(3.10) \quad \lambda_d^\pm = 0, \quad \text{if } d \geq \Omega, \text{ and for any integer } n > 1,$$

$$(3.11) \quad \sum_{d|n} \lambda_d^- \leq 0 \leq \sum_{d|n} \lambda_d^+.$$

Moreover, for any multiplicative function  $g(d)$  with  $0 \leq g(p) < 1$  and satisfying the dimension condition

$$\prod_{w \leq p < z} (1 - g(p))^{-1} \leq \left( \frac{\log z}{\log w} \right)^\kappa \left( 1 + \frac{K}{\log w} \right)$$

for all  $w, z$  such that  $2 \leq w < z \leq y$  we have

$$\sum_{d|P(z)} \lambda_d^\pm g(d) = \left( 1 + O \left( e^{-s} \left( 1 + \frac{K}{\log z} \right)^{10} \right) \right) \prod_{p < z} (1 - g(p)),$$

where  $s = \frac{\log \Omega}{\log z}$ .

*Proof.* This is Fundamental Lemma 6.3 in [7, p. 159].  $\square$

For  $x < \tilde{p} \leq P(x)$  let  $\rho(v, \tilde{p})$  be the number of solutions (in  $s$ ) of the congruence

$$s(s^k + v - 1) \equiv 0 \pmod{\tilde{p}}.$$

We apply Lemma 3.8 with  $g$  given by

$$g(\tilde{p}) = \begin{cases} 1/\tilde{p}, & \text{for } 1 < \tilde{p} \leq x \\ \rho(v, \tilde{p})/\tilde{p} & \text{for } x < \tilde{p} \leq P(x). \end{cases}$$

The range of definition of  $g$  is extended to all integers by multiplicativity.

- $\lambda_d^-$  is not needed.
- $\lambda_d^+$  is chosen satisfying (3.8)–(3.11), with  $\Omega = P(x)^{D/100}$ .

From (3.7) and Lemma 3.8, we obtain

$$(3.12) \quad N(v) \leq \sum_{\tau|P(x)} \sum_{\substack{n \in \mathcal{N}^{c_6}, n \leq 2P(x)^D, \\ n \equiv m_0 + 1 \pmod{P(x)}, \\ n(n^k + v - 1) \equiv 0 \pmod{\tau}}} 1$$

Let  $S(\tau, v)$  be the solution set of  $s(s^k + v - 1) \equiv 0 \pmod{\tau}$ . The system of congruences in the inner sum of (3.12) is equivalent to a union of congruence classes

$$n \equiv z(s) \pmod{P(x)}, \quad s \in S(\tau, v).$$

Thus

$$(3.13) \quad N(v) \leq \sum_{\tau|P(x)} \lambda^+(\tau) \sum_{\substack{n \in \mathcal{N}^{c_6}, n \leq 2P(x)^D \\ n \equiv z(s) \pmod{P(x)}}} 1$$

The inner sum in (3.13) can now be evaluated by the use of Lemma 3.7.

From (3.4), (3.5), (3.6), (3.7) and (3.13) it follows that

$$\#\mathcal{R}_2 = o(\#\mathcal{R}_1),$$

which finishes the proof of Theorem 1.3.  $\square$

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Helmut MAIER  
 Department of Mathematics, University of Ulm  
 Helmholtzstrasse 18  
 89081 Ulm, Germany  
*E-mail:* [helmut.maier@uni-ulm.de](mailto:helmut.maier@uni-ulm.de)

Michael Th. RASSIAS  
 Department of Mathematics and Engineering Sciences  
 Hellenic Military Academy  
 16673 Vari Attikis, Greece  
*E-mail:* [mthrassias@yahoo.com](mailto:mthrassias@yahoo.com)