# Gap distribution among sequences arising in the theory of modular forms 

A thesis<br>submitted in partial fulfillment of the requirements of the degree of<br>Doctor of Philosophy

By

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20123233

## Declaration

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Certified that the work incorporated in the thesis entitled "Gap distribution among sequences arising in the theory of modular forms" submitted by Sudhir Kumar Pujahari was carried out by the candidate, under my supervision. The work presented here or any part of it has not been included in any other thesis submitted previously for the award of any degree or diploma from any other University or institution.

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## Chapter 1

## Introduction

The rich story of equidistribution started in the years 1909-1910 by the work of P. Bohl [5], H. Weyl [50] and W. Sierpinski [43] where they studied the distribution of the sequence $\{n \alpha\}(\bmod 1)$, (for an irrational $\alpha)$ in the unit interval. Let us recall that a sequence of real numbers $\left\{x_{n}\right\}$ lying in the interval $[0,1] \subseteq \mathbb{R}$ is said to be uniformly distributed or equidistributed with respect to Lebesgue measure if for any interval $[\alpha, \beta] \subseteq[0,1]$, we have

$$
\lim _{V \rightarrow \infty} \frac{1}{V} \#\left\{n \leq V: x_{n} \in[\alpha, \beta]\right\}=\beta-\alpha
$$

where \# denotes the cardinality. This subject attracted great attention of mathematicians from all branches of mathematics after Hermann Weyl related the study of equidistribution to the study of exponential sums in his 1916 paper [51]. Our work is partly motivated by the following result of Van der Corput (see [27], page no. 176): If for each positive integer $s$, the sequence $\left\{x_{n+s}-x_{n}\right\}$ is uniformly distributed $(\bmod 1)$, then the sequence $\left\{x_{n}\right\}$ is uniformly distributed $(\bmod 1)$.
Let us consider the following classical question:
Question 1.0.1. Is the converse of Van der Corput's result true?

In other words, if $\left\{x_{n}\right\}$ is uniformly distributed $(\bmod 1)$, then is it true that for any positive integer $s$, the sequence $\left\{x_{n+s}-x_{n}\right\}$ is uniformly distributed $(\bmod 1)$ ?

The answer to the above question is surprisingly no. For example, consider the well-studied sequence $\{n \alpha\}(\bmod 1), \alpha$ is irrational. Write the sequence as follows:
For a natural number $N$, define,

$$
A_{\alpha}(N)=\{\alpha n(\bmod ) 1: 1 \leq n \leq N\} \subset\{\{n \alpha\}(\bmod 1): n \in \mathbb{N}\}
$$

and write them as increasing order as follows:

$$
\begin{equation*}
A_{\alpha}(N)=\left\{0<x_{1} \leq x_{2} \leq x_{3} \leq \cdots \leq x_{N}<1\right\} \tag{1.1}
\end{equation*}
$$

where $x_{N+1}=1+x_{1}$.
In 1957, Steinhaus conjectured the following fact:

$$
\#\left\{x_{i+1}-x_{i}: 1 \leq i \leq N\right\} \leq 3
$$

where \# denotes the cardinality of the set. There are several proofs of the above conjecture available in the literature, but the first proof was given by Vera Sós [44] and [45] in 1958. The above statement is popularly known as "The three gap theorem".
In 2002, Vâjâitu and Zaharescu [48] investigated the following question:
Question 1.0.2. Let $A_{\alpha}(N)$ be as defined in (1.1). Remove as many elements of $A_{\alpha}(N)$ as one likes. Then, how large is the cardinality of the consecutive differences of the resulting set?

More generally, they prove the following:
For any subset $\Omega(N)$ of $A_{\alpha}(N)$, there are no more than $(2+\sqrt{2}) \sqrt{N}$ distinct
consecutive differences, that is, if

$$
B(\Omega(M))=\left\{0 \leq y_{1} \leq y_{2} \leq \cdots \leq y_{M}<1 \leq y_{m+1}:=1+y_{1}\right\} \subset A_{\alpha}(N)
$$

then

$$
\#\left\{y_{i+1}-y_{i}: 1 \leq i \leq M, y_{i} \in \Omega(M)\right\} \leq 2 \sqrt{2 N}+1
$$

In 2015, using additive combinatorics A. Balog, A.Granville and J. Solymosi [2] improved the bound of above result [48] to $2 \sqrt{2 N}+1$ for any finite subset of $\mathbb{R} / \mathbb{Z}$, where $\mathbb{R}$ is the set of real numbers and $\mathbb{Z}$ is the set of integers. In particular for our concerned sequence, they proved that

$$
\# B(\Omega(N)) \leq 2 \sqrt{2 N}+1 .
$$

From the above result, we can conclude that for any subsequence say $\left\{y_{n}\right\}=\left\{y_{1}, y_{2}, y_{3}, \ldots\right\}$ (arranged in ascending order) of $\left\{x_{n}\right\}=\left\{x_{1}, x_{2}, x_{3}, \ldots\right\}$, the consecutive difference $\left\{y_{i+1}-y_{i}\right\}$ is not uniformly distributed $(\bmod 1)$. In this thesis we will show that, if a sequence is equidistributed in $\left[-\frac{1}{2}, \frac{1}{2}\right]$ with respect to a probability measure say $\mu=F(x) d x$ (for the definition see Section 2.1), then the fractional parts of gaps of all elements of the sequence will be equidistributed in $[0,1]$ with respect to the measure $F(x) * F(x) d x$, where $*$ is the convolution of the measures.
More explicitly, if we have two sequences say $\left\{x_{n}\right\}_{n=1}^{\infty}$ and $\left\{y_{m}\right\}_{m=1}^{\infty}$ such that they are equidistributed with respect to probability measures $\mu_{1}=$ $F_{1}(x) d x$ and $\mu_{2}=F_{2}(x) d x$ respectively in $\left[-\frac{1}{2}, \frac{1}{2}\right]$, then the sequence of fractional parts of gaps between elements of $\left\{x_{n}\right\}$ and $\left\{y_{m}\right\}$, that is, $\left\{x_{n}-\right.$ $\left.y_{m}\right\}_{n, m=1}^{\infty}(\bmod 1)$ is equidistributed in $[0,1]$ with respect to $F_{1}(x) * F_{2}(x) d x$. We are also able to predict quantitatively the rate of convergence of the following:

$$
\lim _{V \rightarrow \infty} \frac{1}{V^{2}} \#\left\{1 \leq m, n \leq V:\left\{x_{n}-y_{m}\right\} \bmod 1 \in[\alpha, \beta]\right\}
$$

where $[\alpha, \beta]$ is any subinterval of $[0,1]$, whenever $\left\{x_{n}\right\}_{n=1}^{\infty}$ and $\left\{y_{m}\right\}_{m=1}^{\infty}$ satisfy some conditions that have been described in Theorem 4.1.6. More generally, we have results for $r$ equidistributed sequences. These results are stated and proved in the Chapter 4 as Theorem 4.1.1, 4.1.4 and 4.1.6. We have discussed several applications of our results in Chapter 6.
Let $S(N, k)$ be the space of all holomorphic cusp forms of weight $k$ with respect to $\Gamma_{0}(N)$. For any positive integer $n$, let $T_{n}(N, k)$ be the $n^{\text {th }}$ Hecke operator acting on $S(N, k)$. Let $s(N, k)$ denote the dimension of the vector space $S(N, k)$. For a positive integer $n \geq 1$, let

$$
\left\{a_{n, i, N}, 1 \leq i \leq s(N, k)\right\}
$$

denote the eigenvalues of $T_{n}$, counted with multiplicity. For any positive integer $n$, let $T_{n}^{\prime}$ be the normalized Hecke operator acting on $S(N, k)$, defined as follows

$$
T_{n}^{\prime}:=\frac{T_{n}}{n^{\frac{k-1}{2}}} .
$$

Consider

$$
\left\{\frac{a_{n, i, N}}{n^{\frac{k-1}{2}}}, 1 \leq i \leq s(N, k)\right\},
$$

the eigenvalues of $T_{n}^{\prime}$ counted with multiplicity. Let $p$ be a prime number such that $p$ and $N$ are coprime. Then by the theorem of Deligne (see [14]) proving the Ramanujan-Petersson inequality, we know that

$$
a_{p, i, N} \in\left[-2 p^{\frac{k-1}{2}}, 2 p^{\frac{k-1}{2}}\right] .
$$

For each $i$, choose $\theta_{p, i, N} \in[0, \pi]$ such that

$$
\frac{a_{p, i, N}}{p^{\frac{k-1}{2}}}=2 \cos \theta_{p, i, N}
$$

The $\theta_{p, i, N}$ are called eigenangles.

In Chapter 2, we discuss all the preliminaries we require. In Chapter 3, we discuss Eichler-Selberg trace formula and Kuznetsov trace formula. In Chapter 4, we discuss the distribution of gaps between equidistributed sequences. In Chapter 5, we prove a variant of the Erdös-Turán Inequality. In Chapter 6, we study the distribution of gaps between eigenangles of Hecke operators. Using results of Murty and Sinha [32], Murty and Srinivas [33] have recently proved the following results

$$
\begin{gathered}
\#\left\{(i, j), 1 \leq i, j \leq s(N, k): \theta_{p, i, N} \pm \theta_{p, j, N}=0\right\} \\
=\mathrm{O}\left((s(N, k))^{2}\left(\frac{\log p}{\log k N}\right)\right)
\end{gathered}
$$

Note that taking $k$ and $N$ sufficiently large, the above result gives a little evidence towards the Maeda and Tsaknias conjectures.

For $\mathrm{N}=1$, the famous Maeda conjecture (see [19]) predicts that the polynomial $\prod_{i=1}^{s(N, k)}\left(x-a_{p, i, 1}\right)$ is irreducible over $\mathbb{Q}$. It also predicts that the Galois group of minimal polynomial of $T_{n}$ is the full symmetric group $S_{d}$, where $d$ is the dimension of $S(1, k)$. Based on computational data, Tsaknias (see [46]) predicts that for a fixed level $N>1$, the above polynomial is a product of bounded numbers of irreducible polynomials viewed as a function of $k$. In this thesis, we can get the measure with respect to which the differences of eigenangles of Hecke operators are equidistributed. As a special case to our result, we obtain the result of Murty and Srinivas. We could also get an error term (see Chapter 6, Theorems 6.1.1, 6.1.3). We also discuss similar results for primitive Maass forms (see Theorems 6.2.1, 6.2.2). In the case of primitive Maass forms, we have assumed the Ramanujan bound.

Chapter 7 is a report on recent joint work with M. Ram Murty, where we have given a multiplicity one theorem for Hecke eigenforms. We also dis-
cuss a prove of the Joint Sato-Tate conjecture. The contents of this chapter are independent of the content of the previous chapters.

## Chapter 2

## Preliminary concepts

The study of Fourier coefficients of modular forms is of great interest to number theorists and others because the coefficients carry interesting pieces of information about modular forms. In this thesis, we will study some statistical properties of the Fourier coefficients of modular forms and Maass forms.

### 2.1 Equidistribution and its Extensions

Definition 2.1.1. A sequence of real numbers $\left\{x_{n}\right\}_{n=1}^{\infty}$ is said to be uniformly distributed modulo 1 or equidistributed mod 1 with respect to the Lebesgue measure if for every pair of real numbers $a, b$ with $0 \leq a \leq b \leq 1$, we have

$$
\lim _{N \rightarrow \infty} \frac{\sharp\left\{n \leq N:\left(x_{n}\right) \in[a, b]\right\}}{N}=b-a,
$$

where $\left(x_{n}\right):=x_{n}-\left[x_{n}\right]$ denotes the fractional part of $x_{n}$.
Note that in the above expression, $b-a$ is the Lebesgue measure of the interval $[a, b]$. More generally equidistribution can be defined for any probability measure as follows:

Definition 2.1.2. A sequence of real numbers $\left\{x_{n}\right\}_{n=1}^{\infty}$ is said to be equidistributed (mod 1) with respect to a probability measure $\mu$ if for every pair of real numbers $a, b$ with $0 \leq a \leq b \leq 1$, we have

$$
\lim _{N \rightarrow \infty} \frac{\sharp\left\{n \leq N:\left(x_{n}\right) \in[a, b]\right\}}{N}=\mu([a, b]),
$$

where $\left(x_{n}\right):=x_{n}-\left[x_{n}\right]$ denotes the fractional part of $x_{n}$.
The above definition can be generalized to any unit interval as follows:
Definition 2.1.3. Let $[\alpha, \beta] \subset \mathbb{R}$ be any interval of unit length. A sequence of real numbers $\left\{x_{n}\right\}_{n=1}^{\infty}$ is said to be equidistributed in $[\alpha, \beta]$ with respect to a probability measure $\mu$ if for every pair of real numbers $a, b$ with $\alpha \leq a \leq b \leq \beta$, we have

$$
\lim _{N \rightarrow \infty} \frac{\sharp\left\{n \leq N:\left(x_{n}\right) \in[a, b]\right\}}{N}=\mu([a, b]) .
$$

Remark 2.1.4. In this thesis, the sequences of interest to us are equidistributed in the interval $\left[-\frac{1}{2}, \frac{1}{2}\right]$.

The above definition can be rewritten as follows:
A sequence $\left\{x_{n}\right\}$ of real numbers is said to be uniformly mod 1 if and only if for every interval $I \subset[0,1]$,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{n=1}^{N} \chi_{I}\left(x_{n}\right)=\int_{0}^{1} \chi_{I}(x) d \mu(x)
$$

Equivalently, for all (complex valued) Riemann integrable functions $f(x)$ of period 1 ,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{n=1}^{N} f\left(x_{n}\right)=\int_{0}^{1} f(x) d \mu(x)
$$

In particular, whenever $d \mu(x)=d x$, for all non-zero $m \in \mathbb{Z}$,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{n=1}^{N} e^{2 \pi i m x_{n}}=0
$$

In 1916, Weyl (see [51]) proved that the above condition is sufficient criterion for uniform distribution. More explicitly, he proves the following result now known as Weyl's Criterion:

Theorem 2.1.5. A sequence $\left\{x_{n}\right\}$ is uniformly distributed mod 1 if and only if

$$
c_{m}:=\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{n=1}^{N} e\left(m x_{n}\right)=0
$$

for every $m \in \mathbb{Z}, m \neq 0, e(t)=e^{2 \pi i t}$.
The necessary part is clear from the previous discussion. To prove the sufficient part we use Weierstrass approximation theorem that says every continious function can be approximated by trigonometric polynomials. For detail of the proof, see [27], page 172.
As an application of Weyl's theorem, one can derive the following:
Theorem 2.1.6. If $\theta$ is irrational then $\{(n \theta)\}$ is uniformly distributed mod 1 .
For a proof see [27], exercise 11.1.8. The definition of equidistribution of a sequence can be extended to the notion of equidistribution of sets as follows

Definition 2.1.7. A sequence of finite multisets $A_{n}$ with $\sharp A_{n} \rightarrow \infty$ counted with multiplicity is said to be equidistributed mod 1 with respect to a probability measure $\mu$ if for every continuous function $f$ on $[0,1]$, we have

$$
\lim _{n \rightarrow \infty} \frac{1}{\sharp A_{n}} \sum_{t \in A_{n}} f(t)=\int_{0}^{1} f(x) d \mu
$$

In this case the Weyl limits can be defined as follows:
Definition 2.1.8. For $m \in \mathbb{Z}$,

$$
c_{m}:=\lim _{n \rightarrow \infty} \frac{1}{\sharp A_{n}} \sum_{t \in A_{n}} e(m t) .
$$

Remark 2.1.9. When $\mu$ is the Lebesgue measure and $A_{n}=\left\{x_{1}, x_{2}, \ldots, x_{n}\right\}$, the above definition is same as the classical definition of uniform distribution mod 1 as in Definition 2.1.1

In 1928, Schoenberg [37] and in 1924 Wiener [52] independently gave a necessary and sufficient criterion for equidistribution of a sequence in $[0,1]$ with respect to some positive continuous measure in terms of the Weyl limit (see [22, Theorem 7.5] or [27, 11.3.3]). In the context of set equidistribution we also have the Wiener-Schoenberg theorem (see [27, Theorem 11.6.16]) as follows. We follow the presentation of page 195 of [27].

Theorem 2.1.10. Let $\left\{A_{n}\right\}$ be a sequence of sets as mentioned above. The sequence $\left\{A_{n}\right\}$ is equidistributed with respect to some positive continuous measure if and only if for all $m \in \mathbb{Z}$ the Weyl limits $c_{m}$ defined as above exist and

$$
\sum_{m=1}^{N}\left|c_{m}\right|^{2}=\mathrm{o}(N) \text { as } N \rightarrow \infty
$$

Proof. Consider the sequence of finite multisets $\left\{A_{n}\right\}$ such that $\# A_{n} \rightarrow \infty$ as $n \rightarrow \infty$ is equidistributed in $[0,1]$ with respect to some positive continious measure $\mu$. Then by definition, we know that for any continuous function $f:[0,1] \rightarrow \mathbb{C}$,

$$
\lim _{n \rightarrow \infty} \frac{1}{\sharp A_{n}} \sum_{t \in A_{n}} f(t)=\int_{0}^{1} f(x) d \mu \text {. }
$$

In particular, choosing $f=e(m x)$, we have

$$
c_{m}=\lim _{n \rightarrow \infty} \frac{1}{\sharp A_{n}} \sum_{t \in A_{n}} e(m t)=\int_{0}^{1} e(m x) d \mu .
$$

Hence, the Weyl limit $c_{m}$ exist for all $m \in \mathbb{Z}$. Moreover,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{m=1}^{N}\left|c_{m}\right|^{2}=\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{m=1}^{N} \int_{0}^{1} \int_{0}^{1} e(m(x-y)) d \mu(x) d \mu(y)
$$

Interchanging the sum and integral, the above is equal to

$$
\lim _{N \rightarrow \infty} \int_{0}^{1} \int_{0}^{1}\left(\frac{1}{N} \sum_{m=1}^{N} e(m(x-y)) d \mu(x) d \mu(y)\right)
$$

Note that,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{m=1}^{N} e(m(x-y))= \begin{cases}0 & \text { if } x-y \notin \mathbb{Z} \\ 1 & \text { otherwise }\end{cases}
$$

Since $\left|\frac{1}{N} \sum_{m=1}^{N} e(m(x-y))\right|$ is bounded by 1, by dominated convergence theorem, we can interchange the limit and integral. After interchanging the limit and integral the required limit is zero, because the measure of the set $\left\{(x, y) \in[0,1]^{2}: x-y \in \mathbb{Z}\right\}$ is zero.
Conversely, suppose that

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{m=1}^{N}\left|c_{m}\right|^{2}=0
$$

then by the Riesz representation theorem there exist a measurable function $g(x)$ such that

$$
c_{m}=\int_{0}^{1} e(m x) d g(x)
$$

Hence,

$$
\begin{aligned}
0 & =\lim _{N \rightarrow \infty} \frac{1}{N} \sum_{m=1}^{N}\left|c_{m}\right|^{2} \\
& =\int_{0}^{1} \int_{0}^{1} \lim _{N \rightarrow \infty} \frac{1}{N} e(m x) e(-m y) d g(x) d g(y) \\
& =\int_{0}^{1} \int_{0}^{1} \lim _{N \rightarrow \infty} \frac{1}{N} e(m(x-y)) d g(x) d g(y)
\end{aligned}
$$

that is

$$
\int_{0}^{1} \int_{0}^{1} f(x-y) d g(x) d g(y)=0
$$

where

$$
f(x-y)= \begin{cases}0 & x-y \notin \mathbb{Z} \\ 1 & \text { otherwise }\end{cases}
$$

Now we will show that $g$ is continuous. Let us assume that $g$ is not continuous. In particular, let $g$ has a jump discontinuity at $a \in(0,1)$. Then,

$$
\int_{0}^{1} \int_{0}^{1} f(x-y) d g(x) d g(y) \geq\left[\mu\left(a^{+}\right)-\mu\left(a^{-}\right)\right]^{2}>0
$$

which is a contradiction to the fact that the double integral is zero. Hence, $g$ is continuous. Now choosing $\mu(x)=\sum_{m=-\infty}^{\infty} c_{m} e(m x)$, we have for any $[\alpha, \beta]$ contained in $[0,1]$,

$$
\lim _{N \rightarrow \infty} \frac{1}{N} \#\left\{n \leq N: x_{n} \in[\alpha, \beta]\right\}=\int_{\alpha}^{\beta} \mu(x) d x
$$

Hence, the sequence of multisets $\left\{A_{n}\right\}$ is equidistributed with respect to the positive continuous measure $\mu$.

### 2.2 Fourier Analysis

According to our need, let us recall some facts from Fourier analysis in this section. The reader may refer to [35] for detailed study. For our convenience, let us define:

$$
e(x):=e^{2 \pi i x}
$$

Let $f$ be a periodic and integrable function of period 1 on $\mathbb{R}$. The Fourier series of $f$ is given by

$$
f(x)=\sum_{n=-\infty}^{\infty} \hat{f}(n) e(n x)
$$

where $\hat{f}(n)$ are called the Fourier coefficients, defined as

$$
\hat{f}(n):=\int_{0}^{1} f(x) e(-n x), n \in \mathbb{Z}
$$

Let $f_{i}, 1 \leq i \leq r$ be $r$ integrable functions on $\mathbb{R}$ of period 1 . Define the convolution of $r$ periodic integrable functions of period 1 on $\mathbb{R}$, denoted as

$$
f_{1} * f_{2} * \cdots * f_{r}: \mathbb{R} \rightarrow \mathbb{C}
$$

as follows:

$$
\begin{gather*}
f_{1} * \cdots * f_{r}(y)  \tag{2.1}\\
:=\int_{0}^{1} \cdots \int_{0}^{1} f_{1}\left(y_{1}\right) f_{2}\left(y_{2}\right) \ldots f_{r}\left(y-y_{1}-y_{2}-\cdots-y_{r-1}\right) d y_{r-1} d y_{r-2} \ldots d y_{1}
\end{gather*}
$$

Among the many interesting properties that convolution of periodic integrable functions satisfies, the following property serves our purpose:

$$
\begin{equation*}
\left(f_{1} * f_{2} * \cdots * f_{r}\right)^{\wedge}(n)=\hat{f}_{1}(n) \hat{f}_{2}(n) \ldots \hat{f}_{r}(n) \tag{2.2}
\end{equation*}
$$

In particular, for $r=2$,

$$
\left(f_{1} * f_{2}\right)^{\wedge}(n)=\hat{f}_{1}(n) \hat{f}_{2}(n)
$$

As the above property is important for us, we will give a proof of the above fact for $r=2$.

$$
\begin{aligned}
\left(f_{1} * f_{2}\right)^{\wedge}(n) & =\int_{0}^{1} \int_{0}^{1} f_{1}\left(y_{1}\right) f_{2}\left(y-y_{1}\right) e^{-2 \pi i n y} d y_{1} d y \\
& =\int_{0}^{1} \int_{0}^{1} f_{1}\left(y_{1}\right) e^{-2 \pi i n y_{1}} f_{2}\left(y-y_{1}\right) e^{-2 \pi i n\left(y-y_{1}\right)} d y_{1} d y
\end{aligned}
$$

Using Fubini's theorem, the above equals

$$
\int_{0}^{1} f_{1}\left(y_{1}\right) e^{-2 \pi i n y_{1}} d y_{1} \int_{0}^{1} f_{2}\left(y-y_{1}\right) e^{-2 \pi i n\left(y-y_{1}\right)} d y
$$

$$
=\hat{f}(n) \hat{g}(n)
$$

For $r>2$, we can prove similarly.
The following theorem can be concluded from the famous Riesz-Fischer theorem (see [35], page 91):

Theorem 2.2.1. [Riesz-Fischer]
Let $\left\{a_{n}\right\}$ be a sequence of real numbers. If

$$
\sum_{n=-\infty}^{\infty}\left|a_{n}\right|^{2}<\infty
$$

then there exists a unique periodic square Lebesgue integrable function $f$, that is $f \in L^{2}[0,1]$ such that

$$
f(x)=\sum_{n=-\infty}^{\infty} a_{n} e(n x), \text { that is } \hat{f}(n)=a_{n}
$$

### 2.3 Beurling-Selberg Polynomials

Let $\chi_{I}(x)$ be the characteristic function of the interval $I=[a, b]$ contained in $[0,1]$. For a positive integer $M$, define $\triangle_{M}(x)$ to be the Fejer's Kernel, defined as below:

$$
\Delta_{M}(x)=\sum_{|n|<M}\left(1-\frac{|n|}{M}\right) e(n x)=\frac{1}{M}\left(\frac{\sin \pi M x}{\sin \pi x}\right)^{2}
$$

The Mth order Beurling polynomials are defined as follows:

$$
\begin{gathered}
B^{*}{ }_{M}(x)=\frac{1}{M+1} \sum_{n=1}^{M}\left(\frac{n}{M+1}-\frac{1}{2}\right) \Delta_{M}\left(x-\frac{n}{M+1}\right) \\
+\frac{1}{2 \pi(M+1)} \sin (2 \pi(M+1) x)-\frac{1}{2 \pi} \Delta_{M+1}(x) \sin 2 \pi x+\frac{1}{2(M+1)} \Delta_{M+1}(x) .
\end{gathered}
$$

For an interval $[a, b]$, the Mth order Selberg polynomials are defined as below:

$$
S_{M}^{+}(x)=b-a+B_{M}^{*}(x-b)+B_{M}^{*}(a-x)
$$

and

$$
S_{M}^{-}(x)=b-a+B_{M}^{*}(b-x)+B_{M}^{*}(x-a)
$$

It is clear that both the above polynomials are trigonometric polynomial of degree at most $M$.

We will investigate some properties of Beurling-Selberg polynomials. The following result was observed by Beurling. It was Vaaler [47] who presented Beurling's result in the following language. We follow the presentation of [27, Theorem 11.4.3].

Theorem 2.3.1. Let

$$
\operatorname{sgn}(z)= \begin{cases}1 & \text { if } \operatorname{Re}(z) \geq 0 \\ -1 & \text { otherwise }\end{cases}
$$

If

$$
B(z)=\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}+\frac{2}{z}\right)
$$

then,

1. $B(z)$ is entire.
2. $B(x) \geq \operatorname{sgn}(x)$, for all real $x$.
3. $B(z)=\operatorname{sgn}(z)+\mathrm{O}\left(\frac{e^{2 \pi|I m z|}}{|z|}\right)$.
4. $\int_{-\infty}^{\infty}(B(x)-\operatorname{sgn}(x)) d x=1$.

Proof: To prove 1., it is enough to notice that $(\sin \pi z)^{2}$ is entire with zeros of order 2 at all integers and

$$
\left(\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}+\frac{2}{z}\right)
$$

is holomorphic in $\mathbb{C}-\{\mathbb{Z}-\{0\}\}$. It has poles of order two at all non-zero integers. All the poles of

$$
\left(\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}+\frac{2}{z}\right)
$$

are getting cancelled by the zeros of $\sin \pi z$. Hence, $B(z)$ is entire.
Now we will prove 2.. To prove the above claim we first observe the following:

$$
\left(\frac{\sin \pi z}{\pi}\right)^{2} \sum_{n=-\infty}^{n=\infty} \frac{1}{(z-n)^{2}}=1, z \notin \mathbb{Z}
$$

For $x>0$ we also have

$$
\begin{gathered}
\sum_{n=1}^{\infty} \frac{1}{(x+n)^{2}} \leq \sum_{n=1}^{\infty} \int_{x+n-1}^{x+n+1} \frac{d u}{u^{2}} \\
=\int_{x}^{\infty} \frac{d u}{u^{2}}=\frac{1}{x}=\sum_{n=0}^{\infty} \int_{x+n}^{x+n+1} \frac{d u}{u^{2}} \\
\leq \sum_{n=0}^{\infty} \frac{1}{(x+n)^{2}}
\end{gathered}
$$

Now for Re $z \geq 0$,

$$
\begin{gathered}
B(z)-\operatorname{sgn}(z)=\left(\frac{\sin \pi z}{\pi}\right)^{2} \sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}} \\
-\left(\frac{\sin \pi z}{\pi}\right)^{2} \sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}+\left(\frac{\sin \pi z}{\pi}\right)^{2} \frac{2}{z}-\operatorname{sgn} z
\end{gathered}
$$

$$
\begin{aligned}
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}+\frac{2}{z}\right)-\operatorname{sgn}(z) \\
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\frac{2}{z}+\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-\sum_{n=-\infty}^{-1} \frac{1}{(z-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}\right)-1 \\
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\frac{2}{z}+\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-2 \sum_{n=-\infty}^{-1} \frac{1}{(z-n)^{2}}\right)-1 \\
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\frac{2}{z}-2 \sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}\right)
\end{aligned}
$$

Now when $z$ is real (say $x$ ), we have

$$
\begin{aligned}
B(x)-\operatorname{sgn}(x) & \geq\left(\frac{\sin \pi x^{2}}{\pi}\right)\left(\frac{2}{x}-\frac{2}{x}\right) \\
& \geq 0
\end{aligned}
$$

Similarly, for $R e z<0$,

$$
\begin{aligned}
& \quad B(z)-\operatorname{sgn}(z) \\
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\frac{2}{z}+2 \sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}\right) \\
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\frac{2}{z}+2\left\{\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}+\sum_{n=-\infty}^{-1} \frac{1}{(z-n)^{2}}-\sum_{n=-\infty}^{-1} \frac{1}{(z-n)^{2}}\right\}\right) \\
& =\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\frac{2}{z}-2 \sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}\right)+2 .
\end{aligned}
$$

Now if $z$ is real (say $x$ ), we have

$$
\begin{aligned}
B(x)-\operatorname{sgn}(x) & \geq\left(\frac{\sin \pi x}{\pi}\right)^{2}\left(\frac{2}{x}-\frac{2}{x}\right)+2 \\
& \geq 0
\end{aligned}
$$

Hence,

$$
B(x)-\operatorname{sgn}(x) \geq 0 \text { for real } x .
$$

To prove 3., observe that

$$
\sin ^{2} \pi z=\mathrm{O}\left(e^{2 \pi|\operatorname{Im}(z)|}\right)
$$

In addition, for $x, y>0$, we have

$$
\begin{aligned}
\sum_{n=0}^{\infty} \frac{1}{(x+n)^{2}+y^{2}} & \leq \frac{1}{x^{2}+y^{2}}+\min \left(\int_{0}^{\infty} \frac{d t}{(x+t)^{2}}, \frac{d t}{\left(t^{2}+y^{2}\right)^{2}}\right) \\
& =\frac{1}{x^{2}+y^{2}}+\min \left(\frac{1}{x}, \frac{\pi}{2 y}\right) .
\end{aligned}
$$

Hence

$$
\sum_{n=1}^{\infty} \frac{1}{|z+n|^{2}}=\mathrm{O}\left(\frac{1}{|z|}\right), \text { for } \operatorname{Re}(z) \geq 0
$$

and

$$
\sum_{n=0}^{\infty} \frac{1}{|z-n|^{2}}=\mathrm{O}\left(\frac{1}{|z|}\right), \text { for } \operatorname{Re}(z)<0
$$

Combining all above facts, we get 3 ..
We now prove 4.
Proof. We know,

$$
\int_{-\infty}^{\infty}(B(x)-\operatorname{sgn}(x)) d x=\lim _{A \rightarrow \infty} \int_{-A}^{A}(B(x)-\operatorname{sgn}(x)) d x
$$

Now,

$$
\begin{aligned}
\int_{-A}^{A}(B(x)-\operatorname{sgn}(x)) d x & =\int_{-A}^{0}(B(x)-\operatorname{sgn}(x)) d x+\int_{0}^{A}(B(x)-\operatorname{sgn}(x)) d x \\
& =\int_{0}^{A}(B(x)+B(-x)) d x
\end{aligned}
$$

$$
\begin{aligned}
= & \int_{0}^{A}\left(\frac{\sin \pi x}{\pi}\right)^{2}\left\{\left(\sum_{n=0}^{\infty} \frac{1}{(x-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(x+n)^{2}}+\frac{2}{x}\right)\right. \\
& \left.+\left(\sum_{n=0}^{\infty} \frac{1}{(-x-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(-x+n)^{2}}-\frac{2}{x}\right)\right\} d x \\
= & \int_{0}^{A}\left(\frac{\sin \pi x}{\pi}\right)^{2}\left\{\frac{1}{x^{2}}+\sum_{n=1}^{\infty} \frac{1}{(x-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(x+n)^{2}}\right. \\
& \left.+\frac{1}{x^{2}}+\sum_{n=1}^{\infty} \frac{1}{(x+n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(x-n)^{2}}\right\} d x \\
= & \int_{0}^{A}\left(\frac{\sin \pi x}{\pi}\right)^{2}\left(\frac{2}{x^{2}}\right) d x \\
= & 2 \int_{0}^{A}\left(\frac{\sin \pi x}{\pi}\right)^{2} d x
\end{aligned}
$$

Hence,

$$
\lim _{A \rightarrow \infty} \int_{0}^{A}(B(x)-\operatorname{sgn}(x)) d x=2 \int_{0}^{A}\left(\frac{\sin \pi x}{\pi}\right)^{2} d x=1
$$

Selberg (see [38, page 213-218]) used Beurling polynomials to approximate the characteristic function of an interval. More precisely,

Theorem 2.3.2. Selberg,(1970) Let $I=[a, b]$ be an interval and $\chi_{I}$ be its characteristic function. Then there are continuous functions $S^{+}(x)$ and $S^{-}(x)$ in $L^{1}(\mathbb{R})$ such that

$$
S^{-}(x) \leq \chi_{I}(x) \leq S^{+}(x)
$$

with $\hat{S}^{ \pm}(t)=0$, for $|t| \geq 1$. In addition,

$$
\int_{-\infty}^{\infty}\left(\chi_{I}(x)-S^{-}(x)\right) d x=1
$$

and

$$
\int_{-\infty}^{\infty}\left(S^{+}(x)-\chi_{I}(x)\right) d x=1
$$

Proof: We know that

$$
B(z)=\left(\frac{\sin \pi z}{\pi}\right)^{2}\left(\sum_{n=0}^{\infty} \frac{1}{(z-n)^{2}}-\sum_{n=1}^{\infty} \frac{1}{(z+n)^{2}}+\frac{2}{z}\right) .
$$

Let

$$
S^{+}(x)=\frac{1}{2}(B(x-a)+B(b-x))
$$

Thus,

$$
S^{+}(x) \geq \frac{1}{2}(\operatorname{sgn}(x-a)+\operatorname{sgn}(b-x))=\chi_{I}(x)
$$

as we have

$$
\begin{gathered}
B(x-a) \geq \operatorname{sgn}(x-a) \text { and } B(b-x) \geq \operatorname{sgn}(b-x) . \\
\frac{1}{2}(B(x-a)+B(x-b)) \geq \frac{1}{2}(\operatorname{sgn}(x-a)+\operatorname{sgn}(b-x))
\end{gathered}
$$

and so

$$
S^{+}(x) \geq \frac{1}{2}(\operatorname{sgn}(x-a)+\operatorname{sgn}(b-x)) .
$$

Hence,

$$
\int_{-\infty}^{\infty}\left(S^{+}(x)-\chi_{I}(x)\right) d x=1
$$

as

$$
\int_{-\infty}^{\infty}\left(S^{+}(x)-\chi_{I}(x)\right) d x=\frac{1}{2} \int_{-\infty}^{\infty}\left[(B(x-a)+B(x-b))-\chi_{I}(x)\right] d x
$$

But again,

$$
\begin{equation*}
\chi_{I}(x)=\frac{1}{2}\{\operatorname{sgn}(x-a)+\operatorname{sgn}(b-x)\} . \tag{2.3}
\end{equation*}
$$

Using (2.3), we have

$$
\begin{aligned}
\int_{-\infty}^{\infty}\left(S^{+}(x)-\chi_{I}(x)\right) d x & =\frac{1}{2} \int_{-\infty}^{\infty}[(B(x-a)+B(x-b))-\{\operatorname{sgn}(x-a)+\operatorname{sgn}(b-x)\}] d x \\
& =\frac{1}{2}+\frac{1}{2} \\
& =1
\end{aligned}
$$

We get the above equality by using Theorem 2.3.1. And from above it is also clear that $S^{+}(x) \in L^{1}(R)$. Since $S^{+}(x)$ is a restriction of an entire function, it is continuous.
Now we will show that for $t>1$,

$$
\begin{equation*}
\hat{S}^{+}(t)=\int_{-\infty}^{\infty} S^{+}(x) e(-t x) d x=0 \tag{2.4}
\end{equation*}
$$

Let $J(A, B)=\int_{-A}^{B} S^{+}(x) e(-t x) d x$.
To prove (2.4), we will begin by proving the following:

$$
J(A, B)=\int_{-A}^{B} S^{+}(x) e(-t x) d x=\mathrm{O}\left(\frac{1}{A}+\frac{1}{B}\right) \text { as } A, B \rightarrow \infty
$$

By contour integration,

$$
\begin{align*}
& \int_{-A}^{B} S^{+}(x) e(-t x) d x=\int_{-A}^{-A-i T} S^{+}(x) e(-t x) d x \\
+ & \int_{-A-i T}^{B-i T} S^{+}(x) e(-t x) d x+\int_{B-i T}^{B} S^{+}(x) e(-t x) d x \tag{2.5}
\end{align*}
$$

First we will estimate the second integral on the right-hand side of (2.5).

$$
\begin{align*}
\int_{-A-i T}^{B-i T} S^{+}(x) e(-t x) d x & =\frac{1}{2} \int_{-A-i T}^{B-i T}(B(x-a)+B(b-x)) e(-t x) d x \\
& =\frac{1}{2} \int_{-A-i T}^{B-i T} B(x-a) e(-t x) d x+\frac{1}{2} \int_{-A-i T}^{B-i T} B(b-x) e(-t x) d x . \tag{2.6}
\end{align*}
$$

Substituting $x-i T+a$ and $-x+i T+b$ in place of $x$ in the first and second integral respectively of (2.6), we have

$$
\begin{aligned}
\int_{-A-i T}^{B-i T} S^{+}(x) e(-t x) d x & \ll \int_{A^{\prime}}^{B^{\prime}}|B(x-i T)| e^{-2 \pi t T} d x \\
& \ll \int_{A^{\prime}}^{B^{\prime}} e^{2 \pi T} e^{-2 \pi t T} d x
\end{aligned}
$$

where $A^{\prime}=A+\max \{|a|,|b|\}$ and $B^{\prime}=B+\max \{|a|,|b|\}$.

The last inequality is by using the previous theorem, that is

$$
B(z)=\operatorname{sgn}(z)+\mathrm{O}\left(\frac{e^{2 \pi|\operatorname{Im}(z)|}}{|z|}\right) .
$$

Hence,

$$
\int_{-A-i T}^{B-i T} S^{+}(x) e(-t x) d x \ll \int_{A^{\prime}}^{B^{\prime}} e^{2 \pi T} e^{-2 \pi t T} d x
$$

Now, since $t>1$, the above integral goes to 0 , whenever $T \rightarrow \infty$.
Now we estimate the first and third integrals.
For $z=-A+i y$ and $A>|a|$,
$B(z-a)=B(-A+i y-a)=B(-A-a+i y)=\operatorname{sgn}(-A-a+i y)+\mathrm{O}\left(\frac{e^{-2 \pi y}}{A}\right)$.
For $A>|b|$,

$$
B(b-z)=B(b+A-i y)=\operatorname{sgn}(b+A-i y)+\mathrm{O}\left(\frac{e^{-2 \pi y}}{A}\right)
$$

So,

$$
S^{+}(z) \ll \frac{e^{-2 \pi y}}{A}
$$

The first integral on the right hand side of (2.5) is

$$
\ll \frac{1}{A} \int_{-\infty}^{0} e^{-2 \pi y} e^{2 \pi t y} d y \ll \frac{1}{A} .
$$

Similarly we can estimate the last integral on the right-hand side of (2.5).

$$
J(A, B)=\mathrm{O}\left(\frac{1}{A}+\frac{1}{B}\right) .
$$

Now letting $A, B \rightarrow \infty$, we have shown that $\hat{S}^{+}(t)=0$, for $t>1$.
For $t<1$, we can use the fact $\hat{S}^{+}(-t)=\overline{\hat{S}}^{+}(t)$ and deduce the desired result.

Using continuity of $S^{+}$, the proof for $t= \pm 1$ follows. Now defining

$$
S^{-}(x)=-\frac{1}{2}(B(x-a)+B(x-b))
$$

and proceeding analogously we get the required result.
Note that, for any $\delta>0$, choosing $S^{ \pm}(x)$ for the interval [ $\left.\delta a, \delta b\right]$, we have

$$
\begin{gathered}
\hat{S}^{ \pm}(t)=0 \text { for }|t|>\delta \text { and } \\
\int_{\infty}^{\infty}\left(\chi_{I}(x)-S^{-}(x)\right)=\int_{\infty}^{\infty}\left(S^{+}(x)-\chi_{I}(x)\right)=\frac{1}{\delta}
\end{gathered}
$$

Since $S^{ \pm}$is continuous, the crux of the above discussion is that the characteristic function of any interval in $\mathbb{R}$ can be approximated by a continuous function. For detail see Exercise 11.4.5 of [27]. Moreover, from the work of Vaaler (see [27]), we have the following facts:
For all $M \geq 1$
(a) For a subinterval $[a, b]$ of $[0,1]$,
$S_{M}^{-}(x) \leq \chi_{I}(x) \leq S_{M}^{+}(x)$.

$$
S_{M}^{ \pm}(x)=\sum_{0 \leq|m| \leq M} \hat{S}_{m}^{ \pm}(m) e(m x)
$$

(b)

$$
\hat{S}_{M}^{ \pm}(0)=b-a \pm \frac{1}{M+1}
$$

When $M \neq 0$,
$\left|\hat{S}_{M}^{ \pm}(m)\right| \leq \frac{1}{M+1}+\min \left\{b-a, \frac{1}{\pi|m|}\right\}$.
(c)

$$
\left\|S_{M}^{ \pm}(x)-\chi(x)\right\|_{L_{1}} \leq \frac{1}{M+1}
$$

Using above, we can conclude that

$$
\left|\hat{S}_{M}^{ \pm}(n)-\hat{\chi}_{I}(n)\right| \leq\left\|S_{M}^{+}-\chi(x)\right\|_{L_{1}} \leq \frac{1}{M+1}
$$

(d)

For $n \neq 0$, note that

$$
\left|\hat{\chi}_{I}(n)\right|=\left|\frac{\sin \pi n(b-a)}{\pi k}\right| \leq \min \left(b-a, \frac{1}{\pi|n|}\right)
$$

For $0<|n|<M$,

$$
\left|\hat{S}_{M}^{ \pm}(n)\right| \leq \frac{1}{M+1}+\min \left(b-a, \frac{1}{\pi|n|}\right)
$$

The above result follows from following facts combining with (c):

### 2.4 Erdös-Turán inequality

In 1948 Erdös and Turán gave an effective version of the Weyl's criterion. The following version of the Erdös-Turán inequality can be found in [27]. Let $\left\{x_{n}\right\}_{n=1}^{\infty}$ be a sequence. For any natural number $N$ and an interval $(a, b)$ in $\mathbb{R}$, define the discrepancy $D_{N}$ of the sequence $\left\{x_{n}\right\}_{n=1}^{\infty}$ as

$$
D_{N}:=\sup _{0 \leq a<b \leq 1}\left|\frac{\#\left\{n \leq N: a \leq x_{n} \leq b\right\}}{N}-(b-a)\right| .
$$

Theorem 2.4.1. (Erdös-Turán, 1948) For any integer $M \geq 1$,

$$
D_{N} \leq \frac{1}{M+1}+3 \sum_{m=1}^{M} \frac{1}{N m}\left|\sum_{n=1}^{N} e^{2 \pi i m x_{n}}\right|
$$

There are several proofs of the above inequality. In [25] Montgomery gives a proof using Beurling-Selberg polynomials (see [27, Theorem 11.4.8]). Generalizing the idea implicit in Montgomery's work, Murty and Sinha have proved the following variant of Erdös-Turán inequality in ( [32, Theorem 8]).

Theorem 2.4.2. [Murty-Sinha] Let $I=[a, b] \subset[0,1]$ and $c_{m}$ 's are the Weyl limits given by Definition 2.1.8. Let

$$
\mu=F(x) d x
$$

where

$$
F(x)=\sum_{m=-\infty}^{\infty} c_{m} e(m x)
$$

Define

$$
N_{I}(V):=\#\left\{n \leq V: x_{n} \in I\right\}
$$

and

$$
D_{I, V}(\mu):=\left|N_{I}(V)-V \mu(I)\right| .
$$

Then,
$D_{I, V}(\mu) \leq \frac{V| | \mu| |}{M+1}+\sum_{1 \leq|m| \leq M}\left(\frac{1}{M+1}+\min \left(b-a, \frac{1}{\pi|m|}\right)\right)\left|\sum_{n=1}^{r} e\left(m x_{n}\right)-V c_{m}\right|$ where $\|\mu\|=\underset{x \in[0,1]}{\text { Sup }}|F(x)|$.

We will follow the presentation of [32].
Proof. Let $\chi_{I}$ be the characteristic function of interval $I$. Then by Theorem 2.3.2

$$
\begin{equation*}
S_{M}^{-}\left(x_{n}\right) \leq \chi_{I}\left(x_{n}\right) \leq S_{M}^{+}\left(x_{n}\right) \tag{2.7}
\end{equation*}
$$

Taking sum over all terms of (2.7), we get

$$
\sum_{n \leq V} S_{M}^{-}\left(x_{n}\right) \leq \sum_{n \leq V} \chi_{I}\left(x_{n}\right) \leq \sum_{n \leq V} S_{M}^{+}\left(x_{n}\right) .
$$

Since $S_{M}^{ \pm}(x)$ is a periodic and entire function, hence by Fourier's theory, we have

$$
\sum_{n \leq V} S_{M}^{ \pm}\left(x_{n}\right)=\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) \sum_{n \leq V} e\left(m x_{n}\right)
$$

Subtracting $\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) V c_{m}$ from both the sides, we get

$$
\sum_{n \leq V} S_{M}^{ \pm}\left(x_{n}\right)-V \sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) c_{m}=\sum_{1 \leq|m| \leq M} \hat{S}_{M}^{ \pm}(m) \sum_{n \leq V}\left(e\left(m x_{n}\right)-V c_{m}\right)
$$

Taking the absolute value on both sides, we get

$$
\left|\sum_{n \leq V} S_{M}^{ \pm}\left(x_{n}\right)-V \sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) c_{m}\right|=\left|\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) \sum_{n \leq V}\left(e\left(m x_{n}\right)-V c_{m}\right)\right| .
$$

So,

$$
\left|\sum_{n \leq V} S_{M}^{ \pm}\left(x_{n}\right)-V \sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) c_{m}\right| \leq\left|\sum_{1 \leq|m| \leq M \leq V} \hat{S}_{M}^{ \pm}(m) \sum_{n \leq V}\left(e\left(m x_{n}\right)-V c_{m}\right)\right| .
$$

Since $\hat{S}_{M}^{ \pm}(m) c_{m}=0$ for $|m|>M$, without loss of generality we can extend the sum $\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) c_{m}$ to $\sum_{m \in \mathbb{Z}} \hat{S}_{M}^{ \pm}(m)$. Then,

$$
\sum_{m} \hat{S}_{M}^{ \pm}(m)=\sum_{m} \int_{0}^{1} S_{m}^{ \pm}(x) e(-m x) d x
$$

Now interchanging the sum and integral, we have

$$
\sum_{m} \hat{S}_{M}^{ \pm}(m)=\int_{0}^{1} S_{M}^{ \pm}(x) d \mu
$$

Now from (c), we have

$$
\left|\int_{0}^{1}\left(S_{M}^{ \pm}(x)-\chi_{I}(x)\right) d \mu\right| \leq \frac{\|\mu\|}{M+1}
$$

and

$$
\hat{S}_{M}^{ \pm}(m) \leq \frac{1}{M+1}+\min \left(b-a, \frac{1}{\pi|n|}\right)
$$

In Section 2.5 and Section 2.6 we review definitions and properties of cusp forms and Maass forms respectively. We also review the definition and properties of Hecke operators acting on all these spaces respectively. Based on our need, we will state several results without giving proof, but for a concerned reader, we will give reference for details.

### 2.5 Modular Forms and Hecke operators

In this section, we will review basic definitions and facts about modular forms. For details, the reader may consult [28] and [20]. Let us start with the following notation:
Let $\mathcal{H}$ be the upper half plane, defined as

$$
\mathcal{H}:=\{x+i y: x, y \in \mathbb{R}, y>0\}
$$

For a commutative ring $R$ with unity, define

$$
G L_{2}^{+}(R):=\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right): a, b, c, d \in R \text { and } a d-b c>0\right\} .
$$

Let $\mathbb{Q}$ be the set of rational numbers. Define the action of $G L_{2}^{+}(\mathbb{Q})$ on $\mathcal{H}$ as follows:
For any $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in G L_{2}^{+}(\mathbb{Q})$ and $z \in \mathcal{H}$,

$$
\gamma z:=\frac{a z+b}{c z+d} .
$$

Now, for any positive integer $k$ and a meromorphic function $f$ defined on $\mathcal{H}$, define

$$
\left.f\right|_{k} \gamma(z)=(\operatorname{det} \gamma)^{\frac{k}{2}} f(\gamma z)(c z+d)^{-k} .
$$

Define

$$
S L_{2}(\mathbb{Z}):=\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right): a, b, c, d \in \mathbb{Z} \text { and } a d-b c=1\right\} .
$$

Let $\Gamma$ be a subgroup of finite index in $S L_{2}(\mathbb{Z})$. Let $f$ be a holomorphic function on $\mathcal{H}$ such that

$$
\left.f\right|_{k} \gamma(z)=f(z) \text { for all } \gamma \in \Gamma
$$

Since $\Gamma$ has finite index in $S L_{2}(\mathbb{Z})$, there exists a positive integer, say $l$, such that

$$
\left(\begin{array}{ll}
1 & 1 \\
0 & 1
\end{array}\right)^{l}=\left(\begin{array}{ll}
1 & l \\
0 & 1
\end{array}\right) \in \Gamma
$$

Hence,

$$
f(z+l)=f(z) \text { for all } z \in \mathcal{H}
$$

That is $f$ is of period $l$. So, $f$ has a Fourier series expansion at $\infty$,

$$
\begin{equation*}
f(z)=\sum_{n=-\infty}^{\infty} a_{n}(f) q^{\frac{n}{l}}, q=e^{2 \pi i z} . \tag{2.8}
\end{equation*}
$$

Definition 2.5.1. A function $f$ as defined above is said to be holomorphic at infinity if

$$
a_{n}=0 \text { for all } n<0 \text { in (2.8). }
$$

Definition 2.5.2. A function $f$ as defined above is said to vanish at infinity if $a_{0}=0$ in (2.8).

Since the above defined action is a transitive action in $\Gamma$, for any $\gamma^{\prime} \in$ $S L_{2}(\mathbb{Z}),\left(\gamma^{\prime}\right)^{-1} \Gamma \gamma^{\prime}$ has finite index in $S L_{2}(\mathbb{Z})$ and $\left.f\right|_{k} \gamma^{\prime}(z)$ is also periodic. Hence it has a Fourier series expansion at infinity,

$$
\left.f\right|_{k} \gamma^{\prime}(z)=\sum_{n=-\infty}^{\infty} a_{n}\left(\left.f\right|_{k} \gamma^{\prime}\right) q^{\frac{n}{l}}
$$

With the above notations and definitions, we will define modular forms.
Definition 2.5.3. Let $k \in \mathbb{Z}$. A function $f$ on $\mathcal{H}$ is said to be a modular form of weight $k$ with respect to $\Gamma$ if
(i) $f$ is holomorphic on $\mathcal{H}$.
(ii) $f(\gamma z)=(c z+d)^{k} f(z)$ for all $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in \Gamma$.
(iii) $f_{k} \mid \gamma(z)$ is holomorphic at $\infty$ for all $\gamma \in S L_{2}(\mathbb{Z})$.

Remark 2.5.4. Let $\Gamma$ be a congruent subgroup of $S L_{2}(\mathbb{Z})$ of level $N$. Let $q_{N}=$ $e^{\frac{2 \pi i z}{N}}$ for $z \in \mathcal{H}$. Let $f: \mathcal{H} \Rightarrow \mathbb{C}$ satisfies (i) and (ii) in Definition 2.5.3 and in the Fourier expansion $f(z)=\sum_{n=0}^{\infty} a_{n} q_{N}^{n}$, the coefficients satisfies

$$
\left|a_{n}\right| \leq c n^{r} \text { for some constant } c, r \in \mathbb{R}^{+}
$$

Then $f_{k} \mid \gamma(z)$ is holomorphic at $\infty$ for all $\gamma \in S L_{2}(\mathbb{Z})$. (see Proposition 1.2.4 of [15].)

Definition 2.5.5. We say that $f$ vanishes at the cusps if $\left.f\right|_{k} \gamma^{\prime}$ vanishes at cusps for all $\gamma^{\prime} \in S L_{2}(\mathbb{Z})$.

Definition 2.5.6. A function $f$ as defined above is said to be a cusp form of weight $k$ with respect to $\Gamma$ if
(i) $f$ is a modular form of weight $k$ with respect to $\Gamma$.
(ii) $f$ vanishes at the cusps.

Remark 2.5.7. To check whether $f$ vanishes at the cusps, it is sufficient to check that the constant term in the Fourier series expansion (2.8) of $f_{k} \mid \gamma$ is zero for all $\gamma \in S L_{2}(\mathbb{Z})$.

Now we will define one of the most important families of linear operators in the study of modular forms known as Hecke operators.
We will define these operators using the theory of double cosets. Let $\Gamma_{1}$ and $\Gamma_{2}$ be two congruent subgroup of $S L_{2}(\mathbb{Z})$. Let $\alpha$ be an element of $G L_{2}^{+}(\mathbb{Q})$. A double coset in $G L_{2}^{+}(\mathbb{Q})$ is defined as follows:

$$
\Gamma_{1} \alpha \Gamma_{2}:=\left\{g_{1} \alpha g_{2}: g_{1} \in \Gamma_{1}, g_{2} \in \Gamma_{2}\right\}
$$

We can decompose the double coset as a disjoint union of $\Gamma_{1}$-orbits, that is

$$
\Gamma_{1} \alpha \Gamma_{2}=\coprod_{j} \Gamma_{1} \alpha \gamma_{j} \text { for some } \gamma_{j} \in \Gamma_{2}
$$

Let $S^{\times}$be a multiplicative subgroup of $(\mathbb{Z} / N \mathbb{Z})^{\times}$. Let $S^{+}$be an additive subgroup of $\mathbb{Z}$.
For a positive integer $n$, let

$$
X_{n}=X_{n}\left(N, S^{\times}, S^{+}\right):=\left\{\left(\begin{array}{ll}
a & b \\
0 & d
\end{array}\right): a, b, c, d \in \mathbb{Z}, a d-b c=n, N \mid c, a \in S^{\times}, b \in S^{+}\right\}
$$

Note that $X_{1}$ is a congruent subgroup of $S L_{2}(\mathbb{Z})$ and the set $X_{n}$ is invariant under the left and right action of $X_{1}$. We also note that there is a finite set of orbit representatives $\alpha_{i} \in X_{n}$ such that

$$
X_{n}=\coprod_{i} X_{1} \alpha_{i}
$$

Definition 2.5.8. For any modular form $f$ of weight $k$ on $X_{1}$, the $n$-th Hecke operator $T_{n}$ that maps to a modular form of weight $k$ on $X_{1}$ is defined as

$$
T_{n}(f):=\left.n^{\frac{k}{2}-1} \sum_{i} f\right|_{k} \alpha_{i} .
$$

These operators take modular forms to modular forms and cusp forms to cusp forms. In this thesis, we are interested in a particular kind of subgroup of $S L_{2}(\mathbb{Z})$ :

$$
\Gamma_{0}(N):=\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right) \in S L_{2}(\mathbb{Z}): c \equiv 0(\bmod N)\right\}
$$

for some positive integer $N$. Note that $X_{1}\left(N,(\mathbb{Z} / N \mathbb{Z})^{\times}, \mathbb{Z}\right)=\Gamma_{0}(N)$ and it has index

$$
\psi(N):=N \prod_{p \mid N}\left(1+\frac{1}{p}\right)
$$

in $S L_{2}(\mathbb{Z})$, for example (see Exercise 2.3.4, page 19 in [28]). Since the subgroup $\Gamma_{0}(N)$ is a finite index subgroup of $S L_{2}(Z)$ and $\left(\begin{array}{ll}1 & 1 \\ 0 & 1\end{array}\right) \in \Gamma_{0}(N)$, modular forms with respect to $\Gamma_{0}(N)$ are periodic with period 1 and will have Fourier series expansion of the form

$$
f(z)=\sum_{n=0}^{\infty} a_{n}(f) q^{n}, \quad q=e^{2 \pi i z}
$$

The Fourier series expansion at $i \infty$ of a cusp form is

$$
f(z)=\sum_{n=1}^{\infty} a_{n}(f) q^{n}
$$

For our further use, let $M(N, k)$ be the space of modular forms of weight $k$ with respect to $\Gamma_{0}(N)$ and $S(N, k)$ be the space of cusp forms of weight $k$ with respect to $\Gamma_{0}(N)$. We know that $M(N, k)$ and $S(N, k)$ are finite dimensional vector spaces over $\mathbb{C}$.

Our particular interest is when the Hecke operators act on the spaces $S(N, k)$. The above mentioned Hecke operators have several interesting properties. Here we will mention some of them as the following theorem:

Theorem 2.5.9. (Hecke) For positive integers $n$, the Hecke operators $T_{n}$ acting on $S(N, k)$ satisfy the following properties:
(i) For $m, n \geq 1$,

$$
T_{m} T_{n}=T_{n} T_{m} .
$$

(ii) If $(m, n)=1$, then

$$
T_{m} T_{n}=T_{m n} .
$$

(iii) For a prime power $p^{r}, r \geq 1$,

$$
T_{p^{r}} T_{p}= \begin{cases}T_{p^{r+1}}+p^{k-1} T_{p^{r-1}} & \text { if }(p, N)=1 \\ \left(T_{p}\right)^{r} & \text { otherwise. }\end{cases}
$$

For a proof of Theorem 4.1.1, see [20], page 60.
Theorem 2.5.10. Let $n \geq 1$ be coprime to $N$. Then the Hecke operators $T_{n}$ are self-adjoint operators with respect to the Petersson inner product, where the Petersson inner product on $S(N, k)$ is defined as follows:
For any $f, g \in S(N, k)$,

$$
\begin{equation*}
\langle f, g\rangle=\int_{\mathcal{H} / \Gamma_{0}(N)} f(z) \bar{g}(z) y^{k} \frac{d x d y}{y^{2}} . \tag{2.9}
\end{equation*}
$$

For a proof, see [20] page 75. From Theorem 2.5.10, we can conclude that the eigenvalues of the Hecke operators $T_{n}$ are real, whenever $(n, N)=1$. Moreover, the collection of all Hecke operators acting on $S(N, k)$ generate an algebra over $\mathbb{C}$. For a positive integer $n \geq 1$, let

$$
\left\{\lambda_{n, i, N}, 1 \leq i \leq s(N, k)\right\}
$$

denote the set of eigenvalues of $T_{n}$. For any positive integer $n$, let

$$
T_{n}^{\prime}:=\frac{T_{n}}{n^{\frac{k-1}{2}}}
$$

be the normalized Hecke operator acting on $S(N, k)$ with eigenvalues

$$
\left\{a_{n, i, N}=\frac{\lambda_{n, i, N}}{n^{\frac{k-1}{2}}}, 1 \leq i \leq s(N, k)\right\} .
$$

By the celebrated theorem of Deligne [14], which proves the famous Ramanujan conjecture, we know that for any prime $p$ not dividing $N$, the eigenvalues of $T_{p}^{\prime}$ lie in the interval $[-2,2]$. Since the Hecke operators are diagonalizable and commutative there exists an ordered basis for $S(N, k)$ such that every operator $T_{n}$ is represented on that basis by a diagonal matrix. An element of such a basis is called an eigenform. A basis of $S(N, k)$ consisting of eigenforms is called a Hecke eigenbasis. An eigenform $f$ is said to be a normalized eigenform, if the first Fourier coefficient $a_{1}(f)=1$.

### 2.6 Maass Forms

In this section, we will follow the presentation of [24]. Let $\mathcal{H}$ be the upper half-plane. Consider $\Gamma=S L_{2}(\mathbb{Z})$. The non-Euclidean Laplace operator on the space of smooth functions on $\mathcal{H}$ is given by

$$
\triangle:=-y^{2}\left(\frac{\delta^{2}}{\delta x^{2}}+\frac{\delta^{2}}{\delta y^{2}}\right)
$$

If $f(z)=u+i v$, then

$$
\triangle(f(z))=v^{2}\left(\frac{\delta^{2}}{\delta x^{2}}+\frac{\delta}{\delta y^{2}}\right) f(z)
$$

The operator $\triangle$ is invariant under the action of $S L_{2}(\mathbb{Z})$ on $\mathcal{H}$, that is, for any smooth function $f$ on $\mathcal{H}$ and for any $z \in \mathcal{H}$,

$$
(\triangle f)(\gamma z)=\triangle(f(\gamma z))
$$

where $\gamma=\left(\begin{array}{ll}a & b \\ c & d\end{array}\right) \in S L_{2}(\mathbb{Z})$, and the action of $S L_{2}(\mathbb{Z})$ on $\mathcal{H}$ defined as follows:

$$
\gamma z=\frac{a z+b}{c z+d} .
$$

Definition 2.6.1. A smooth function $f \neq 0$ on $\mathcal{H}$ is called a Maass form with respect to the group $\Gamma$ if
(i) for all $\gamma \in \Gamma$ and $z \in \mathcal{H}$;

$$
f(\gamma z)=f(z)
$$

(ii) $f$ is an eigenfunction of $\triangle$, that is,

$$
\triangle f=\lambda f \text { for } \lambda \in \mathbb{C}
$$

(iii) there exists a positive integer $N$ such that

$$
f(z) \ll y^{N}, \text { as } y \rightarrow \infty
$$

Definition 2.6.2. A Maass form $f$ is said to be a cusp form if, for all $z \in \mathcal{H}$,

$$
\int_{0}^{1} f\left(\left(\begin{array}{ll}
1 & x \\
0 & 1
\end{array}\right) z\right) d x=\int_{0}^{1} f(z+x) d x=0
$$

We know that the Maass cusp forms span a subspace $\mathcal{C}(\Gamma \backslash \mathcal{H})$ in $L^{2}(\Gamma \backslash \mathcal{H})$, where $L^{2}(\Gamma \backslash \mathcal{H})$ denotes the square integrable functions on $\Gamma \backslash \mathcal{H}$ and the $L^{2}$ norm is induced by the Petersson inner product.

### 2.7 Hecke Operators on Maass forms

For any positive integer $n$, define

$$
Y_{n}=\left\{\left(\begin{array}{ll}
a & b \\
c & d
\end{array}\right): a, b, c, d \in \mathbb{Z}, a d-b c=n\right\}
$$

Note that there is a natural action of $Y_{1}=S L_{2}(\mathbb{Z})$ on $Y_{n}$. For any positive integer $n$ and $f \in \mathcal{C}(\Gamma \backslash \mathcal{H})$, define the $n^{\text {th }}$ Hecke operator acting on $\mathcal{C}(\Gamma \backslash \mathcal{H})$ as follows:

$$
T_{n}(f):=\frac{1}{\sqrt{n}} \sum_{a d=n} \sum_{0 \leq b<d} f\left(\frac{a z+b}{d}\right) .
$$

In the case of Maass forms these Hecke operators also satisfy several interesting properties. According to our need, we record some of the properties.

Theorem 2.7.1. For $n \geq 1$, let $T_{n}$ be the $n^{\text {th }}$ Hecke operator acting on $\mathcal{C}(\Gamma \backslash \mathcal{H})$. Then
(i) For $m, n \geq 1$,

$$
T_{m} T_{n}=\sum_{d \mid(m, n)} T_{\frac{m n}{d^{2}}} .
$$

In particular,

$$
T_{m} T_{n}=T_{n} T_{m}
$$

(ii) The Hecke operators commute with the Laplace operator $\triangle$.
(iii) The Hecke operators are self-adjoint operators.

Now again, since the Hecke operators $T_{n}$ together with the Laplacian $\triangle$ form a commutative family $\hbar$ of Hermitian operators on $L^{2}(\Gamma \backslash \mathcal{H})$ and the Hecke operators are self-adjoint in $L(\Gamma \backslash \mathcal{H})$, there exists a basis of $\mathcal{C}(\Gamma \backslash \mathcal{H})$ consisting of simultaneous eigenfunctions for all $T_{n}$.
Consider $\left\{u_{j}: j \geq 0\right\}$ to be a complete orthonormal basis for the space $\mathcal{C}(\Gamma \backslash \mathcal{H})$ consisting of the common eigenfunctions of $\hbar$, where $u_{0}$ is a constant function. Then

$$
\Delta u_{j}=\left(\frac{1}{4}+\left(t_{j}\right)^{2}\right) u_{j} \text { and } T_{n} u_{j}=\lambda_{j}(n) u_{j}
$$

where $\lambda_{j}(n) \in \mathbb{R}$. Selberg predicted that $t_{j}>0$. When the group is $S L_{2}(\mathbb{Z})$, it has been known for some time (Selberg, Roelcke) that $t_{j}>\frac{1}{4}$. The Fourier expansion of a Maass form, for $z=x+i y \in \mathcal{H}$

$$
u_{j}(z)=\sqrt{y} \rho_{j}(1) \sum_{n \neq 0} \lambda_{j}(n) K_{i t_{j}}(2 \pi|n| y) e(n x)
$$

where $\rho_{j}(1) \neq 0$ and $K_{v}$ is the modified Bessel function of the third kind, that is

$$
K_{v}=\frac{\pi}{2} \frac{I_{-v}(x)-I_{v}(x)}{\sin \pi \alpha}
$$

where

$$
I_{v}(x)=\sum_{m=0}^{\infty} \frac{1}{m!\Gamma(m+v+1)}\left(\frac{x}{2}\right)^{2 m+v} .
$$

A primitive Maass forms is a Maass form normalized by $\rho_{j}(1)$, that is, $\frac{u_{j}(z)}{\rho_{j}(1)}$. We know :

$$
\Omega(T):=\#\left\{j: 0<t_{j} \leq T\right\}=\frac{1}{4 \pi} \operatorname{vol}(\Gamma \backslash \mathcal{H}) T^{2}+\mathrm{O}(T \log T)
$$

where $\operatorname{vol}(\Gamma \backslash \mathcal{H})$ denote the volume of $\Gamma / \mathcal{H}$. The Ramanujan conjecture predicts that for any prime $p$,

$$
\left|\lambda_{j}(p)\right| \leq 2
$$

At present we are far from the above bound. The best bound towards the Ramanujan's conjecture for Maass forms is due to Kim and Sarnak (see [21]), that is,

$$
\left|\lambda_{j}(p)\right| \leq p^{\frac{7}{64}}+p^{-\frac{7}{64}} .
$$

Assuming the Ramanujan's conjecture, we can write

$$
\begin{equation*}
\lambda_{j}(p)=2 \cos \theta_{p, j}, \quad \theta_{p, j} \in[0, \pi] . \tag{2.10}
\end{equation*}
$$

### 2.8 Equidistribution of Hecke eigenvalues

Let $E=\left\{(x, y) \in \mathbb{R}^{2}: y^{2}=x^{3}+A x+B, A, B \in \mathbb{Q}\right\}$ be an elliptic curve with rational coefficients. Let us define the $j$-invariants of elliptic curves.

Definition 2.8.1. For any elliptic curve $E$, the $j$-invariant of $E$ is defined as:

$$
j(E)=1728 \frac{4 A^{3}}{4 A^{3}+27 B^{2}}
$$

Definition 2.8.2. The discriminant $\Delta_{E}$ of an elliptic curve $E$ is defined as

$$
\Delta_{E}=-16\left(4 A^{3}+27 B^{2}\right)
$$

The importance of $j$-invariant is that it determines $E$ up to isomorphism over an algebraically closed field. Let $E$ be an elliptic curve defined as in Definition 2.8.1 with integer coefficients, that is $A, B \in \mathbb{Z}$. For any prime $p$, let $N_{p}(E)=\left|E\left(\mathbb{F}_{p}\right)\right|$ be the number of solutions of $(E \bmod p)$ and define an integer $a_{p}(E)$ such that

$$
N_{p}(E)=p+1-a_{p}(E)
$$

Definition 2.8.3. A prime $p$ is said to be of good reduction if $p$ does not divide the discriminant $\Delta$.

From the result of Hasse we know that

$$
\left|a_{p}(E)\right| \leq 2 \sqrt{p}
$$

In the years 1950 and 1951, Sato and Tate independently studied the behaviour of $a_{p}^{\prime} \mathrm{s}$ varying $p$. Write

$$
a_{p}(E)=2 \sqrt{p} \cos \theta_{p}(E), \theta_{p}(E) \in[0, \pi] .
$$

Sato and Tate independently formulated the following conjecture:
Conjecture (Sato-Tate): Suppose $E$ has no complex multiplication. Then $a_{p}(E)\left(\right.$ resp. $\left.\theta_{p}(E)\right)$ are equidistributed in $[-1,1]$ (resp. in $\left.[0, \pi]\right)$ as $p \rightarrow \infty$ with respect to the measure

$$
\frac{2}{\pi} \sqrt{1-x^{2}} d x\left(\text { resp. } \frac{2}{\pi} \sin ^{2} \theta d \theta\right)
$$

In [40], Serre generalized the above conjecture as follows:
Conjecture (Sato-Tate): For a fixed $i$ such that $1 \leq i \leq s(N, k)$, Let $a_{p, i, N}$ be an eigenvalue of the $p^{t h}$ normalized Hecke operator $T_{p}$ acting on a non C.M $f \in S(N, k)$ such that $p$ is coprime to $N$. Then the family $\left\{a_{p, i, N}\right\}$ is equidistributed in $[-2,2]$ as $p \rightarrow \infty$ with respect to the Sato-Tate measure

$$
d \mu_{\infty}=\frac{1}{2 \pi} \sqrt{4-x^{2}} d x
$$

More precisely, the Sato-Tate conjecture is saying that for any continuous function $\phi:[-2,2] \rightarrow \mathbb{R}$ and interval $[\alpha, \beta] \subset[-2,2]$

$$
\lim _{V \rightarrow \infty} \frac{1}{V} \sum_{p \leq V} \phi\left(a_{p, i, N}\right)=\int_{\alpha}^{\beta} \phi d \mu_{\infty}
$$

Recently the above conjecture has been proved by Barnet-Lamb, Geraghty, Harris and Taylor [3], [10], [18] in a series of papers. In 1997, Serre [39] studied the "vertical" Sato-Tate conjecture by fixing a prime $p$ and varying $N$ and $k$. In particular, he proved the following theorem:

Theorem 2.8.4 (Serre). Let $N_{\lambda}, k_{\lambda}$ be positive integers such that $k_{\lambda}$ is even. $N_{\lambda}+k_{\lambda} \rightarrow \infty$ and $p$ is a prime not dividing $N_{\lambda}$ for any $\lambda$. Then the family of eigenvalues of the normalized $p^{\text {th }}$ Hecke operator

$$
T_{p}^{\prime}\left(N_{\lambda}, k_{\lambda}\right)=\frac{T_{p}\left(N_{\lambda}, k_{\lambda}\right)}{p^{\frac{k_{\lambda}-1}{2}}}
$$

is equidistributed in the interval $\Omega=[-2,2]$ with respect to the measure

$$
d \mu_{p}:=\frac{p+1}{\pi} \frac{\sqrt{1-\frac{x^{2}}{4}}}{\left(p^{\frac{1}{2}}+p^{-\frac{1}{2}}\right)^{2}-x^{2}} d x .
$$

Remark Also in 1997, Conrey, Duke and Farmer [13] studied a special case of above result by fixing $N=1$. In 2009, Murty and Sinha [32] investigated the effective / quantitative version of Serre's results, in which they give an explicit estimate on the rate of convergence. They proved the following theorem:

Theorem 2.8.5 (Murty-Sinha). Let $p$ be a fixed prime. Let $\{(N, k)\}$ be a sequence of pairs of positive integers such that $k$ is even and $N$ is not divisible by $p$. For any interval $[\alpha, \beta] \subset[-2,2]$,

$$
\frac{1}{s(N, k)} \sharp\left\{1 \leq i \leq s(N, k): a_{p, i, N} \in[\alpha, \beta]\right\}=\int_{\alpha}^{\beta} d \mu_{p}+\mathrm{O}\left(\frac{\log p}{\log k N}\right) \text {, }
$$

where the constant can be computed effectively.
In the case of Maass forms, the Sato-Tate conjecture is still open. In 1987, Sarnak [36, Theorem 1.2] studied a vertical version of Sato-Tate conjecture for primitive Maass forms. Assuming Ramanujan's bound in the case of primitive Maass forms, Sarnak's theorem can be interpreted as follows:

Theorem 2.8.6. [Sarnak] For a fixed prime $p$, the eigenvalues $\lambda_{j}(p), 1 \leq j \leq$ $r(T)$ of Hecke operators are equidistributed in $[-2,2]$ with respect to the measure $\mu_{p}$ as $T \rightarrow \infty$, where $r(T)$ denotes the number of Laplacian eigenvalues up to $T^{2}$ and

$$
\mu_{p}=\frac{p+1}{\pi} \frac{\sqrt{1-\frac{x^{2}}{4}}}{\left(p^{\frac{1}{2}}+p^{-\frac{1}{2}}\right)^{2}-x^{2}} d x
$$

Recently Lau and Wang [24], using the Kuznetsov trace formula made Sarnak's [36] result effective, that is, they proved the joint distribution of eigenvalues of the Hecke operators quantitatively for primitive Maass forms of level 1 and stated the same for primitive holomorphic cusp forms. They proved the following theorem:

Theorem 2.8.7 (Lau-Wang). Let $p$ be a prime. Let $k$ be positive even integer such that $\log (p) \leq \delta \log k$, for some small absolutely constant $\delta$. Let $\left\{\lambda_{j}(p)\right\}$ be the eigenvalues of normalized Hecke operators $T_{p}^{\prime}$ acting on space of primitive Maass forms. For any $I=[\alpha, \beta] \subset[-2,2]$

$$
\begin{gathered}
\frac{1}{\Omega(T)} \sharp\left\{0 \leq t_{j} \leq T:\left(\lambda_{j}(p)\right) \in I\right\} \\
=\int_{I} d \mu_{p}+\mathrm{O}\left(\frac{\log (p)}{\log k}\right),
\end{gathered}
$$

where $d \mu_{p}=\frac{p+1}{\pi} \frac{\sqrt{1-\frac{x^{2}}{4}}}{\left(p^{\frac{1}{2}}+p^{-\frac{1}{2}}\right)^{2}-x^{2}} d x$.
They have mentioned that their methods will work for higher level.

## Chapter 3

## Trace formulas

### 3.1 Eichler-Selberg trace formula

In this thesis, we will use the Eichler-Selberg trace formula as one of our important tools to prove the main theorems. This is a formula for the trace of Hecke operators $T_{n}$ acting on $S(N, k)$ in terms of class numbers of binary quadratic forms and certain arithmetic functions. In this section, we follow the presentation of [32]. Let $\Delta$ be a non-negative integer congruent to 0 or $1(\bmod 4)$, let $B(\triangle)$ be the set of all positive definite binary quadratic forms with discriminent $\Delta$, that is

$$
B(\Delta)=\left\{a x^{2}+b x y+c y^{2}: a, b, c \in \mathbb{Z}, a>0, b^{2}-4 a c=\Delta\right\} .
$$

We denote the set of primitive forms by

$$
b(\Delta)=\{f(x, y) \in B(\Delta): \operatorname{gcd}(a, b, c)=1\} .
$$

Let us define an action of the full modular group $S L_{2}(\mathbb{Z})$ on $B(\Delta)$ as follows:

$$
f(x, y)\left(\begin{array}{ll}
\alpha & \beta \\
\gamma & \delta
\end{array}\right):=f(\alpha x+\beta y, \gamma x+\delta y)
$$

Note that this action respects primitive forms. We know that the above action has finitely many orbits. We define $h(\Delta)$ to be the number of orbits
of $b(\Delta)$.
Let $h_{w}$ be defined as follows:

$$
\begin{aligned}
h_{w}(-3) & =\frac{1}{3} \\
h_{w}(-4) & =\frac{1}{2} \\
h_{w}(\Delta) & =h(\Delta) \text { for } \Delta<-4
\end{aligned}
$$

Theorem 3.1.1. For any positive integer $n$, the trace $\operatorname{Tr}$ of $T_{n}$ acting on $S(N, k)$ is given by

$$
\operatorname{Tr} T_{n}=A_{1}(n)+A_{2}(n)+A_{3}(n)+A_{4}(n),
$$

where $A_{i}(n) s$ are as follows:

$$
A_{1}(n)= \begin{cases}n^{\frac{k}{2}-1} \frac{k-1}{12} \psi(N) & \text { if } n \text { is a square }, \\ 0 & \text { otherwise. }\end{cases}
$$

where $\psi(N)=N \prod_{p \mid N}\left(1+\frac{1}{p}\right)$;

$$
A_{2}(n)=-\frac{1}{2} \sum_{t \in \mathbb{Z}, t^{2}<4 n} \frac{\varrho^{k-1}-\varrho^{k-1}}{\varrho-\bar{\varrho}} \sum_{f} h_{w}\left(\frac{t^{2}-4 n}{f^{2}}\right) \mu(t, f, n) ;
$$

where $\varrho$ and $\varrho$ are the zeros of the polynomial $x^{2}-t x+n$ and the inner sum runs over positive divisors $f$ of $\frac{\left(t^{2}-4 n\right)}{f^{2}} \in \mathbb{Z}$ is congruent to 0 or $1(\bmod 4) . \mu(t, f, n)$ is given by

$$
\mu(t, f, n)=\frac{\psi(N)}{\psi\left(\frac{N}{N_{f}}\right)} M\left(t, n, N N_{f}\right)
$$

where $N_{f}=\operatorname{gcd}(N, f)$ and $M(t, n, K)$ denote the number of solutions of the congruence $x^{2}-t x+n \equiv 0(\bmod K)$;

$$
\begin{equation*}
A_{3}(n)=-\sum_{d \mid n, 0<d \leq \sqrt{n}} d^{k-1} \sum_{c \mid N} \phi\left(\operatorname{gcd}\left(c, \frac{N}{c}\right)\right) \tag{3.1}
\end{equation*}
$$

Here, $\phi$ denotes Euler's function and in the first summation, if there is a contribution from the term $d=\sqrt{n}$, it should be multiplied by $\frac{1}{2}$. In the inner sum, we
also need the condition that $\operatorname{gcd}\left(c, \frac{N}{c}\right)$ divides $\operatorname{gcd}\left(N, \frac{n}{d}-d\right)$;

$$
A_{4}(n)= \begin{cases}\sum_{t \mid n, t>0} t & \text { if } k=2 \\ 0 & \text { otherwise } .\end{cases}
$$

Serre [39] and Murty-Sinha [32] used the above Eichler-Selberg trace formula to compute Weyl limits for the family of Hecke eigenangles. They also used it to estimate

$$
\mid \sum_{i=1}^{s(N, k)}\left(2 \cos m \theta_{p, i_{j}, N}-c_{m}(s(N, k)) \mid .\right.
$$

### 3.2 Kuznetsov trace formula

The following weighted version of Kuznetsov trace formula can be found in [24, Lemma 3.1]:

Theorem 3.2.1. [Kuznetsov] Let $m, n$ be two positive integers. Then for any $\epsilon>0$,

$$
\sum_{t_{j} \leq T} \alpha_{j} \lambda_{j}(m) \lambda_{j}(n)=\frac{T^{2}}{\pi^{2}} \delta_{m, n}+\mathrm{O}\left(T^{1+\epsilon}(m n)^{\frac{7}{64}}+(m n)^{\frac{1}{4}+\epsilon}\right)
$$

where $\alpha_{j}=\frac{\left(\rho_{j}(1)\right)^{2}}{\cosh \pi t_{j}}$ and $\delta_{m, n}$ is the Kronecker symbol.
Below we will state the unweighted version of Kuznetsov trace formula (see [24, Lemma 3.3])

Theorem 3.2.2. Let $k_{0}=\frac{11}{155}, \eta_{0}=\frac{43}{620}$. Let $m$, $n$ be positive integers. For $\epsilon>0$,
$\sum_{t_{j} \leq T} \lambda_{j}(m) \lambda_{j}(n)=\frac{1}{12} T^{2} \delta_{m, n=\square} \frac{\sigma((m, n))}{\sqrt{m n}}+\mathrm{O}_{\epsilon}\left(T^{2-k_{0}+\epsilon}(m n)^{\frac{7}{64}}+(m n)^{\eta_{0}+\epsilon}\right)$, where $\sigma(l)=\sum_{d \mid l} d$ and $\delta_{l=\square}=1$ if $l$ is a square and $\delta_{l=\square}=0$ otherwise.

In [24], Lau-Wang used the above Kuznetsov trace formula to get the estimates

$$
\left|\sum_{i=1}^{\Omega(T)} 2 \cos m \theta_{i}-c_{m} \Omega(T)\right| .
$$

## Chapter 4

## Distribution of gaps between equidistributed sequences

### 4.1 Distribution of gaps between equidistributed sequences

Let us start with a result that predicts the Weyl limits of gaps of equidistributed families. Henceforth, we will use $[x]$ as fractional part of $x$.

$$
[x]:=x-\lfloor x\rfloor,
$$

where $\lfloor x\rfloor$ is the largest integer less than equal to $x$.
Theorem 4.1.1. Consider $\left\{X_{1_{n}}\right\}_{n=1}^{\infty},\left\{X_{2_{n}}\right\}_{n=1}^{\infty}, \ldots,\left\{X_{r_{n}}\right\}_{n=1}^{\infty}$ to be collection of $r$ sequences of multisets in $[0,1]$ such that $1 \leq i \leq r, \# X_{i_{n}} \rightarrow \infty$ as $n \rightarrow \infty$. For every $m \in \mathbb{Z}$, let $c_{i_{m}}$ be the $m^{\text {th }}$ Weyl limit of $X_{i_{n}}$ respectively, that is,

$$
c_{i_{m}}:=\lim _{n \rightarrow \infty} \frac{1}{\# X_{i_{n}}} \sum_{t \in X_{i_{n}}} e(m t) .
$$

Assume that $c_{i_{m}}$ exists for each i. If $C_{m}$ is the $m^{\text {th }}$ Weyl limit of the family

$$
\left\{\left[x_{1}+x_{2}+\cdots+x_{r}\right], x_{i} \in X_{i_{n}}, 1 \leq i \leq r\right\}
$$

that is, for $m \in \mathbb{Z}$,

$$
C_{m}:=\lim _{n \rightarrow \infty} \frac{1}{\prod_{i=1}^{r} \# X_{i_{n}}} \sum_{\substack{x_{i} \in X_{i_{n}} \\ 1 \leq i \leq r}} e\left(m\left[x_{1}+x_{2}+\cdots+x_{r}\right]\right) .
$$

Then the Weyl limit

$$
\begin{equation*}
C_{m}=\prod_{i=1}^{r} c_{i_{m}} \tag{4.1}
\end{equation*}
$$

Remark 4.1.2. To consider the gaps, we can take the family

$$
\left\{\left[x_{1}-x_{2}-\cdots-x_{r}\right], x_{i} \in X_{i_{n}}, 1 \leq i \leq r\right\} .
$$

in the above theorem.
Further, if we consider the multisets such that if $x \in A_{i_{n}}$ then $-x \in A_{i_{n}}$. For simplicity, let us write $A_{i_{n}}=\left\{ \pm x_{i}\right\}$, we get the following corollary.

Corollary 4.1.3. In particular, if we consider the family $A_{i_{n}}=\left\{ \pm x_{i}\right\} \subseteq\left[-\frac{1}{2}, \frac{1}{2}\right]$, then we have

$$
c_{i_{m}}=\lim _{n \rightarrow \infty} \frac{1}{\# A_{i_{n}}} \sum_{t \in A_{i_{n}}} e(m[t]) .
$$

Let $C_{m}$ be the $m^{\text {th }}$ Weyl limit of the family

$$
\left\{\left[x_{1}+x_{2}+\cdots+x_{r}\right], x_{i} \in A_{i_{n}}, 1 \leq i \leq r\right\}
$$

that is for $m \in \mathbb{Z}$,

$$
C_{m}:=\lim _{n \rightarrow \infty} \frac{1}{\prod_{i=1}^{r} \# A_{i_{n}}} \sum_{\substack{x_{i} \in A_{i_{n}} \\ 1 \leq i \leq r}} e\left(m\left[x_{1}+x_{2}+\cdots+x_{r}\right]\right) .
$$

Then the Weyl limit

$$
\begin{equation*}
C_{m}=\prod_{i=1}^{r} c_{i_{m}} \tag{4.2}
\end{equation*}
$$

Proof. By the definition of the Weyl limits, we know

$$
C_{m}:=\lim _{n \rightarrow \infty} \frac{1}{\prod_{i=1}^{r} \# X_{i_{n}}} \sum_{\substack{x_{i} \in X_{i_{n}} \\ i=1,2, \ldots, r}} e\left(m\left[x_{1}+x_{2}+\cdots+x_{r}\right]\right)
$$

Observe that

$$
e(m x)=e(m[x]) .
$$

Using the above observation, we have

$$
\begin{aligned}
C_{m} & =\lim _{n \rightarrow \infty} \frac{1}{\prod_{i=1}^{r} \# X_{i_{n}}} \prod_{i=1}^{r} \sum_{x_{i} \in X_{i_{n}}} e\left(m x_{i}\right) \\
& =\prod_{i=1}^{r} c_{i_{m}} .
\end{aligned}
$$

In particular, for $A_{i_{n}}=\left\{ \pm x_{i}\right\} \subseteq\left[-\frac{1}{2}, \frac{1}{2}\right]$ the above calculation follows immediately, that is for any non-zero integer $m$, if $c_{i_{m}}$ be the $m^{\text {th }}$ Weyl limit of the family $\left\{ \pm x_{i}, x_{i} \in X_{i_{n}}, 1 \leq i \leq r\right\}$ and $C_{m}$ be the $m^{t h}$ Weyl limit of the family $\left\{\left[x_{1}+x_{2}+\cdots+x_{r}\right], x_{i} \in A_{i_{n}}, 1 \leq i \leq r\right\}$, then

$$
C_{m}=\prod_{i=1}^{r} c_{i_{m}}
$$

In the next theorem, we will predict the measure with respect to which the above mentioned family in Theorem 4.1.1 is equidistributed.

Theorem 4.1.4. Consider $\left\{A_{1_{n}}\right\}_{n=1}^{\infty},\left\{A_{2_{n}}\right\}_{n=1}^{\infty}, \ldots,\left\{A_{r_{n}}\right\}_{n=1}^{\infty} \subseteq\left[-\frac{1}{2}, \frac{1}{2}\right]$ to be sequences of multisets such that $-x_{i} \in A_{i_{n}}$ whenever $x_{i} \in A_{i_{n}}$, and $\# A_{i_{n}} \rightarrow$ $\infty$ as $n \rightarrow \infty$ for $i=1,2,3, \ldots, r$. If $\left\{A_{i_{n}}\right\}, i=1,2, \ldots, r$ are equidistributed in $\left[-\frac{1}{2}, \frac{1}{2}\right]$ with respect to the measure $F_{i}(x) d x$ respectively, where

$$
F_{i}(x)=\sum_{m=-\infty}^{\infty} c_{i_{m}} e(m x)
$$

then the family

$$
\left\{\left[x_{1}+x_{2}+\cdots+x_{r}\right], x_{i} \in A_{i_{n}}\right\}
$$

is equidistributed in $[0,1]$ with respect to the measure

$$
\mu=F(x) d x
$$

where

$$
F(x)=\sum_{m=-\infty}^{\infty} C_{m} e(m x)
$$

Moreover, if

$$
\begin{equation*}
\sum_{m=-\infty}^{\infty}\left|c_{i_{m}}\right|^{2}<\infty \text { for all } 1 \leq i \leq r \tag{4.3}
\end{equation*}
$$

then the above function $F(x)$ equals

$$
F_{1} * F_{2} * \cdots * F_{r}(x)
$$

where

$$
\begin{gathered}
F_{1} * F_{2} * \cdots * F_{r}(y) \\
=\int_{0}^{1} \ldots \int_{0}^{1} F_{1}\left(y_{1}\right) F_{2}\left(y_{2}\right) \ldots F_{r}\left(y-y_{1}-y_{2}-\cdots-y_{r-1}\right) d y_{r-1} d y_{r-2} \ldots d y_{1}
\end{gathered}
$$

Remark 4.1.5. In general an equidistributed family may not satisfy (4.3).
Proof. Since $\left\{A_{i_{n}}\right\}$ are equidistributed in $\left[-\frac{1}{2}, \frac{1}{2}\right]$, with respect to the measure $F_{i}(x) d x$, by Theorem 2.1.10

$$
\lim _{V \rightarrow \infty} \frac{1}{V} \sum_{|m| \leq V}\left|c_{i_{m}}\right|^{2}=0 \text { for } 1 \leq i \leq r
$$

Observe that, if there exist a subsequence say $\left\{y_{n}\right\}$ of $\left\{\left|c_{i_{m}}\right|^{2}\right\}$, such that $\left|y_{n}\right| \geq 1$ for all $n$, then

$$
\sum_{n \leq V}\left|y_{n}\right|^{2} \geq V
$$

and hence,

$$
\sum_{n \leq V}\left|c_{i_{m}}\right|^{2} \geq V
$$

which is a contradiction. Hence, $\left|c_{i_{m}}\right|<1$, except possibly for finitely many $m$.
Now using Theorem 4.1.1 and above fact, we have

$$
\left|C_{m}\right|<\left|c_{i_{m}}\right|
$$

except possibly for finitely many $m$ and for all $1 \leq i \leq r$. Hence,

$$
\sum_{|m| \leq V}\left|C_{m}\right|^{2} \leq \sum_{|m| \leq V}\left|c_{i_{m}}\right|^{2}
$$

Now taking the limit, we have

$$
\begin{equation*}
\lim _{V \rightarrow \infty} \frac{1}{V} \sum_{|m| \leq V}\left|C_{m}\right|^{2}=0 \tag{4.4}
\end{equation*}
$$

Hence, by Theorem 2.1.10, we can conclude that

$$
\left\{\left[x_{1}+x_{2}+\cdots+x_{r}\right]\right\}
$$

is equidistributed in $[0,1]$ with respect to the measure

$$
\mu=F(x) d x
$$

where

$$
F(x)=\sum_{m=-\infty}^{\infty} C_{m} e(m x)
$$

In addition, if the concerned family satisfies (4.3), that is

$$
\sum_{m=-\infty}^{\infty}\left|c_{i_{m}}\right|^{2}<\infty, 1 \leq i \leq r
$$

then $\left|c_{i_{m}}\right|<1$, except possibly for finitely many $m$.
Using Theorem 4.1.1 and the last line, we have

$$
\left|C_{m}\right|^{2} \leq\left|c_{i_{m}}\right|^{2} \text { except possibly finitely many } m
$$

Hence,

$$
\sum_{m=-\infty}^{\infty}\left|C_{m}\right|^{2}<\infty
$$

Hence, by Theorem 2.2.1, there exist a function $F \in L^{2}([0,1])$ such that

$$
\hat{F}(m)=C_{m} \text { for all } \mathrm{m}
$$

By uniqueness of $F$, the above function has to be

$$
F(x)=\sum_{m} C_{m} e(m x) .
$$

But note that

$$
C_{m}=\prod_{i=1}^{r} c_{i_{m}} \text { and } \prod_{i=1}^{r} \hat{F}_{i}(m)=\left(F_{1} * F_{2} * \cdots * F_{r}\right)^{\wedge}(m)
$$

Hence,

$$
F(x)=F_{1} * F_{2} * \cdots * F_{r}(x) .
$$

Let $\left\{A_{i_{n}}\right\}_{n=1}^{\infty}$ be equidistributed sequences of finite multisets. If we know the distribution effectively, then our next result will help us to predict the effective equidistribution of family of gaps of equidistributed families.

Theorem 4.1.6. Let $\left\{A_{1_{n}}\right\}_{n=1}^{\infty},\left\{A_{2_{n}}\right\}_{n=1}^{\infty}, \ldots,\left\{A_{r_{n}}\right\}_{n=1}^{\infty}$ be sequences of finite multisets as defined in Theorem 4.1.4. Let $\left\{A_{i_{n}}\right\}_{n=1}^{\infty}, 1 \leq i \leq r$ be equidistributed sequences in $\left[-\frac{1}{2}, \frac{1}{2}\right]$ with respect to the measure $F_{i}(x) d x$, where

$$
F_{i}(x)=\sum_{m=-\infty}^{\infty} c_{i_{m}} e(m x)
$$

and $c_{i_{m}}$ are as defined in Theorem 4.1.1.
Consider $\underline{x}=\left(x_{1}, x_{2}, . ., x_{r}\right), A_{\underline{n}}=A_{1_{n}} \times A_{2_{n}} \times \cdots \times A_{r_{n}}$.
Then, for any positive integer $M$ and any $I=[\alpha, \beta] \subseteq[0,1]$, we have

$$
\begin{gathered}
\left|\frac{1}{\prod_{i=1}^{r}\left(\# A_{i_{n}}\right)} \#\left\{\underline{x} \in A_{\underline{n}}:\left[x_{1}+x_{2}+\cdots+x_{r}\right] \in I\right\}-\int_{I} \mu\right| \\
\leq \frac{\prod_{i=1}^{r} \# A_{i_{n}}}{M+1} \\
+\sum_{|m| \leq M}\left(\frac{1}{M+1}+\min \left(\beta-\alpha, \frac{1}{\pi|m|}\right)\right)\left(\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left(m x_{i}\right)-\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}\right|\right)
\end{gathered}
$$

where $\mu$ is as defined in Theorem 4.1.4.

For proof of the above theorem see the next chapter.

## Chapter 5

## A Variant of the Erdös-Turán Inequality

The following variant of Erdös-Turán inequality will be very useful for getting the effective equidistribution of eigenvalues of Hecke operators.

Theorem 5.0.7. Let $I=[a, b]$ be an interval of $\mathbb{R}$. For $1 \leq i \leq r$, let $\left\{A_{i_{n}}\right\}$ be defined as in Theorem 4.1.4. Let $V_{n}:=\prod_{i=1}^{r} \# A_{i_{n}}$ and

$$
x_{\underline{n}}:=\left[ \pm x_{1} \pm x_{2} \pm \ldots \pm x_{r}\right], x_{i} \in A_{i_{n}} .
$$

Let $C_{\underline{m}}$ the Weyl limit defined as in Definition 2.1.8. Let

$$
\mu=F(x) d x
$$

where

$$
F(x)=\sum_{m=-\infty}^{\infty} C_{m} e(m x)
$$

Define

$$
N_{I}\left(V_{n}\right):=\#\left\{\left(x_{1}, x_{2}, \ldots, x_{r}\right) \in A_{1_{n}} \times A_{2_{n}} \times \ldots \times A_{r_{n}}: x_{\underline{n}} \in I\right\}
$$

and

$$
D_{I, V_{n}}(\mu):=\left|N_{I}\left(V_{n}\right)-V_{n} \mu(I)\right| .
$$

Then,

$$
\begin{aligned}
& \qquad D_{I, V_{n}}(\mu) \leq \frac{V_{n}\|\mu\|}{M+1} \\
& +\sum_{1 \leq|m| \leq M}\left(\frac{1}{M+1}+\min \left(b-a, \frac{1}{\pi|m|}\right)\right)\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left( \pm m x_{i}\right)-\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}\right|, \\
& \text { where }\|\mu\|=\operatorname{Sup}_{x \in[0,1]}|F(x)| .
\end{aligned}
$$

Proof. Let $\chi_{I}$ be the characteristic function of the interval $I$. Then by (a) of Section 2.3, we have

$$
\sum_{\substack{x_{i} \in A_{i} \\ 1 \leq i \leq r}} S_{M}^{-}\left(x_{\underline{n}}\right) \leq \sum_{\substack{x_{i} \in A_{i} \\ 1 \leq i \leq r}} \chi_{I}\left(x_{\underline{n}}\right) \leq \sum_{\substack{x_{i} \in i_{i} \\ 1 \leq i \leq r}} S_{M}^{+}\left(x_{\underline{n}}\right) .
$$

Now using the Fourier expansion of $S_{M}^{ \pm}\left(x_{\underline{n}}\right)$, we know that

$$
\sum_{\substack{x_{i} \in A_{i} \\ 1 \leq i \leq r}} S_{M}^{ \pm}\left(x_{\underline{n}}\right)=\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m)\left(\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left( \pm m x_{i}\right)\right) .
$$

Subtracting $\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}$ from the inner exponential sums, we get

$$
\begin{align*}
& \sum_{\substack{x_{i} \in A_{i_{n}} \\
1 \leq i \leq r}} S_{M}^{ \pm}\left(x_{\underline{n}}\right)-\left(\prod_{i=1}^{r} \# A_{i_{n}}\right) \sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) C_{m}  \tag{5.1}\\
= & \sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m)\left(\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left( \pm m x_{i}\right)-\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}\right) .
\end{align*}
$$

Since $c_{i_{0}}=1$,

$$
\left(\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left( \pm m x_{i}\right)-\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}\right)=0 \text { for } m=0
$$

Taking the absolute value on both sides we get

$$
\begin{aligned}
& \left|\sum_{\substack{x_{i} \in i_{i n} \\
1 \leq i \leq r}} S_{M}^{ \pm}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r} \# A_{i_{n}} \sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) C_{m}\right| \\
\leq & \sum_{|m| \leq M}\left|\hat{S}_{M}^{ \pm}(m)\right|\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left( \pm m x_{i}\right)-\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}\right| .
\end{aligned}
$$

Now let us consider the sum

$$
\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m) C_{m}
$$

Since for all $|m|>M, \hat{S}_{M}^{ \pm}(m)=0$, without loss of generality let us extend the range of sums to $\mathbb{Z}$. Then, we have

$$
\sum_{m} \hat{S}_{M}^{ \pm}(m) C_{m}=\sum_{m} C_{m} \int_{0}^{1} S_{M}^{ \pm}(x) e(-m x) d x
$$

Now interchanging the sum and integral and using the definition of $\mu$, the above quantity equals

$$
\int_{0}^{1} S_{M}^{ \pm}(x) d \mu
$$

Using (c) of Section 2.3, we have

$$
\begin{equation*}
\left|\int_{0}^{1}\left(S_{M}^{ \pm}(x)-\chi_{I}(x)\right) d \mu\right| \leq \frac{\|\mu\|}{M+1} . \tag{5.2}
\end{equation*}
$$

Note that

$$
\begin{aligned}
& D_{I, V_{n}}(\mu)=\left|N_{I}(V)-\left(\prod_{i=1}^{r} A_{i_{n}}\right) \mu(I)\right| \\
= & \left|\sum_{\substack{x_{i} \in A_{i} \\
1 \leq i \leq r}} \chi_{I}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} \chi_{I}(x) d \mu\right| .
\end{aligned}
$$

Now adding and subtracting $\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu$ to the above expression, we get

$$
\begin{aligned}
& =\mid \prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} \chi_{I}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu \\
+ & \prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} \chi_{I}(x) d \mu \mid .
\end{aligned}
$$

Using triangle inequality, we get

$$
\begin{aligned}
& D_{I, V}(\mu) \leq\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} \chi_{I}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu\right| \\
& +\left|\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} \chi_{I}(x) d \mu\right| \\
& \leq\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} \chi_{I}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu\right| \\
& +\left|\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1}\left(S_{M}^{+}(x)-\chi_{I}(x)\right) d \mu\right|
\end{aligned}
$$

Now using (5.2), the above is

$$
\leq \frac{\prod_{i=1}^{r}\left(\# A_{i_{n}}\|\mu\|\right)}{M+1}+\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} \chi_{I}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) \int_{0}^{1} S_{M}^{+}(x) d \mu\right|
$$

Using (a) of Section 2.3, we have

$$
\begin{gathered}
\leq \frac{\prod_{i=1}^{r}\left(\# A_{i_{n}}\|\mu\|\right)}{M+1} \\
+\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} S_{M}^{+}\left(x_{\underline{n}}\right)-\prod_{i=1}^{r} \# A_{i_{n}} \int_{0}^{1} S_{M}^{+}(x) d \mu\right| .
\end{gathered}
$$

Using (5.1) and the fact that

$$
\int_{0}^{1} S_{M}^{+}(x) d \mu=\sum_{|m| \leq M} \hat{S}_{M}^{+}(m) C_{m}
$$

$D_{I, V}(\mu)$ is

$$
\leq \frac{\prod_{i=1}^{r}\left(\# A_{i_{n}}\|\mu\|\right)}{M+1}+\sum_{|m| \leq M} \hat{S}_{M}^{ \pm}(m)\left(\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left( \pm m x_{i}\right)-\prod_{i=1}^{r} \# A_{i_{n}} c_{i_{m}}\right)
$$

Now using (b) of Section 2.3, we have

$$
\begin{gathered}
\left|D_{I, V}(\mu)\right| \leq \frac{\prod_{i=1}^{r}\left(\# A_{i_{n}}\|\mu\|\right)}{M+1} \\
+\sum_{|m| \leq M} \frac{1}{M+1}+\min \left(b-a, \frac{1}{\pi|m|}\right)\left|\prod_{i=1}^{r} \sum_{x_{i} \in A_{i_{n}}} e\left(m x_{i}\right)-\prod_{i=1}^{r}\left(\# A_{i_{n}}\right) C_{m}\right|
\end{gathered}
$$

## Chapter 6

## Distribution of gaps between eigenangles of Hecke operators

In the following sections we will state several applications of Theorem 4.1.1, 4.1.4 and 4.1.6.

### 6.1 Distribution of gaps between eigenangles for cusp forms

In this section, we use the notations from the Section 2.5. In this section we study the distribution of gaps between eigenangles of Hecke operators acting on the space $S(N, k)$. In this case, let $n=s(N, k)$ and for all $1 \leq j \leq$ $r$,

$$
A_{j_{n}}=\left\{\frac{\theta_{p, i, N}}{2 \pi}, 1 \leq i \leq s(N, k)\right\}
$$

Note that

$$
\# A_{j_{n}}=n=s(N, k)
$$

So as $N+k \rightarrow \infty$, we have $n \rightarrow \infty$ that is we have infinite number of multi sets. Since each sets are same and we are going to study the distributions of gaps between the elements of the multisets $A_{j_{n}}$, it is better to introduce
another subscript $j$ to $i$ that is we will write

$$
A_{j_{n}}=\left\{\frac{\theta_{p, i_{j}, N}}{2 \pi}, 1 \leq i \leq s(N, k)\right\} .
$$

The subscript $j$ to $i$ is for the convenience of the reader.
Theorem 6.1.1. Let $N$ be a positive integer and $p$ a prime not dividing $N$. For an interval $[\alpha, \beta] \subseteq[0,1], r \leq s(N, k)$,

$$
\begin{gathered}
\frac{1}{s(N, k)^{r}} \#\left\{1 \leq i_{1}, \ldots, i_{r} \leq s(N, k):\left[\frac{ \pm \theta_{p, i_{1}, N} \pm \cdots \pm \theta_{p, i_{r}, N}}{2 \pi}\right] \in[\alpha, \beta]\right\} \\
=\int_{[\alpha, \beta]} \nu_{p}+\mathrm{O}\left(\frac{\log p}{\log k N}\right)
\end{gathered}
$$

where

$$
\nu_{p}=F(x) * F(x) * \cdots * F(x) d x
$$

$\theta_{p, i, N} \in[0, \pi]$ such that $a_{p, i, N}=2 \cos \theta_{p, i, N}$ and

$$
F(x)=4(p+1) \frac{\sin ^{2} 2 \pi x}{\left(p^{\frac{1}{2}}+p^{-\frac{1}{2}}\right)^{2}-\cos ^{2} 2 \pi x}
$$

Here the implied constant is effectively computable.
Remark 6.1.2. For $r=2$, the above mentioned measure is

$$
\nu_{p}=\frac{2(1+\cos 4 \pi x)\left(1-\frac{1}{p^{2}}\right)+\frac{4}{p}\left(\frac{1}{p^{2}}-\cos 4 \pi x\right)}{1+\frac{1}{p^{4}}-\frac{2}{p^{2}} \cos 4 \pi x} d x .
$$

The following theorem can be deduced from Theorem 6.1.1.
Theorem 6.1.3. For any $\alpha \in[0,1]$,

$$
\begin{aligned}
\#\left\{1 \leq i_{1}, i_{2}, \ldots, i_{r}\right. & \left.\leq s(N, k):\left[ \pm \theta_{p, i_{1}, N} \pm \theta_{p, i_{2}, N} \pm \cdots \pm \theta_{p, i_{r}, N}\right]=\alpha\right\} \\
& =\mathrm{O}\left((s(N, k))^{r}\left(\frac{\log p}{\log k N}\right)\right)
\end{aligned}
$$

where the implied constant is effectively computable.

In the above theorem for $r=2$, we have an interesting consequence, namely

Theorem 6.1.4. For any $\alpha \in[0,1]$,

$$
\#\left\{\left(i_{1}, i_{2}\right):\left[\frac{ \pm \theta_{p, i_{1}, N} \pm \theta_{p, i_{2}, N}}{2 \pi}\right]=\alpha\right\}=\mathrm{O}\left((s(N, k))^{2}\left(\frac{\log p}{\log k N}\right)\right) .
$$

In particular, taking $\alpha=0$ and using the fact that

$$
\#\left\{\left(i_{1}, i_{2}\right):\left(\frac{ \pm \theta_{p, i_{1}, N} \pm \theta_{p, i_{2}, N}}{2 \pi}\right)=\alpha\right\} \leq \#\left\{\left(i_{1}, i_{2}\right):\left\{\frac{ \pm \theta_{p, i_{1}, N} \pm \theta_{p, i_{2}, N}}{2 \pi}\right\}=\alpha\right\}
$$

we have the following corollary which recovers Theorem 1 in [33].

## Corollary 6.1.5.

$$
\begin{aligned}
\#\left\{\left(i_{1}, i_{2}\right), 1\right. & \left.\leq i_{1}, i_{2} \leq s(N, k):\left(\theta_{p, i_{1}, N} \pm \theta_{p, i_{2}, N}\right)=0\right\} \\
& =\mathrm{O}\left((s(N, k))^{2}\left(\frac{\log p}{\log k N}\right)\right)
\end{aligned}
$$

Using the results from Chapter 2, we know that $\left\{\frac{ \pm \theta_{p, i_{j}, N}}{2 \pi}, 1 \leq j \leq r\right\}$ is equidistributed in $\left[-\frac{1}{2}, \frac{1}{2}\right]$ with respect to the measure

$$
\mu_{p}=F(x) d x
$$

where

$$
\begin{equation*}
F(x)=4(p+1) \frac{\sin ^{2} 2 \pi x}{\left(p^{\frac{1}{2}}+p^{-\frac{1}{2}}\right)^{2}-4 \cos ^{2} 2 \pi x} \tag{6.1}
\end{equation*}
$$

So using Theorem 4.2, we can conclude that, the concerned family

$$
\left[\frac{ \pm \theta_{p, i_{1}, N} \pm \cdots \pm \theta_{p, i_{r}, N}}{2 \pi}\right]
$$

is equidistributed in $[0,1]$ with respect to the measure $\underbrace{F * F \cdots * F}_{r \text { times }} d x$.
To proceed further, let us prove the following proposition:

Proposition 6.1.6. For any positive integer $m$,

$$
\begin{aligned}
& \left|\prod_{j=1}^{r} \sum_{\substack{i \\
1 \leq i \leq s(N, k)}} 2 \cos m \theta_{p, i_{j}, N}-C_{m}(s(N, k))^{r}\right| \\
& \ll p^{\frac{3 r m}{2}} m^{r}(\log p)^{r} 2^{r \nu(N)}+(\sqrt{N} d(N))^{r}
\end{aligned}
$$

where $\nu(N)$ is the number of distinct prime factor of $N$ and $d(N)$ is the divisor function. The constant is effectively computable.

Proof. Estimating each term of the Eichler-Selberg trace formula, Murty and Sinha (see [32] Theorem 18 and (11)) prove the following: For any positive integer $m$, let $c_{i_{m}}$ be the Weyl limits of the family

$$
\left\{ \pm \theta_{p, i_{j}, N}, 1 \leq j \leq s(N, k)\right\}
$$

For $1 \leq i \leq r$, the Weyl limits $c_{i_{m}}$ are given by

$$
c_{i_{m}}= \begin{cases}1 & \text { if } m=0  \tag{6.2}\\ \left(\frac{1}{p^{\frac{m}{2}}}-\frac{1}{p^{\frac{m-2}{2}}}\right) & \text { if } m \text { is even } \\ 0 & \text { otherwise }\end{cases}
$$

Moreover, for $m \geq 1$ and $1 \leq j \leq r$,

$$
\left\lvert\, \sum_{i=1}^{s(N, k)}\left(2 \cos m \theta_{p, i_{j}, N}-c_{i_{m}}(s(N, k)) \left\lvert\, \ll p^{\frac{3 m}{2}} 2^{\nu(N)} \log p^{m}+d(N) \sqrt{N} .\right.\right.\right.
$$

Using the fact that

$$
\begin{equation*}
2 \cos m \theta_{p, i_{j}, N}=X_{m}\left(2 \cos \theta_{p, i_{j}, N}\right)-X_{m-2}\left(2 \cos \theta_{p, i_{j}, N}\right), \quad m \geq 2 \tag{6.3}
\end{equation*}
$$

where

$$
X_{m}(2 \cos \theta)=\frac{\sin (m+1) \theta}{\sin \theta}
$$

we have

$$
\prod_{j=1}^{r} \sum_{\substack{i \\ 1 \leq i \leq s(N, k)}} 2 \cos m \theta_{p, i_{j}, N}=\left(\operatorname{Tr} T_{p^{m}}^{\prime}-\operatorname{Tr} T_{p^{m-2}}^{\prime}\right)^{r}
$$

Now, using the estimates of Eichler-Selberg trace formula on page 696 of [32] and the well-known inequality

$$
\begin{equation*}
(a-b)^{r} \leq r\left(a^{r}+b^{r}\right) \tag{6.4}
\end{equation*}
$$

we have

$$
\begin{gathered}
\left(\operatorname{Tr} T_{p^{m}}^{\prime}-\operatorname{Tr} T_{p^{m-2}}^{\prime}\right)^{r} \ll r\left(\left(\frac{k-1}{12}\right) \psi(N)\left(\frac{1}{p^{\frac{m}{2}}}-\frac{1}{p^{\frac{m-1}{2}}}\right)\right)^{r} \\
+\left(p^{\frac{3 m}{2}} m(\log p) 2^{\nu(N)}+\sqrt{N} d(N)\right)^{r} .
\end{gathered}
$$

Using (6.4) and the fact that $s(N, k)=\mathrm{O}(k N)$, we get

$$
\begin{aligned}
& \left|\prod_{j=1}^{r} \sum_{\substack{i \\
1 \leq i \leq s(N, k)}} 2 \cos m \theta_{p, i_{j}, N}-\left(c_{m}(s(N, k))\right)^{r}\right| \\
& \quad \ll p^{\frac{3 r m}{2}} m^{r}(\log p)^{r} 2^{r \nu(N)}+(\sqrt{N} d(N))^{r} .
\end{aligned}
$$

Proof of Theorem 6.1.1 Using Theorem 4.1.6, the concerned quantity is

$$
\ll \frac{(s(N, k))^{r}}{M+1}+p^{\frac{3 r M}{2}} M^{r}(\log p)^{r} 2^{r \nu(N)}+(\sqrt{N} d(N))^{r} .
$$

Now we want to choose $M$ such that

$$
\frac{(s(N, k))^{r}}{M+1} \sim p^{\frac{3 r M}{2}} .
$$

And that can be achieved by choosing $M=c \frac{\log k N}{\log p}$ for a sufficiently small constant $c$. The result then follows.

### 6.2 Distribution of gaps between eigenangles for primitive Maass forms

For notations used in this section, the reader may see Section 2.6. In this case, for $1 \leq i \leq r$,

$$
A_{j_{i}}=\left\{0<t_{j} \leq T: \frac{ \pm \theta_{j}}{2 \pi}\right\} .
$$

Note the cardinality of each $A_{j_{i}}$ are $\Omega(T)$. Since each multi sets are same, for our convenience let us introduce a suffix to $j$ that is

$$
A_{j_{i}}=\left\{0<t_{j_{i}} \leq T: \frac{ \pm \theta_{j_{i}}}{2 \pi}\right\} .
$$

Theorem 6.2.1. There exists a small constant $\delta>0$ such that for all large $T$,

$$
\begin{gathered}
\frac{1}{(\Omega(T))^{r}} \#\left\{0<t_{j_{i}} \leq T, 1 \leq i \leq r:\left[\frac{ \pm \theta_{j_{1}}(p) \pm \theta_{j_{2}}(p) \cdots \pm \theta_{j_{r}}(p)}{2 \pi}\right] \in I\right\} \\
=\int_{I} \nu_{p} d x+\mathrm{O}\left(\frac{\log p}{\log T}\right)
\end{gathered}
$$

holds uniformly for integers $1 \leq r \leq \Omega(T)$ and for primes $p$ satisfying

$$
r \log p \leq \delta \log T
$$

and uniformly for any interval $I=[a, b] \subseteq[0,1]$. Here

$$
\nu_{p}=\underbrace{F_{p}(x) * F_{p}(x) * \cdots * F_{p}}_{r \text { times }}(x),
$$

and $F_{p}(x)$ is as defined in (6.1). Here the implied constant is effectively computable and $\theta_{j_{i}}(p)$ are called as eigenangles of Hecke operators acting on the space of primitive Maass forms.

As a consequence of the above Theorem 6.2.1, we have the following result

Theorem 6.2.2. For any $\alpha \in[0,1]$,

$$
\begin{aligned}
\#\left\{0<t_{j_{i}} \leq T, 1\right. & \left.\leq i \leq r:\left[\frac{ \pm \theta_{j_{1}}(p) \pm \theta_{j_{2}}(p) \cdots \pm \theta_{j_{r}}(p)}{2 \pi}\right]=\alpha\right\} \\
& =\mathrm{O}\left((\Omega(T))^{r}\left(\frac{\log p}{\log T}\right)\right)
\end{aligned}
$$

Corollary 6.2.3. For any $\alpha=0$ and $r=2$,

$$
\begin{gathered}
\#\left\{0<t_{j_{i}} \leq T, 1 \leq i \leq r:\left[\frac{ \pm \theta_{j_{1}}(p) \pm \theta_{j_{2}}(p)}{2 \pi}\right]=0\right\} \\
=\mathrm{O}\left((\Omega(T))^{r}\left(\frac{\log p}{\log T}\right)\right)
\end{gathered}
$$

Remark 6.2.4. From above corollary, we can derive similar result like Corollary 6.1.5.

In order to prove Theorem 6.2.1, we first need the following result.
Proposition 6.2.5. For $m \geq 1,1 \leq i \leq r \leq \Omega(T), 0<k<\frac{11}{155}$ and $\eta>\frac{43}{620}$,

$$
\begin{gathered}
\left|\prod_{j=1}^{r} \sum_{i=1}^{\Omega(T)} \cos m \theta_{j_{i}}(p)-\Omega(T)^{r} C_{m}\right| \\
\ll\left(T^{2-k} p^{m \eta}\right)^{r}
\end{gathered}
$$

where the implied constant depends only on $\eta$.
Proof. The following result is a special case of [24], Lemma 4.1:
For any positive integer $m$, let $c_{m}$ be the Weyl limits of the family

$$
\left\{\frac{ \pm \theta_{i}(p)}{2 \pi}, 1 \leq i \leq \Omega(T)\right\}
$$

Then the Weyl limits $c_{m}$ are given by

$$
c_{m}= \begin{cases}1 & \text { if } m=0  \tag{6.5}\\ \left(\frac{1}{p^{\frac{m}{2}}}-\frac{1}{p^{\frac{m-2}{2}}}\right) & \text { if } m \text { is even } \\ 0 & \text { otherwise }\end{cases}
$$

Moreover, for $m \geq 1$,

$$
\left|\sum_{i=1}^{\Omega(T)} 2 \cos m \theta_{i}-c_{m} \Omega(T)\right| \ll T^{2-k} p^{m \eta}
$$

Lau and Wang [24] proved the above result using Kuznetsov trace formula. Now proceeding as in the proof of Propositions 6.1.6, the result follows.

Proof. From Sarnak's theorem (see [36, Theorem 1.2]), we can conclude that the family $\left\{\frac{ \pm \theta_{i}(p)}{2 \pi}, 1 \leq i \leq r\right\}$ is equidistributed in $\left[-\frac{1}{2}, \frac{1}{2}\right]$ with respect to $F_{p}(x) d x$.

Now proceeding just like proof of Theorems 6.1.1 and choosing $M=$ $c \frac{\log p}{\log T}$, we have Theorem 6.2.1.

## Chapter 7

## Distinguishing Hecke eigenforms

The content of this chapter is joint work with M. Ram Murty (see [30]). This work is independent of the previous chapters. In this chapter, we will follow closely the presentation of [30].

### 7.1 Introduction

Let $E=\left\{(X, Y): Y^{2}=X^{3}+A X+B, A, B \in \mathbb{Z}\right\}$ be an elliptic curve. For any prime $p$, let $\left|E\left(\mathbb{F}_{p}\right)\right|$ be the number of solutions of

$$
Y^{2}=X^{3}+A X+B(\bmod p)
$$

together with the point at $\infty$.
Definition 7.1.1. A set of primes $\mathcal{P}$ is said to have positive upper density if

$$
\limsup _{x \rightarrow \infty} \frac{\#\{p \in \mathcal{P}: p \leq x\}}{\#\{p \leq x\}}>0
$$

Recently, using Galois theory, Kulkarni, Patankar and Rajan [23] proved that if $E_{1}$ and $E_{2}$ are two elliptic curves defined over a number field $K$ with at least one of them not of CM type such that

$$
\# E_{1}\left(\mathbb{F}_{p}\right)=\# E_{2}\left(\mathbb{F}_{p}\right)
$$

for a set of primes of positive lower density, then $E_{1}$ and $E_{2}$ become isogenous after base change. Given two normalized Hecke eigenforms $f_{1}$ and $f_{2}$ of weights $k_{1}, k_{2}$ and levels $N_{1}, N_{2}$ respectively, let

$$
\begin{equation*}
f_{i}(z)=\sum_{n=1}^{\infty} a_{n}\left(f_{i}\right) n^{\left(k_{i}-1\right) / 2} q^{n}, \quad q=e^{2 \pi i z}, \quad i=1,2, \tag{7.1}
\end{equation*}
$$

be the Fourier expansions at infinity. By the work of Deligne [14], we can write

$$
a_{n}\left(f_{i}\right)=2 \cos \theta_{n}^{(i)}, \theta_{n}^{(i)} \in[0, \pi] .
$$

Definition 7.1.2. A Dirichlet character is a group homomorphism from $(\mathbb{Z} / n \mathbb{Z})^{*}$ to $\mathbb{C}^{*}$, where $(\mathbb{Z} / n \mathbb{Z})^{*}$ is the group of units of $(\mathbb{Z} / n \mathbb{Z})$ and $\mathbb{C}^{*}$ is the set of non zero complex numbers.

Definition 7.1.3. An eigenform $f$ is said to be of $C M$ type if there exists an imaginary quadratic field $K / \mathbb{Q}$ such that $a_{p}(f)=0$ whenever $p$ is inert in $K$.

In this chapter our goal is to prove the following theorem:
Theorem 7.1.4. Suppose that at least one of $f_{1}, f_{2}$ is not of $C M$ type. If

$$
\limsup _{x \rightarrow \infty} \frac{\#\left\{p \leq x: a_{p}\left(f_{1}\right)=a_{p}\left(f_{2}\right)\right\}}{x / \log x}>0
$$

then $f_{1}=f_{2} \otimes \chi$ for some Dirichlet character $\chi$.
Using the celebrated work on the modularity of elliptic curves over $\mathbb{Q}$ due to Wiles [53], Breuil, Conrad, F. Diamond and Taylor [7], we get the result of Kulkarni, Patankar and Rajan [23] for elliptic curves over $\mathbb{Q}$.

### 7.2 Preliminaries

In this section, we collect the relevant facts that will be needed in various stages of our proof.

Proposition 7.2.1. If $f_{1}, f_{2}$ are normalized Hecke eigenforms, with at least one not of $C M$ type, such that $f_{1} \neq f_{2} \otimes \chi$ for some Dirichlet character $\chi$, then for any two positive integers $m, n$,

$$
\sum_{p \leq x} \frac{\sin (m+1) \theta_{p}^{(1)}}{\sin \theta_{p}^{(1)}} \frac{\sin (n+1) \theta_{p}^{(2)}}{\sin \theta_{p}^{(2)}}=o(x / \log x)
$$

as $x$ tends to infinity. Here the summation is over primes.
Proof. This is essentially Theorem 2.4 of [17] combined with the standard Tauberian theorem. However, for the sake of completeness, we give a proof in section 7.3.

Proposition 7.2.2. Let $0<\delta<\pi$. Let $f_{\delta}(x)$ be the "tent" function defined on $[-\pi, \pi]$ be given by

$$
f_{\delta}(x)= \begin{cases}1-|x| / \delta & \text { if }|x| \leq \delta \\ 0 & \text { if }|x|>\delta\end{cases}
$$

Then, for any $M \geq 1$, we have

$$
f_{\delta}(x)=\frac{\delta}{2 \pi}+2 \sum_{n=1}^{M} \frac{1-\cos n \delta}{\pi n^{2} \delta} \cos n x+O\left(\frac{1}{M \delta}\right)
$$

where the implied constant is absolute.
Proof. To prove the above proposition, let us compute the Fourier expansions of $f_{\delta}(x)$. The Fourier expansions of $f_{\delta}(x)$ is given by

$$
\frac{\delta}{2 \pi}+2 \sum_{n=1}^{\infty} \frac{1-\cos n \delta}{\pi n^{2} \delta} \cos n x
$$

For any $n \neq 0$,

$$
\hat{f}_{\delta}(n)=\frac{1}{2 \pi} \int_{-\delta}^{\delta}\left(1-\frac{|x|}{\delta} e^{-i n x}\right) d x
$$

$$
=\frac{1}{2 \pi} \int_{-\delta}^{0}\left(1+\frac{x}{\delta}\right) e^{-i n x} d x+\frac{1}{2 \pi} \int_{0}^{\delta}\left(1-\frac{x}{\delta}\right) e^{-i n x} d x
$$

Replacing $x$ by $-x$ in the first integral, the above equals

$$
\begin{align*}
\frac{1}{2 \pi} \int_{0}^{\delta}\left(1-\frac{x}{\delta}\right)(2 \cos n x) d x & =\frac{\sin n \delta}{\pi n}-\frac{1}{\pi \delta}\left(\delta \frac{\sin n \delta}{n}+\frac{\cos n \delta}{n^{2}}-\frac{\cos n 0}{n^{2}}\right)  \tag{7.2}\\
& =\left(\frac{1-\cos n \delta}{n^{2} \pi \delta}\right)
\end{align*}
$$

Now

$$
\sum_{n=-\infty}^{\infty}\left(\hat{f}(n) e^{i n x}+\hat{f}(-n) e^{-i n x}\right)=\frac{\delta}{2 \pi} \sum_{n=1}^{\infty}\left(\frac{1-\cos n \delta}{n^{2} \pi \delta} 2 \cos n x\right)
$$

Note that the above sum equals:

$$
\begin{equation*}
\sum_{n=1}^{\infty}\left(\frac{2 \cos n x}{n^{2} \pi \delta}-\frac{\cos n \delta 2 \cos n x}{n^{2} \pi \delta}\right)=\sum_{n=1}^{\infty} \frac{2 \cos n x}{n^{2} \pi \delta}-\sum_{n=1}^{\infty} \frac{\cos n \delta 2 \cos n x}{n^{2} \pi \delta} \tag{7.3}
\end{equation*}
$$

Let us estimate the truncation of the first sum on the right-hand side of (7.3) at $M$, namely

$$
\sum_{n=M+1}^{\infty} \frac{2 \cos n x}{n^{2} \pi \delta}
$$

Since, $\cos n x \leq 1$, the above sum is less than or equal to

$$
\sum_{n=M+1}^{\infty} \frac{2}{n^{2} \delta}
$$

We will use Euler Maclaurin summation formula (see [27, Theorem 2.1.9]) to estimate the above sum. Choosing $f(x)=\frac{1}{\delta x^{2}}$ in the Euler Maclaurin summation formula, we have

$$
\sum_{n=M+1}^{\infty} \frac{2}{n^{2} \pi \delta}=\int_{M+1}^{\infty} \frac{1}{\delta x^{2}}=\frac{2}{\delta(M+1)}
$$

Hence

$$
\sum_{n=M+1}^{\infty} \frac{2 \cos n x}{n^{2} \pi \delta}=\mathrm{O}\left(\frac{1}{M \delta}\right)
$$

Similarly, the second sum also will be $\mathrm{O}\left(\frac{1}{M \delta}\right)$ after truncating at $M$. So we can conclude that the tail terms will be $\mathrm{O}\left(\frac{1}{M \delta}\right)$ and we get the required result.

Proof of Theorem 7.1.4 Let us assume that $f_{1} \neq \chi \otimes f_{2}$ for any character $\chi$. Let $\pi>\delta>0$ and take $f_{\delta}(x)$ as in Proposition 7.2.2. Note that if $\theta_{p}^{(1)}=\theta_{p}^{(2)}$, then

$$
f_{\delta}\left(\theta_{p}^{(1)}-\theta_{p}^{(2)}\right)=1
$$

Hence,

$$
\#\left\{p \leq x: \theta_{p}^{(1)}=\theta_{p}^{(2)}\right\} \leq \sum_{p \leq x} f_{\delta}\left(\theta_{p}^{(1)}-\theta_{p}^{(2)}\right)+f_{\delta}\left(\theta_{p}^{(1)}+\theta_{p}^{(2)}\right) .
$$

By Proposition 7.2.2, the right hand side is equal to

$$
\frac{\delta \pi(x)}{\pi}+4 \sum_{n=1}^{M} \frac{1-\cos n \delta}{\pi n^{2} \delta} \sum_{p \leq x} \cos n \theta_{p}^{(1)} \cos n \theta_{p}^{(2)}+O\left(\frac{\pi(x)}{M \delta}\right)
$$

upon using the trigonometric identity

$$
\cos (A+B)+\cos (A-B)=2 \cos A \cos B
$$

The inner sum corresponding to $n=1$ is

$$
\sum_{p \leq x} \cos \theta_{p}^{(1)} \cos \theta_{p}^{(2)}
$$

In [26, Lemma 5], Murty has showed that if $f_{i}, i=1,2$ is defined as in (7.1) and $N_{1}=N_{2}=1$, then there there exists a $v$ such that $h=x^{\theta}, v<\theta<1$,

$$
\sum_{x \leq p \leq x+h} a_{p}\left(f_{1}\right) a_{p}\left(f_{2}\right)=\mathrm{o}(h) .
$$

In particular, writing $a_{p}\left(f_{i}\right)=2 \cos \theta_{p}^{(i)}$, we have

$$
\sum_{p \leq x} \cos \theta_{p}^{(1)} \cos \theta_{p}^{(2)}=\mathrm{o}(x)
$$

To treat $n \geq 2$, we use the identity

$$
\begin{equation*}
2 \cos n \theta=\frac{\sin (n+1) \theta}{\sin \theta}-\frac{\sin (n-1) \theta}{\sin \theta} \tag{7.4}
\end{equation*}
$$

so that we can rewrite our sum as

$$
\begin{gathered}
\sum_{n=2}^{M} \frac{1-\cos n \delta}{\pi n^{2} \delta} \\
\times \sum_{p \leq x}\left(\frac{\sin (n+1) \theta_{p}^{(1)}}{\sin \theta_{p}^{(1)}}-\frac{\sin (n-1) \theta_{p}^{(1)}}{\sin \theta_{p}^{(1)}}\right)\left(\frac{\sin (n+1) \theta_{p}^{(2)}}{\sin \theta_{p}^{(2)}}-\frac{\sin (n-1) \theta_{p}^{(2)}}{\sin \theta_{p}^{(2)}}\right) .
\end{gathered}
$$

Dividing by $\pi(x)$ and taking lim sup as $x$ tends to infinity, we obtain upon applying Proposition 7.2.1, that the inner sums go to zero. Thus, we obtain

$$
\limsup _{x \rightarrow \infty} \frac{\#\left\{p \leq x: \theta_{p}^{(1)}=\theta_{p}^{(2)}\right\}}{\pi(x)} \leq \frac{\delta}{\pi}+O\left(\frac{1}{\delta M}\right)
$$

Letting $M$ tend to infinity, we see that this density can be made arbitrarily small since $\delta$ is arbitrary. This contradicts our hypothesis and completes the proof.

### 7.3 Joint Sato-Tate distribution for two Hecke eigenforms

For the sake of completeness, let us review the Sato-Tate conjecture. For detail see Section 2.8.
Conjecture: If $f(p)$ is a $p^{t h}$ normalized Hecke eigenvalue then the family
$\{f(p)\}$ is equidistributed in $[-2,2]$ as $p \rightarrow \infty$ with respect to the Sato-Tate measure

$$
d \mu_{\infty}=\frac{1}{2 \pi} \sqrt{4-x^{2}} d x
$$

There are already several readable expositions of the proof of the Sato-Tate conjecture deduced from the potential automorphy of symmetric power $L$ functions (i.e the symmetric power $L$-function will be automorphic after a suitable change of the base field) (see for example, [17] and section 6 of Chapter 12 of [29]). What has not been explicitly presented in the literature is that the joint Sato-Tate distribution holds for two Hecke eigenforms, provided that one is not the Dirichlet twist of the other. In this section we prove the following theorem:

Theorem 7.3.1. For any rectangle $I \subset[-2,2]^{2}$,
$\lim _{x \rightarrow \infty} \frac{1}{\pi(x)} \#\left\{p \leq x:\left(f_{1}(p), f_{2}(p)\right) \in I\right.$ where $f_{1}(p)$ is not a character multiple of $\left.f_{2}(p)\right\}$

$$
=\int_{I} d \mu \times d \mu
$$

where $d \mu$ is the Sato-Tate measure.
Assume that $f_{1}(p)$ is not a character multiple of $f_{2}(p)$ and $f_{i}(p)=2 \cos \theta_{i}(p)$, $\theta_{i}(p) \in[0, \pi]$. By a change of variable, the above theorem is equivalent to the following:
For any rectangle $I \subset[0, \pi]^{2}$,

$$
\begin{aligned}
\lim _{x \rightarrow \infty} \frac{1}{\pi(x)} \# & \left\{p \leq x:\left(\theta_{1}(p), \theta_{2}(p)\right) \in I\right\} \\
& =\int_{I} d \nu \times d \nu
\end{aligned}
$$

where $\nu=\frac{2}{\pi} \sin ^{2} \theta$. Note that if $\theta_{i}(p)$ is in the family then $-\theta_{i}(p)$ is in the family. Now to prove the above theorem, we will compute the Weyl limits
of the family $\left\{\left( \pm \theta_{1}(p), \pm \theta_{2}(p)\right)\right\}$ and that is for any two positive integer $m$ and $n$,

$$
\lim _{x \rightarrow \infty} \frac{1}{\pi(x)} \sum_{p \leq x} e\left( \pm m \theta_{1}(p) \pm n \theta_{2}(p)\right)=\lim _{x \rightarrow \infty} \frac{1}{\pi(x)} \sum_{p \leq x} 2 \cos m \theta_{1}(p) 2 \cos n \theta_{2}(p)
$$

Using (7.4), the above equals

$$
=\sum_{p \leq x}\left(\frac{\sin (m+1) \theta_{p}^{(1)}}{\sin \theta_{p}^{(1)}}-\frac{\sin (m-1) \theta_{p}^{(1)}}{\sin \theta_{p}^{(1)}}\right)\left(\frac{\sin (n+1) \theta_{p}^{(2)}}{\sin \theta_{p}^{(2)}}-\frac{\sin (n-1) \theta_{p}^{(2)}}{\sin \theta_{p}^{(2)}}\right) .
$$

From the above discussion, it is clear that to prove Theorem 7.3.1 it is sufficient to prove Proposition 7.2.1. To do so, according to our need, let us review some basic representation theory. Let $G$ be a finite group and $V$ be a complex vector space. For details, the reader may consult [8] and [41].

Definition 7.3.2. A group homomorphism $\pi$ of $G$ to $G L(V)$ is said to be a representation $(\pi, V)$ of $G$, where $G L(V)$ is the group of isomorphisms of $V$ on to itself.

Definition 7.3.3. A character $\chi$ associated to a representation $(\pi, V)$ of $G$ is defined by

$$
\chi_{\pi}(g)=\operatorname{Tr}_{V}(\pi(g))
$$

Let $W$ be a subspace of $V$. The subspace $W$ is said to be $G$-invariant if

$$
\pi(g) W \subset W \text { for all } g \in G
$$

Definition 7.3.4. A representation of $G$ on $W, \pi_{W}: G \rightarrow G L(W)$ is said to be a subrepresentation of $(\pi, V)$ if

$$
\pi_{W}(g)=\pi(g) \mid W \text { for all } g \in G
$$

Definition 7.3.5. Any representation $(\pi, V)$ of $G$ is said to be an irreducible representation if the only $G$-invariant subspaces of $G$ are $\{0\}$ and $V$.

Definition 7.3.6. Any character of an irreducible representation is called irreducible character.

A linear combination of characters is known as a generalised character. A group $G$ is said to be an elementary group if it is a direct product of cyclic groups and groups whose order is a power of a prime.

Definition 7.3.7. Let $G$ be a topological group and $V$ be an Hilbert space. A unitary representation of $G$ is an isometric action of $G$ on $V$ so that the action map $G \times V \rightarrow V$ is continuous.

To define automorphic representations, we will closely follow [8, page 92] Let $F$ be a number field and $\mathbb{A}=\mathbb{A}_{F}$ be the adele ring of $F$. Let $Z_{\mathbb{A}}$ be the group of scalar matrices with entries in the group $\mathbb{A}^{*}$, where $\mathbb{A}^{*}$ be the set of units of $\mathbb{A}$. The group $G L_{n}(F)$ is a discrete subgroup of $G L_{n}(\mathbb{A})$ and the quotient $G L_{n}(F) / G L_{n}(\mathbb{A})$, and the quotient $Z_{\mathbb{A}} G L_{n}(F) / G L_{n}(\mathbb{A})$ has finite volume. Let $w$ be a unitary character of $\mathbb{A}^{*} / F^{*}$ i.e $|w(z)|=1$ for all $z \in$ $\mathbb{A}^{*}$. Let $L^{2}\left(G L_{n}(F) / G L_{n}(\mathbb{A}), w\right)$ be the space of measurable functions $f$ on $G L_{n}(\mathbb{A})$ such that

$$
f\left(\left(\begin{array}{llll}
z & & & \\
& \cdot & & \\
& \cdot & & \\
& & & \\
& & & z
\end{array}\right) g\right)=w(z) f(g) \text { for } z \in \mathbb{A}^{*}
$$

and such that

$$
\int_{Z_{\mathbb{A}} G L_{n}(F) / G L_{n}(\mathbb{A})}|f(g)|^{2} d g<\infty,
$$

where $d g$ is the Haar measure.
Definition 7.3.8. A representation $(\pi, V)$ is said to be an automorphic representation if $\pi$ is an irreducible unitary representation of $G L_{n}(\mathbb{A})$ that occurs in $L^{2}\left(G L_{n}(F) / G L_{n}(\mathbb{A}), w\right)$.

Theorem 7.3.9 (Brauer). Let $G$ be a finite group, and let $\chi$ be a generalized character. Then, there exist elementary subgroups $E_{1}, E_{2}, \ldots$ and irreducible characters $\psi_{1}, \psi_{2}, \ldots$ of the $E_{i}$ and integers $a_{i} \in \mathbb{Z}$ such that

$$
\chi=\sum a_{i} \chi_{i}^{G}
$$

Let $(\pi, V)$ and $\left(\sigma, V^{\prime}\right)$ be two representations of the group $G$. Let us denote the set of $G$-invariant maps from $V$ to $V^{\prime}$ as $\operatorname{Hom}_{G}(\pi, \sigma)$.

Theorem 7.3.10 (Frobenius reciprocity law). Let $H$ be a subgroup of a finite group $G$ and $\sigma$ be a representation of $H$. Let $\pi$ be a representation of $G$. Then

$$
\operatorname{Hom}_{G}\left(\pi, \operatorname{Ind}_{H}^{G}(\sigma)\right) \equiv \operatorname{Hom}_{H}\left(\left.\pi\right|_{H}, \sigma\right),
$$

and

$$
\operatorname{Hom}_{G}\left(\operatorname{Ind}_{H}^{G}(\sigma), \pi\right) \equiv \operatorname{Hom}_{H}\left(\sigma,\left.\pi\right|_{H}\right),
$$

where $\operatorname{Ind}_{H}^{G}$ is the induced representation of $H$ on $G$.
The following theorem can be found in [27] or [49]:
Theorem 7.3.11. [Tauberian theorem] Let $F(s)=\sum_{n=0}^{\infty} \frac{a_{n}}{n^{s}}$ be a Dirichlet series with non negative real coefficients $a_{n}$ such that
(a) $F(s)$ is absolutely convergent for $\operatorname{Re}(s)>1$.
(b) $F(s)$ can be extends meromorphically to the region $\operatorname{Re}(s) \geq 1$ with a possible pole at $s=1$ with residue $R$.
(c) $\sum_{n \leq x} a_{n}=\mathrm{O}(x)$. Then

$$
\sum_{n \leq x} a_{n}=R x+\mathrm{o}(x)
$$

we now need to understand the $\ell$-adic Galois representation attached to a cusp form of level $N$ and weight $k$. Let $\ell$ be a prime. $K$ be a finite extension of $\mathbb{Q}_{l}$. Let $V$ be finite-dimensional $K$ vector space.

Definition 7.3.12. A continuous representation $\rho: \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q}) \rightarrow G L_{K}(V)$ is said to be a $\ell$-adic representation of $\operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$.

Definition 7.3.13. The normal subgroup $I_{p}=\overline{\mathbb{Q}} / \mathbb{Q}_{p}\left(\zeta_{m}: p \nmid m\right)$ defined by any intermediate extensions of $\mathbb{Q}_{p} \subset \mathbb{Q}_{p}\left(\zeta_{m}: p \nmid m\right) \subset \overline{\mathbb{Q}}$ is called the inertia subgroup.

Definition 7.3.14. An $\ell$-adic representation is said to be unramified at $p$ if the restriction to the inertia group is trivial.

Definition 7.3.15. Let $f(z)=\sum_{n} a_{n} e^{2 \pi i n z}$ be a normalized Hecke eigenform of level $N$ and weight $k$. A 2-dimensional $\ell$-adic representation $V$ over $K$ is said to be associated to $f$ if, for every $p \nmid N \ell, V$ is unramified at $p$ and

$$
\operatorname{Tr}\left(\phi_{p}: V\right)=a_{p}(f),
$$

where $\phi_{p} \in \operatorname{Gal}(\overline{\mathbb{Q}} / \mathbb{Q})$,

$$
\phi_{p}(a)=a^{p} .
$$

Let $f=\sum_{n} \frac{a_{n}}{n^{s}}$ is an Hecke eigenform. Then by the famous work of Deligne [14], for any prime $p$,

$$
a(p)=\alpha(p)+\beta(p) \text { and }|\alpha(p)|=|\beta(p)|=|\alpha(p) \beta(p)|=1
$$

Now we will define the corresponding Symmetric $m$-th power ( $s y m^{m}$ ) $L$ function attached to $f$.

Definition 7.3.16. For any positive integer $m$, the $m$-th Symmetric power $L$ function attached to $f$ is defined as

$$
L\left(s, S^{m} m^{m} f\right):=\prod_{p} \sum_{i=0}^{m}\left(1-\frac{\alpha(p)^{i} \beta(p)^{m-i}}{p^{s}}\right)^{-1}
$$

for $\operatorname{Re}(s)>1$.
Definition 7.3.17. For two positive integers $m, n$, the Rankin-Selberg convolution of L-functions attached to $S y m^{m}$ and $S y m^{n}$ is defined as

$$
L\left(s, S_{y m}^{m} f \times \text { Sym }^{n} f\right):=\prod_{p} \prod_{i=0}^{m} \prod_{j=0}^{n}\left(1-\frac{\alpha(p)^{i} \beta(p)^{m-i} \alpha(p)^{j} \beta(p)^{n-j}}{p^{s}}\right)^{-1}
$$

for $\operatorname{Re}(s)>1$.
The following theorem can be found in [12, Proposition 5.1];
Theorem 7.3.18. [Cogdel-Michel] Let $f_{1}$ and $f_{2}$ be two holomorphic primitive
forms of even weight $k=2 l$ and trivial nebentypus, one of them, say $f_{1}$, not of CM type. Suppose that, for some given $m \geq 1$,

$$
L_{p}\left(s, S y m^{m} f_{1}\right)=L_{p}\left(s, S y m^{n} f_{2}\right)
$$

for every prime poutside a set of density 0 . Then there exists a character of order at most 2 such that $f_{2}=f_{1} \otimes \chi$. (If $l$ is odd, then $\chi$ is trivial.)

Remark 7.3.19. In case of non-trivial Nebentypus, Rajan (see [34, Corollary 5.1]) proved the above result.

In 1981, Shahidi [42] proved the following non-vanishing theorem for RankinSelberg $L$-function.

Theorem 7.3.20. [Shahidi] The Rankin-Selberg L-function $L\left(s, f_{1} \times f_{2}\right)$ does not vanish on $\mathbb{R}(s)=1$.

Proof of Proposition 7.3.1 Let $\rho_{1}$ and $\rho_{2}$ be the associated ( $\ell$-adic) Galois representations of $f_{1}, f_{2}$ respectively. Let $F / F_{1}$ be solvable. Then there is a chain of fields $F \supset F_{m} \supset F_{m-1} \supset \cdots \supset F_{1}$ such that $F / F_{m}$ and $F_{j} / F_{j-1}$ for $2 \leq j \leq m$ are all cyclic of prime degree. By [1], the automorphic induction map exists from $F$ to $F_{m}$ and successively from $F_{j}$ to $F_{j-1}$ for $2 \leq j \leq m$ which at the final stage is $F_{1}$. Here most important thing for us is that both the Galois and automorphic representations obtained by descent are Galois invariant at every step (see for example, the comment at the bottom of [9, page 11]). By the work of [3], both $\operatorname{Sym}^{n}\left(\rho_{1}\right)$ and $\operatorname{Sym}^{m}\left(\rho_{2}\right)$ are potentially automorphic over a totally real Galois extension $F$ over $\mathbb{Q}$. By the Arthur-Clozel theory of base change [1], we see that for any sub field $F_{1}$ of $F$ with $F / F_{1}$ solvable, both $\left.\operatorname{Sym}^{m}\left(\rho_{1}\right)\right|_{F_{1}}$ and $\left.\operatorname{Sym}^{n}\left(\rho_{2}\right)\right|_{F_{1}}$ are also automorphic over $F_{1}$. Let $G=\operatorname{Gal}(F / \mathbb{Q})$. By Brauer induction, we can write

$$
1=\sum_{i} a_{i} \operatorname{Ind}_{H_{i}}^{G} \psi_{i},
$$

where the $a_{i}$ 's are integers and $\psi_{i}$ 's are one-dimensional characters of nilpotent subgroups $H_{i}$ of $G$. Thus,
$L\left(s,\left(\operatorname{Sym}^{m}\left(\rho_{1}\right) \otimes \operatorname{Sym}^{n}\left(\rho_{2}\right)\right) \otimes 1\right)=\prod_{i} L\left(s,\left(\operatorname{Sym}^{m} \rho_{1} \otimes \operatorname{Sym}^{n} \rho_{2}\right) \otimes \operatorname{Ind}_{H_{i}}^{G} \psi_{i}\right)^{a_{i}}$.
By Frobenius reciprocity,

$$
\left(\operatorname{Sym}^{m} \rho_{1} \otimes \operatorname{Sym}^{n} \rho_{2}\right) \otimes \operatorname{Ind}_{H_{i}}^{G} \psi_{i}=\operatorname{Ind}_{H_{i}}^{G}\left(\left.\left(\left(\operatorname{Sym}^{m} \rho_{1}\right) \otimes\left(\operatorname{Sym}^{n} \rho_{2}\right)\right)\right|_{L^{H_{i}}} \otimes \psi_{i}\right) G a l(\overline{\mathbb{Q}} / \mathbb{Q}) .
$$

Since $\left.\left(\operatorname{Sym}^{m} \rho_{1}\right)\right|_{L^{H_{i}}}$ and $\left.\left(\operatorname{Sym}^{n} \rho_{2}\right)\right|_{L^{H_{i}}}$ are both automorphic over $L^{H_{i}}$, and $\psi_{i}$ is a Hecke character of $L^{H_{i}}$ by Artin reciprocity, we can form the RankinSelberg convolution:

$$
\begin{equation*}
L\left(s,\left.\left.\left(\operatorname{Sym}^{m}\left(\rho_{1}\right)\right)\right|_{L^{H_{i}}} \otimes\left(\operatorname{Sym}^{n}\left(\rho_{2}\right)\right)\right|_{L^{H_{i}}} \otimes \psi_{i}\right) \tag{7.5}
\end{equation*}
$$

By theorems of Cogdel-Michel (see Theorem 7.3.18) and Rajan (see Remark 7.3.19) which says that the $m$-th symmetric powers of the Galois representations associated with $f_{1}$ and $f_{2}$ are equal then $f_{1}$ is a Dirichlet character twist of $f_{2}$. Thus, if $f_{1} \neq \chi \otimes f_{2}$ for any Dirichlet character $\chi$, we have

$$
\pi_{1}:=\left.\left(\operatorname{Sym}^{m}\left(\rho_{1}\right)\right)\right|_{F^{H_{i}}} \text { and } \pi_{2}:=\left.\left(\operatorname{Sym}^{n}\left(\rho_{2}\right)\right)\right|_{F^{H_{i}}} \otimes \psi_{i}
$$

are such that $\pi_{2} \nexists \pi \otimes|d e t|^{i t}$ for any real number $t$.
By our hypothesis, $f_{1}$ and $f_{2}$ are analytic and non-vanishing for $\operatorname{Re}(s) \geq 1$. By standard Rankin-Selberg theory (see [11, page 69 or 225]) and Shahidi's result (see Theorem 7.3.20), the $L$-function (7.5) is analytic and non-vanishing in the region $\operatorname{Re}(s) \geq 1$.
Thus,

$$
L\left(s,\left(\operatorname{Sym}^{m} \rho_{1}\right) \otimes \operatorname{Sym}^{n} \rho_{2}\right)
$$

extends to an analytic function to $\operatorname{Re}(s) \geq 1$ and is non-vanishing there. By the standard Tauberian theorem (see Theorem 7.3.11) applied to the logarithmic derivative of this $L$-function, we deduce Proposition 7.2.1.

Remark 7.3.21. Our argument extends easily to imply a corresponding result for any two modular forms $f_{1}$ and $f_{2}$ over a totally real field since the results of [3] apply in this context also.

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