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On Atkin–Swinnerton-Dyer congruence relations

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Abstract

In this paper, we exhibit a noncongruence subgroup \( \Gamma \) whose space of weight 3 cusp forms \( S_3(\Gamma) \) admits a basis satisfying the Atkin–Swinnerton-Dyer congruence relations with two weight 3 newforms for certain congruence subgroups. This gives a modularity interpretation of the motive attached to \( S_3(\Gamma) \) by Scholl and also verifies the Atkin–Swinnerton-Dyer congruence conjecture for this space.

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1. Introduction

The theory of modular forms for congruence subgroups is well developed. Given a cuspidal normalized newform \( g = \sum_{n \geq 1} a_n(g)q^n \), where \( q = e^{2\pi i \tau} \), of weight \( k \geq 2 \) level \( N \) and character \( \chi \), the Fourier coefficients of \( g \) satisfy the recursive relation

\[
a_{np}(g) - a_p(g)a_n(g) + \chi(p)p^{k-1}a_{n/p}(g) = 0
\]

\( (1) \)

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for all primes $p$ not dividing $N$ and for all $n \geq 1$. Thanks to the work of Eichler [Eic57], Shimura [Shi59], and Deligne [Del73], there exists a compatible family of $\lambda$-adic representations $\rho_{\lambda,g}$ of the Galois group $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$, unramified outside $lN$, where $\lambda$ divides $l$, such that

$$
\text{Tr}(\rho_{\lambda,g}(\text{Frob}_p)) = a_p(g),
\det(\rho_{\lambda,g}(\text{Frob}_p)) = \chi(p)p^{k-1}
$$

for all primes $p$ not dividing $lN$. Combining both, we see that the characteristic polynomial $H_p(T) = T^2 - A_1(p)T + A_2(p)$ of $\rho_{\lambda,g}(\text{Frob}_p)$ is independent of the $\lambda$’s not dividing $p$, and the Fourier coefficients of $g$ satisfy the relation

$$a_{np}(g) - A_1(p)a_n(g) + A_2(p)a_{n/p}(g) = 0 \quad (2)$$

for all $n \geq 1$ and all primes $p$ not dividing $N$.

The knowledge on modular forms for noncongruence subgroups, however, is far from satisfactory. For example, the paper [Sch97] by Scholl and Serre’s letter to Thompson [Th89] indicate that the obvious definitions of Hecke operators for noncongruence subgroups would not work well. On the positive side, Atkin and Swinnerton-Dyer [ASD71] initiated the study of the arithmetic properties of modular forms for noncongruence subgroups, and they have observed very interesting congruence relations for such forms, which we now explain.

Let $\Gamma$ be a noncongruence subgroup of $\text{SL}_2(\mathbb{Z})$ with finite index. For an integer $k \geq 2$, denote by $S_k(\Gamma)$ the space of cusp forms of weight $k$ for $\Gamma$ and by $d$ its dimension. The modular curve $X_\Gamma$, the compactification by adding cusps of the quotient of the Poincaré upper half plane by $\Gamma$, has a model defined over a number field $K$ in the sense of Scholl [Sch85, Section 5].

As explained in [ASD71,Sch85,Sch87], there exists a subfield $L$ of $K$, an element $\kappa \in K$ with $\kappa^\mu \in L$, where $\mu$ is the width of the cusp $\infty$, and a positive integer $M$ such that $\kappa^\mu$ is integral outside $M$ and $S_k(\Gamma)$ has a basis consisting of $M$-integral forms. Here a form $f$ of $\Gamma$ is called $M$-integral if in its Fourier expansion at the cusp $\infty$

$$f(\tau) = \sum_{n \geq 1} a_n(f)q^{n/\mu}, \quad (3)$$

the Fourier coefficients $a_n(f)$ can be written as $\kappa^\mu c_n(f)$ with $c_n(f)$ lying in the ring $\mathcal{O}_L[1/M]$, where $\mathcal{O}_L$ denotes the ring of integers of $L$.

Based on their numerical data, Atkin and Swinnerton-Dyer [ASD71] made an amazing discovery of congruence relations for certain cusp forms for noncongruence subgroups. It is tempting to extrapolate from their observations and from our own numerical data to formulate the following congruence conjecture.

**Conjecture 1.1** (Atkin–Swinnerton-Dyer congruences). Suppose that the modular curve $X_\Gamma$ has a model over $\mathbb{Q}$ in the sense of [Sch85, Section 5]. There exist a positive
integer $M$ and a basis of $S_k(\Gamma)$ consisting of $M$-integral forms $f_j$, $1 \leq j \leq d$, such that for each prime $p$ not dividing $M$, there exists a nonsingular $d \times d$ matrix $(\lambda_{i,j})$ whose entries are in a finite extension of $\mathbb{Q}_p$, algebraic integers $A_p(j)$, $1 \leq j \leq d$, with $|\sigma(A_p(j))| \leq 2p^{(k-1)/2}$ for all embeddings $\sigma$, and characters $\chi_j$ unramified outside $M$ so that for each $j$ the Fourier coefficients of $h_j := \sum_i \lambda_{i,j} f_i$ satisfy the congruence relation

$$\text{ord}_p(a_{np}(h_j) - A_p(j)a_n(h_j) + \chi_j(p)p^{k-1}a_{n/p}(h_j)) \geq (k-1)(1 + \text{ord}_p n)$$

for all $n \geq 1$; or equivalently, for all $n \geq 1$,

$$a_{np}(h_j) - A_p(j)a_n(h_j) + \chi_j(p)p^{k-1}a_{n/p}(h_j)/(np)^{k-1}$$

is integral at all places dividing $p$.

In other words, the recursive relation (1) on Fourier coefficients of modular forms for congruence subgroups is replaced by the congruence relation (4) for forms of noncongruence subgroups. The meaning of $A_p(j)$’s is mysterious; the examples in [ASD71] suggest that they satisfy the Sato-Tate conjecture.

In [Sch85] Scholl proved a “collective version” of this conjecture.

**Theorem 1.2 (Scholl).** Suppose that $X_\Gamma$ has a model over $\mathbb{Q}$ as before. Attached to $S_k(\Gamma)$ is a compatible family of $2d$-dimensional $l$-adic representations $\rho_l$ of the Galois group $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ unramified outside $lM$ such that for primes $p > k + 1$ not dividing $lM$, the following hold.

(i) The characteristic polynomial

$$H_p(T) = \sum_{0 \leq r \leq 2d} B_r(p)T^{2d-r}$$

of $\rho_l(\text{Frob}_p)$ lies in $\mathbb{Z}[T]$ and is independent of $l$, and its roots are algebraic integers with absolute value $p^{(k-1)/2}$;

(ii) For any $M$-integral form $f$ in $S_k(\Gamma)$, its Fourier coefficients $a_n(f)$, $n \geq 1$, satisfy the congruence relation

$$\text{ord}_p(a_{np^d}(f) + B_1(p)a_{np^{d-1}}(f) + \cdots + B_{2d-1}(p)a_{n/p^{d-1}}(f) + B_{2d}(p)a_{n/p^d}(f))$$

$$\geq (k-1)(1 + \text{ord}_p n)$$

for $n \geq 1$.

Scholl’s theorem establishes the Atkin–Swinnerton-Dyer congruences if $S_k(\Gamma)$ has dimension 1. If the Atkin–Swinnerton-Dyer congruences were established in general,
then

\[ H_p(T) = \prod_{1 \leq j \leq d} \left( T^2 - A_p(j)T + \chi_j(p)p^{k-1} \right). \]

Scholl’s congruence relation (6) may be regarded as a collective replacement for forms for noncongruence subgroup of equality (2) for newforms.

Let \( f = \sum_{n \geq 1} a_n(f)q^n \) be an \( M \)-integral cusp form in \( S_k(\Gamma) \), and let \( g = \sum_{n \geq 1} b_n(g)q^n \) be a normalized newform of weight \( k \) level \( N \) and character \( \chi \).

**Definition 1.3.** The two forms \( f \) and \( g \) above are said to satisfy the Atkin–Swinnerton-Dyer congruence relations if, for all primes \( p \) not dividing \( MN \) and for all \( n \geq 1 \),

\[
(a_n p(f) - b_p(g)a_n(f) + \chi(p)p^{k-1}a_n/p(f))/(np)^{k-1}
\]

is integral at all places dividing \( p \).

In particular, if \( S_3(\Gamma) \) has a basis of \( M \)-integral forms \( f_j \), \( 1 \leq j \leq d \), such that each \( f_j \) satisfies the Atkin–Swinnerton-Dyer congruence relations with some cuspidal newform \( g_j \) of weight 3 for certain congruence subgroup, then this not only establishes the Atkin–Swinnerton-Dyer congruences conjecture for the space \( S_3(\Gamma) \), but also provides an interpretation of the \( A_p(j) \)'s in the conjecture. Geometrically, this means that the motive attached to \( S_3(\Gamma) \) by Scholl comes from modular forms for congruence subgroups. Furthermore, if \( -I \notin \Gamma \) and \( \mathcal{E}_\Gamma \) is the elliptic modular surface associated to \( \Gamma \) in the sense of [Shi72], then \( \mathcal{E}_\Gamma \) is an elliptic surface with base curve \( X_\Gamma \). The product of the \( L \)-functions \( \prod_{1 \leq j \leq d} L(s, g_j) \) occurs in the Hasse–Weil \( L \)-function \( L(s, \mathcal{E}_\Gamma) \) attached to the surface \( \mathcal{E}_\Gamma \), and it is the part arising from the transcendental lattice of the surface. In this case, \( L(s, \mathcal{E}_\Gamma) \) has both its numerator and denominator product of automorphic \( L \)-functions. In other words, the Hasse–Weil \( L \)-function of \( \mathcal{E}_\Gamma \) “is modular”. To-date, only a few such examples are known. In a recent preprint of Livné and Yui [LY03], the \( L \)-functions of some rank 4 motives associated to nonrigid Calabi–Yau threefolds are proven to be the \( L \)-functions of some automorphic forms.

For the noncongruence subgroup \( \Gamma_{7,1,1} \) studied in [ASD71], the space \( S_4(\Gamma_{7,1,1}) \) is 1-dimensional. Let \( f \) be a nonzero 14-integral form in \( S_4(\Gamma_{7,1,1}) \). Scholl proved in [Sch88] that there is a normalized newform \( g \) of weight 4 level 14 and trivial character such that \( f \) and \( g \) satisfy the Atkin–Swinnerton-Dyer congruence relations. In the unpublished paper [Sch93], Scholl obtained a similar result for \( S_4(\Gamma_{4,3}) \) and \( S_4(\Gamma_{5,2}) \); both spaces are also 1-dimensional.

The purpose of this paper is to present an example of 2-dimensional space of cusp forms of weight 3 whose associated \( l \)-adic representation is modular and the existence of an \( M \)-integral basis, independent of \( p \), such that each satisfies the Atkin–Swinnerton-Dyer congruence relations with a cusp form of a congruence subgroup. More precisely, we shall prove

**Theorem 1.4.** Let \( \Gamma \) be the index 3 noncongruence subgroup of \( \Gamma^1(5) \) such that the widths at two cusps \( \infty \) and \(-2\) are 15.
(1) Then $X_\Gamma$ has a model over $\mathbb{Q}$, $\kappa = 1$, and the space $S_3(\Gamma)$ is 2-dimensional with a basis consisting of 3-integral forms

$$f_+(\tau) = q^{1/15} + iq^{2/15} - \frac{11}{3} q^{4/15} - \frac{4}{9} q^{7/15} + i \frac{71}{9} q^{8/15} + \frac{932}{81} q^{10/15} + O(q^{11/15}),$$

$$f_-(\tau) = q^{1/15} - iq^{2/15} - \frac{11}{3} q^{4/15} - \frac{4}{9} q^{7/15} - i \frac{71}{9} q^{8/15} + \frac{932}{81} q^{10/15} + O(q^{11/15}).$$

(2) The 4-dimensional $l$-adic representation $\rho_l$ of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ associated to $S_3(\Gamma)$ constructed by Scholl is modular. More precisely, there are two cuspidal newforms of weight 3 level 27 and character $\chi_{-3}$ given by

$$g_+(\tau) = q - 3iq^2 - 5q^4 + 3iq^5 + 5q^7 + 3iq^8 + 9q^{10} + 15iq^{11} - 10q^{13} - 15iq^{14} - 11q^{16} - 18iq^{17} - 16q^{19} - 15iq^{20} + 45q^{22} + 12i q^{23} + O(q^{24}),$$

$$g_-(\tau) = q + 3iq^2 - 5q^4 - 3iq^5 + 5q^7 - 3iq^8 + 9q^{10} - 15iq^{11} - 10q^{13} + 15iq^{14} - 11q^{16} + 18iq^{17} - 16q^{19} + 15iq^{20} + 45q^{22} - 12iq^{23} + O(q^{24}),$$

such that over the extension by joining $\sqrt{-1}$, $\rho_l$ decomposes into the direct sum of the two $\lambda$-adic representations attached to $g_+$ and $g_-$, where $\lambda$ is a place of $\mathbb{Q}(i)$ dividing $l$.

(3) $f_+$ and $g_+$ (resp. $f_-$ and $g_-$) satisfy the Atkin–Swinnerton-Dyer congruence relations.

Here $\chi_{-3}$ is the quadratic character attached to the field $\mathbb{Q}(\sqrt{-3})$. The precise definition of $\Gamma$ in terms of generators and relations is given at the end of Section 3.

The proof of this theorem occupies Sections 2–7. Here we give a sketch. The modular curve $X_\Gamma$ of $\Gamma$ is a three-fold cover of the congruence modular curve $X_{\Gamma^1(5)}$ ramified only at two cusps of $\Gamma^1(5)$. By explicitly computing the Eisenstein series of weight 3 for $\Gamma^1(5)$, we obtain in Section 4 an explicit basis $f_+$ and $f_-$ of $S_3(\Gamma)$ which are 3-integral, as stated above.

To establish the congruence relations, we take advantage of the existence of an elliptic surface $\mathcal{E}$ over $X_\Gamma$ with an explicit algebraic model. There exists a $\mathbb{Q}$-rational involution $A$ on $X_\Gamma$ which induces an action on $\mathcal{E}$ of order 4, which commutes with the action of the Galois group over $\mathbb{Q}$. In fact, $f_+$ and $f_-$ are eigenfunctions of $A$ with eigenvalues $-i$ and $i$, respectively. The explicit defining equation of $\mathcal{E}$ gives rise to a 4-dimensional $l$-adic representation $\rho_l^\ast$ of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ which is isomorphic to the $l$-adic representation $\rho_l$ Scholl constructed in [Sch85] at most up to a quadratic twist $\phi$. Take $l = 2$. The dyadic representations are unramified outside 2 and 3. Making use of the action of $A$, we may regard $\rho_2^\ast$ as a 2-dimensional representation over $\mathbb{Q}_2(i)$, which is...
isomorphic to the completion of $\mathbb{Q}(i)$ at the place with $1+i$ as a uniformizer, denoted by $\mathbb{Q}(i)_{1+i}$ for convenience. The explicit defining equation allows us to determine the characteristic polynomial of the Frobenius under $\rho_2^*$ over $\mathbb{Q}_2$ for small primes, and that over $\mathbb{Q}(i)_{1+i}$ except for primes congruent to $2$ mod $3$, in which case the trace is determined up to sign (cf. Table 1 in Section 5).

On the other hand, the two cuspidal newforms $g_+$ and $g_-$ combined come from a 4-dimensional 2-adic representation $\tilde{\rho}_2$ of $\text{Gal}(\bar{\mathbb{Q}}/\mathbb{Q})$, on whose space the Atkin–Lehner operator $H_{27}$ acts. It has order 4. So we may also regard $\tilde{\rho}_2$ as a 2-dimensional representation over $\mathbb{Q}(i)_{1+i}$ except for primes congruent to 2 mod 3, in which case the trace is determined up to sign (cf. Table 1 in Section 5).

To prove that $\rho_2^*$ and $\tilde{\rho}_2$ are isomorphic, we show that $\rho_2^*$ satisfies the congruence relations (6) for $n=1$, $p=7$ and $p=13$. This in turn forces $\phi$ to be either $\chi_{-3}$ or trivial. In either case, we have the desired isomorphism. Finally, to obtain the Atkin–Swinnerton-Dyer congruence relations between $f_\pm$ and $g_\pm$, we compare the $p$-adic theory on eigenspaces of $A$ and the dyadic theory on eigenspaces of $H_{27}$, and draw the desired conclusion using Scholl’s proof of Theorem 1.2 in Section 5 of [Sch85].

We end the paper by observing that if the space of cusp forms of weight 3 for a noncongruence subgroup $\Gamma'$ is 1-dimensional with a nonzero $M$-integral form $f$ and there is an elliptic $K3$ surface over the modular curve $X_{\Gamma'}$, then there is a cuspidal newform $g$ such that $f$ and $g$ satisfy the Atkin–Swinnerton-Dyer congruence relations.

2. An elliptic surface

Let $E$ denote the minimal smooth model of the elliptic surface given by

$$y^2 + (1 - t^3)xy - t^3y = x^3 - t^3x^2,$$

where the parameter $t$ runs through the points in the complex projective line $\mathbb{C}P^1$. Viewed as an elliptic curve defined over $\mathbb{C}(t)$, its $j$-invariant is

$$j = \frac{(t^{12} - 12t^9 + 14t^6 + 12t^3 + 1)^3}{t^{15}(t^6 - 11t^3 - 1)}.$$  

Its Mordell–Weil group is isomorphic to $\mathbb{Z}/5\mathbb{Z}$. Indeed, it is a subgroup of $\mathbb{Z}/15\mathbb{Z} \times \mathbb{Z}/15\mathbb{Z}$ and it contains $\mathbb{Z}/5\mathbb{Z}$ as a subgroup. By examining the restriction of the determinant of its transcendental lattice, we find that this group has order a power of 5. We conclude that the group is $\mathbb{Z}/5\mathbb{Z}$ from a result of Cox and Parry [CP80].

It is clear that at a generic $t \in \mathbb{C}$, the fiber of the natural projection

$$\pi : E \to \mathbb{C}P^1,$$

$$\pi \quad \begin{cases} (x, y, t) \mapsto t \end{cases}$$


is an elliptic curve, namely, a smooth compact curve of genus one. This is the case except for 8 values of \( t \): 0, \( \infty \), and the six roots of \( t^6 - 11t^3 - 1 = 0 \). At these 8 exceptional values of \( t \), there are 8 special fibers. They are identified, using Tate algorithm, to be of respective type \( I_{15} \), \( I_{15} \), \( I_1 \), \( I_1 \), \( I_1 \), \( I_1 \), \( I_1 \), \( I_1 \) in Kodaira’s notation.

The surface \( E \) is an elliptic modular surface by [Shi72], [Nor85]. Denote by \( \Gamma \) an associated modular group, which is a subgroup of \( SL_2(\mathbb{Z}) \) of finite index. Let \( \bar{\Gamma} = \pm \Gamma / \pm I \) be its projection in \( PSL_2(\mathbb{Z}) \). We know from the information of the special fibers that the group \( \bar{\Gamma} \) has no elliptic points. In other words, it is a torsion free subgroup of \( SL_2(\mathbb{Z}) \).

Now we consider some topological and geometrical invariants of \( E \). By Kodaira’s formula [Kod63], the Euler characteristic of \( E \) is 36. Its irregularity, which equals the genus of the base curve, is 0. Its geometric genus is 2 by Noether’s formula. Denote by \( h^{i,j} = \dim H^j(E, \mathcal{O}_E^i) \) the \((i, j)\)’s Hodge number of \( E \). Then the Hodge numbers of \( E \) can be arranged into the following Hodge diamond:

\[
\begin{array}{ccc}
1 & & 0 \\
& 0 & & 30 \\
2 & & 2 \\
& 0 & & 0 \\
1 & & & & 1
\end{array}
\]

where on the \( i \)th row, the \( j \)th entry is \( h^{i-j,j-1} \), when \( i = 1, 2, 3 \), and is \( h^{3-j,j+i-4} \) when \( i = 4, 5 \).

As the group \( \bar{\Gamma} \) is not among the genus-zero torsion free congruence subgroups of \( PSL_2(\mathbb{Z}) \) listed by Sebbar [Seb01], we conclude that \( \bar{\Gamma} \) is a noncongruence subgroup of \( SL_2(\mathbb{Z}) \).

Let \( L \) be the free part of the cohomology group \( H^2(E, \mathbb{Z}) \). It is an even unimodular lattice with the bilinear form given by the cup-product. The signature of this lattice is \((5, 29)\) by the Hodge index theory. It follows from the classification of even unimodular lattices that \( L \) is isometric to \( U^5 \oplus E_8(-1)^3 \), where \( U \) denotes the hyperbolic matrix and \( E_8(-1) \) denotes the unique negative definite even unimodular lattice of rank 8.

By the Shioda–Tate formula [Shi72], the Picard number is \( 2 + 2(15 - 1) = 30 \). The Néron–Severi group \( NS(E) \), which is the group of divisors on \( E \) modulo algebraic equivalence, is a torsion-free \( \mathbb{Z} \)-module of rank 30. This group can be imbedded into \( L \) by a cohomology sequence. The determinant of this sublattice, by a formula in [Shi72], is equal to

\[
| \det(NS(E)) | = 15^2 = 9.
\]

The orthogonal complement \( T_E \) of \( NS(E) \) in \( L \), called the transcendental lattice of \( L \), has rank 4 and \( | \det(T_E) | = | \det(NS(E)) | = 9 \) since \( L \) is unimodular.

As we are interested in the arithmetic properties of \( E \), we shall consider the reductions of \( E \). It turns out that for this particular elliptic surface \( E \) the only bad prime is 3. (The prime 5 is good because the 5 torsion points have killed the contribution of 5 from the
special fibres.) Hence we may regard $E$ as a normal connected smooth scheme over $\mathbb{Z}[1/3]$; it is tamely ramified along the closed subscheme formed by the cusps.

3. Determining the noncongruence subgroup $\Gamma$

For any positive integer $N$ let

$$
\Gamma^0(N) = \left\{ \begin{pmatrix} a & b \\ c & d \end{pmatrix} \in SL_2(\mathbb{Z}) : N|b \right\},
$$

$$
\Gamma^1(N) = \left\{ \begin{pmatrix} a & b \\ c & d \end{pmatrix} \in \Gamma^0(N) : a \equiv d \equiv 1 \mod N \right\}.
$$

With $t^3$ in (8) replaced by $t$, the new equation defines an elliptic modular surface $E'$ over the modular curve for the group $\Gamma^1(5)$. The surface $E'$ is a three-fold cover of $E'$, and thus the group $\Gamma$ is a subgroup of $\Gamma^1(5)$ of index three. We proceed to determine $\Gamma$ in terms of generators and relations.

First we decompose the full modular group $SL_2(\mathbb{Z})$ as

$$
SL_2(\mathbb{Z}) = \bigcup_{1 \leq i \leq 6} \Gamma^0(5)\gamma_i,
$$

where $\gamma_i = \begin{pmatrix} 1 & i-1 \\ 0 & 1 \end{pmatrix}$ for $1 \leq i \leq 5$ and $\gamma_6 = \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}$. Further,

$$
\Gamma^0(5) = \bigcup_{1 \leq j \leq 4} \Gamma^1(5)A^j = \pm\Gamma^1(5) \bigcup \pm\Gamma^1(5)A,
$$

where $A = \begin{pmatrix} -2 & -5 \\ 1 & 2 \end{pmatrix}$, $A^2 = -I$. Hence the coset representatives of $\pm\Gamma^1(5)$ in $SL_2(\mathbb{Z})$ may be taken as $\gamma_i$ and $A\gamma_i$ for $1 \leq i \leq 6$.

Listed below are the cusps of $\pm\Gamma^1(5)$ and a choice of generators of their stabilizers in $SL_2(\mathbb{Z})$:

<table>
<thead>
<tr>
<th>Cusps of $\pm\Gamma^1(5)$</th>
<th>Generators of stabilizers</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\infty$</td>
<td>$\gamma = \begin{pmatrix} 1 &amp; 5 \ 0 &amp; 1 \end{pmatrix}$</td>
</tr>
<tr>
<td>0</td>
<td>$\delta = \begin{pmatrix} 1 &amp; 0 \ -1 &amp; 1 \end{pmatrix}$</td>
</tr>
<tr>
<td>$-2$</td>
<td>$A\gamma A^{-1} = \begin{pmatrix} 11 &amp; 20 \ -5 &amp; -9 \end{pmatrix}$</td>
</tr>
<tr>
<td>$-\frac{5}{2}$</td>
<td>$A\delta A^{-1} = \begin{pmatrix} 11 &amp; 25 \ -4 &amp; -9 \end{pmatrix}$</td>
</tr>
</tbody>
</table>
Therefore the group $\Gamma^1(5)$ is generated by $\gamma$, $\delta$, $A\gamma A^{-1}$, $A\delta A^{-1}$ with the relation

$$(A\delta A^{-1})(A\gamma A^{-1})\delta \gamma = I.$$ 

In particular, $\Gamma^1(5)$ is actually generated by $A\delta A^{-1}$, $\delta$, $\gamma$.

For $\Gamma$, we may assume that its cusp at $\infty$ has width 15 so that

$$\Gamma^1(5) = \bigcup_{0 \leq j \leq 2} \Gamma^j.$$ 

The information on types of special fibers of $E$ and the above table give rise to the following information on cusps of $\Gamma$ and a choice of generators of stabilizers of each cusp:

<table>
<thead>
<tr>
<th>Cusps of $\Gamma$</th>
<th>Width</th>
<th>Generators of stabilizers</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\infty$</td>
<td>15</td>
<td>$\gamma^3$</td>
</tr>
<tr>
<td>$-2$</td>
<td>15</td>
<td>$A\gamma^3 A^{-1}$</td>
</tr>
<tr>
<td>0</td>
<td>1</td>
<td>$\delta$</td>
</tr>
<tr>
<td>5</td>
<td>1</td>
<td>$\gamma\delta \gamma^{-1}$</td>
</tr>
<tr>
<td>10</td>
<td>1</td>
<td>$\gamma^2 \delta \gamma^{-2}$</td>
</tr>
<tr>
<td>$-5$</td>
<td>1</td>
<td>$A\delta A^{-1}$</td>
</tr>
<tr>
<td>$4\tau/7$</td>
<td>1</td>
<td>$\gamma A\delta A^{-1} \gamma^{-1}$</td>
</tr>
<tr>
<td>$15/2$</td>
<td>1</td>
<td>$\gamma^2 A\delta A^{-1} \gamma^{-2}$</td>
</tr>
</tbody>
</table>

This shows that $\Gamma$ is generated by $\gamma^3$, $\delta$, $A\gamma^3 A^{-1}$, $A\delta A^{-1}$, $\gamma\delta \gamma^{-1}$, $\gamma A\delta A^{-1} \gamma^{-1}$, $\gamma^2 \delta \gamma^{-2}$, $\gamma^2 A\delta A^{-1} \gamma^{-2}$ with the relation

$$(A\delta A^{-1})(A\gamma^3 A^{-1})\delta (\gamma A\delta A^{-1} \gamma^{-1})(\gamma \delta \gamma^{-1})(\gamma^2 A\delta A^{-1} \gamma^{-2})(\gamma^2 \delta \gamma^{-2})\gamma^3 = I.$$ 

**Remark 3.1.** Similar to the above, if we take elements $\gamma^2$, $\delta$, $A\delta A^{-1}$, $A\gamma^2 A^{-1}$, $\gamma\delta \gamma^{-1}$, $\gamma A\delta A^{-1} \gamma^{-1}$ as generators with the relation

$$(A\delta A^{-1})(A\gamma^2 A^{-1}) (\delta)(\gamma A\delta A^{-1})(\gamma \delta \gamma^{-1})(\gamma^2) = I,$$ 

then we get the noncongruence subgroup $\Gamma_2$ of $\Gamma^1(5)$ of index 2 associated to the elliptic modular surface defined by

$$y^2 + (1 - t^2)xy - t^2 y = x^3 - t^2 x^2.$$ (12)
4. The space of weight 3 cusp forms for $\Gamma$

It follows readily from the dimension formula in [Shi71] that the space $S_3(\Gamma)$ of cusp forms of weight 3 for $\Gamma$ has dimension 2. We shall give a basis of this space in terms of the weight 3 Eisenstein series of $\Gamma^1(5)$ by using Hecke’s construction as described in [Ogg69]. The space of weight 3 Eisenstein series of $\Gamma^1(5)$ has dimension 4, equal to the number of cusps of $\Gamma^1(5)$. We are only interested in the two Eisenstein series that vanish at all but only one of the two cusps $\infty$ and $-2$.

Let $k \geq 3$, and $c, d \in \mathbb{Z}$. The Eisenstein series

$$G_k(\tau; (c, d); N) = \sum_{\substack{m \equiv c \\ n \equiv d \pmod{N}}} (m \tau + n)^{-k}$$

is a weight $k$ modular form for the principal congruence subgroup $\Gamma(N)$. Moreover, for any $L \in SL_2(\mathbb{Z})$, we have

$$G_k(\tau; (c, d); N)|L = G_k(\tau; (c, d)L; N).$$

This Eisenstein series has the following Fourier expansion:

$$G_k(\tau; (c, d); N) = \sum_{\lambda=0}^{\infty} a_{\lambda} z^\lambda, \quad z = e^{2\pi i \tau/N}, \quad (13)$$

where

$$a_0 = \begin{cases} 0 & \text{if } c \not\equiv 0 \pmod{N}, \\ \sum_{n \equiv d \pmod{N}} n^{-k} & \text{if } c \equiv 0 \pmod{N} \end{cases}$$

and for $\lambda \geq 1$,

$$a_{\lambda} = \frac{(-2\pi i)^k}{N^k \Gamma(k)} \sum_{\substack{m \equiv c \\ (m, N)=1}} (\text{sgn} \nu) \nu^{k-1} e^{2\pi i \nu d/N}. \quad (14)$$

The restricted Eisenstein series is defined by

$$G_k^*(\tau; (c, d); N) = \sum_{\substack{m \equiv c \\ n \equiv d \pmod{N}}} (m \tau + n)^{-k},$$

which is a weight $k$ modular form for $\Gamma(N)$, satisfying

$$G_k^*(\tau; (c, d); N)|L = G_k^*(\tau; (c, d)L; N).$$
for all \( L \in SL_2(\mathbb{Z}) \). Let \( \mu(n) \) denote the Möbius function. Then \( G_k^*(\tau; (c, d); N) \) can be expressed in terms of the series \( G_k(\tau; (c, d); N) \):

\[
G_k^*(\tau; (c, d); N) = \sum_{a=1}^{\infty} \mu(a)a^{-k}G_k(\tau; (a'c, a'd); N)
\]

\[
= \sum_{(t,N)=1 \atop t \mod N} c_t \cdot G_k(\tau; (ct, dt); N),
\]

where \( a' \) is chosen such that \( aa' \equiv 1 \pmod{N} \), and \( c_t = \sum_{at \equiv 1 \pmod{N} \atop a > 0} \mu(a)a^{-k} \).

The Eisenstein series \( G_k^*(\tau; (c, d); N) \) has value 1 at the cusp \(-\frac{d}{c}\) and 0 at all other cusps.

In the special case \( k = 3, N = 5 \), for any character \( \chi \) of \((\mathbb{Z}/5\mathbb{Z})^*\) we have

\[
\sum_{(t,N)=1} \tilde{\eta}(t)c_t = \sum_{t} \sum_{a \equiv t^{-1} \atop a > 0} \mu(a)a^{-k} = L^{-1}(k, \chi).
\]

Denote by \( \chi_3 : (\mathbb{Z}/5\mathbb{Z})^* \to \mathbb{C} \) the character given by \( \chi_3(2) = i, \chi_2 = \chi_3^2, \chi_4 = \chi_3^2, \) and \( \chi_1 \) the trivial character of \((\mathbb{Z}/5\mathbb{Z})^*\). Then the constants \( c_t \) can be expressed via the values of L-series:

\[
c_1 = \frac{1}{4} \left( L^{-1}(3, \chi_1) + L^{-1}(3, \chi_2) + L^{-1}(3, \chi_3) + L^{-1}(3, \chi_4) \right),
\]

\[
c_2 = \frac{1}{4} \left( L^{-1}(3, \chi_1) - L^{-1}(3, \chi_2) + iL^{-1}(3, \chi_3) - iL^{-1}(3, \chi_4) \right),
\]

\[
c_3 = \frac{1}{4} \left( L^{-1}(3, \chi_1) - L^{-1}(3, \chi_2) - iL^{-1}(3, \chi_3) + iL^{-1}(3, \chi_4) \right),
\]

\[
c_4 = \frac{1}{4} \left( L^{-1}(3, \chi_1) + L^{-1}(3, \chi_2) - L^{-1}(3, \chi_3) - L^{-1}(3, \chi_4) \right).
\]

Using the functional equation of the L-function \( L(s, \chi) \) and the Bernoulli polynomials, we obtain two explicit L-values

\[
L(3, \chi_3) = \frac{\tau(\chi_3)}{2\pi} \left(-\frac{1}{2}\right) \left(\frac{2\pi}{5}\right)^3 (-\frac{1}{3}) \frac{6}{5} (2 - i),
\]

\[
L(3, \chi_4) = \frac{\tau(\chi_4)}{2\pi} \left(-\frac{1}{2}\right) \left(\frac{2\pi}{5}\right)^3 (-\frac{1}{3}) \frac{6}{5} (2 + i),
\]

where \( \tau(\chi) \) denotes the Gauss sum of the character \( \chi \).

As \( T_4^1(5) = \bigcup_{a=0}^4 T(5) \left(\begin{array}{cc} 1 & 0 \\ a & 1 \end{array}\right) \), we let

\[
E_1(\tau) = \sum_{0 \leq a \leq 4} G_k^*(\tau; (0, 1); 5)\left|\begin{array}{cc} 1 & 0 \\ a & 1 \end{array}\right| = \sum_a G_k^*(\tau; (a, 1); 5)
\]

\[
= c_1 \sum_a G_k(\tau; (a, 1); 5) + c_2 \sum_a G_k(\tau; (2a, 2); 5)
\]

\[
+ c_3 \sum_a G_k(\tau; (3a, 3); 5) + c_4 \sum_a G_k(\tau; (4a, 4); 5).
\]

(16)
Applying the Fourier expansion of the Eisenstein series \( G_k(\tau; (c, d); 5) \) and the equations for \( c_i' \)'s, we have the Fourier expansion of \( E_1(\tau) \):

\[
E_1(\tau) = 1 - \frac{1}{2} \sum_{\lambda=1}^{\infty} \left( 2 \sum_{v|\lambda, v>0} v^2 (\chi_3(v) + \chi_4(v)) + i \sum_{v|\lambda, v>0} v^2 (\chi_4(v) - \chi_3(v)) \right) \frac{q^{\lambda/5}}{5}.
\]  

(17)

In particular, the Fourier coefficients of \( E_1(\tau) \) are rational integers. And from the definition of \( E_1(\tau) \), it has value 1 at the cusp \( \infty \), and 0 at the other cusps 0, \(-2\), \(-\frac{5}{2}\).

Let

\[
E_2(\tau) = E_1(\tau)|A^{-1}
\]

\[
= c_1 \sum_a G_k(\tau; (-2a + 1, -5a + 2); 5) \\
+ c_2 \sum_a G_k(\tau; (-4a + 2, -10a + 4); 5) \\
+ c_3 \sum_a G_k(\tau; (-6a + 3, -15a + 6); 5) \\
+ c_4 \sum_a G_k(\tau; (-8a + 4, -20a + 8); 5).
\]

Then \( E_2(\tau) \) has value \(-1\) at the cusp \(-2\), and 0 at other cusps. In the same way we calculate the Fourier expansion of this Eisenstein series to get

\[
E_2(\tau) = \frac{1}{2} \sum_{\lambda=1}^{\infty} \left( \sum_{v|\lambda, v>0} v^2 (\chi_3(v) + \chi_4(v)) + 2i \sum_{v|\lambda, v>0} v^2 (\chi_3(v) - \chi_4(v)) \right) q^{\lambda/5}.
\]  

(18)

Since \( E_1(\tau) \) has values 1, 0, 0, 0 at the cusps \( \infty, -2, 0, -\frac{5}{2} \), respectively, and \( E_2(\tau) \) has values 0, \(-1\), 0, 0 at these cusps, both modular forms have no other zero points.

Consider the natural covering map

\[
\Gamma\backslash\mathcal{H} \to \Gamma^1(5)\backslash\mathcal{H},
\]

where \( \mathcal{H} \) denotes the Poincaré upper half plane. It ramifies only at the two cusps \( \infty \) and \(-2\), with index 3. Therefore the two functions

\[
f_1 = \sqrt[3]{E_1^2(\tau)E_2(\tau)} \quad \text{and} \quad f_2 = \sqrt[3]{E_1(\tau)E_2^2(\tau)}
\]  

(19)
are well-defined entire modular forms of weight 3 for $\Gamma$. Further, they vanish at every cusp, hence they are cusp forms. It is clear that $f_1$ and $f_2$ are linearly independent, and thus form a basis of the space $S_3(\Gamma)$. Choosing a proper cubic root of one, we may assume that the Fourier coefficients of both $f_1$ and $f_2$ are rational numbers with denominators involving only powers of 3.

Since the action of the matrix $A$ interchanges the cusp $\infty$ with the cusp $-2$, it defines an operator on the space $S_3(\Gamma)$. More precisely, its actions on $f_1$ and $f_2$ are

\[(f_1)|A = f_2 \quad \text{and} \quad (f_2)|A = -f_1.\]

Thus the operator $A$ on $S_3(\Gamma)$ has eigenforms $f_+ = f_1 + if_2$ and $f_- = f_1 - if_2$ with eigenvalues $-i$ and $i$, respectively. The Fourier expansions of these two eigenforms are as follows:

\[
f_+(\tau) = q^{1/15} + iq^{2/15} - \frac{11}{3} q^{4/15} - i \frac{16}{3} q^{5/15} - \frac{4}{9} q^{7/15} + i \frac{71}{9} q^{8/15} + \frac{932}{81} q^{10/15} + i \frac{247}{81} q^{11/15} + \frac{443}{243} q^{13/15} - i \frac{3832}{243} q^{14/15} - \frac{13151}{729} q^{16/15} + i \frac{9131}{729} q^{17/15} + O(q^{18/15}),
\]

\[
f_-(\tau) = q^{1/15} - iq^{2/15} - \frac{11}{3} q^{4/15} + i \frac{16}{3} q^{5/15} - \frac{4}{9} q^{7/15} - i \frac{71}{9} q^{8/15} + \frac{932}{81} q^{10/15} - i \frac{247}{81} q^{11/15} + \frac{443}{243} q^{13/15} + i \frac{3832}{243} q^{14/15} - \frac{13151}{729} q^{16/15} - i \frac{9131}{729} q^{17/15} + O(q^{18/15}).
\]

This proves the first assertion of Theorem 1.4.

Let $X$ be the modular curve of the noncongruence subgroup $\Gamma$, that is, $\mathcal{X}(\mathbb{C}) = \Gamma \backslash \mathcal{H}$. As seen in Section 2, it is a projective line over $\mathbb{C}$. The two cusp forms constructed in Section 4 give rise to a Hauptmodul of $X$:

\[t = \frac{f_1}{f_2} = \sqrt[3]{\frac{E_1}{E_2}}.\]

Since the Fourier coefficients of $t$ are in $\mathbb{Q}$, the curve $X$ is defined over $\mathbb{Q}$. It is easy to check from the generators and relation exhibited in Section 3 that the matrix $A$ lies in the normalizer of the noncongruence subgroup $\Gamma$ in $SL_2(\mathbb{Z})$. Therefore $A$ induces a $\mathbb{Q}$-rational involution on the modular curve $X$, given by

\[A(t) = -\frac{1}{t}. \quad (20)\]
Further, $A$ induces an order 4 $\mathbb{Q}$-rational action on the elliptic surface $E$. We first notice that the $j$-function \((9)\) is invariant if we send $t$ to $-1/t$. To see the action more explicitly, make the following change of variables:

\[
x = t^3 X - 1/12t^6 + 1/2t^3 - 1/12,
\]
\[
y = t^4 Y + 1/2t^6 X - 1/2t^3 X - 1/24t^9 + 7/24t^6 + 5/24t^3 + 1/24
\]

so that the original defining equation \((8)\) becomes

\[
Y^2 = t \left( X^3 - \frac{1 + 12t^3 + 14t^6 - 12t^9 + t^{12}}{48t^6} X \\
+ \frac{1 + 18t^3 + 75t^6 + 75t^{12} - 18t^{15} + t^{18}}{864t^9} \right).
\]

The action of $A$ sends $t$ to $-1/t$, $Y$ to $Y/t$, and $X$ to $-X$. Hence it is defined over $\mathbb{Q}$ and has order 4.

5. The $l$-adic representation attached to $S_3(\Gamma)$

As explained in Section 1, given a noncongruence subgroup $\Gamma'$ of $SL_2(\mathbb{Z})$ of finite index with the modular curve $X_{\Gamma'}$ defined over $\mathbb{Q}$, Scholl [Sch85] defined a compatible family of $l$-adic representations $\rho_l$ of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ associated to $S_k(\Gamma')$, $k \geq 3$, from which he derived the congruence relations Theorem 1.2. When $k = 3$, the representation $\rho_l$ is defined as follows. Choose an integer $N \geq 3$ such that $\pm \Gamma' \Gamma(N) = SL_2(\mathbb{Z})$. Denote by $X(N)$ the compactified modular curve for the principal congruence subgroup $\Gamma(N)$, and by $X(N)^0$ the part of $X(N)$ with cusps removed. Let $G(N) = SL(\mu_N \times \mathbb{Z}/N)$, let $f_{\text{univ}} : E_{\text{univ}} \to X(N)^0$ be the restriction to $X(N)^0$ of the universal elliptic curve of $X(N)$, and let $V(N)$ (resp. $V(N)^0$) be the normalization of the fiber product $X_{\Gamma'} \times_{X(1)} X(N)$ (resp. $X_{\Gamma'} \times_{X(1)} X(N)^0$). The finite group scheme $G(N)$ acts on the second factor of $V(N)$, $E_{\text{univ}}$, and the sheaf $\mathcal{F}_{\text{univ}} = R^1 f_{\text{univ}}^* \mathbb{Q}_l$, respectively. We have the projection map $\pi_0 : V(N)^0 \to X(N)^0$ and the inclusion map $i_N : V(N)^0 \to V(N)$. The representation $\rho_l$ of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ is the action of the Galois group on the $\mathbb{Q}_l$-space

\[
H^1(V(N) \otimes \overline{\mathbb{Q}}, (i_N)_* \pi_0^* \mathcal{F}_{\text{univ}})^{G(N)}.
\]

The reason that an auxiliary modular curve $X(N)$ is involved is that the curve $X_{\Gamma'}$ does not have a universal elliptic curve, while $X(N)$ for $N \geq 3$ does. As shown above, the $l$-adic sheaf comes from this universal elliptic curve. At the end, $G(N)$ invariants are taken to rid the dependence of $X(N)$. 

When there is an elliptic surface $E'$ over the modular curve $X_{\Gamma'}$ with $h': E' \to X_{\Gamma'}$ tamely ramified along the cusps and elliptic points, inspired by [Sch88], we introduce another $l$-adic representation $\rho_l^*$ of the Galois group of $\mathbb{Q}$ using $E'$ as follows. Let $X_{\Gamma'}^0$ be the part of $X_{\Gamma'}$ with the cusps and elliptic points removed. Denote by $i$ the inclusion from $X_{\Gamma'}^0$ into $X_{\Gamma'}$, and by

$$h : E' \to X_{\Gamma'}^0$$

the restriction map, which is smooth. For almost all prime $l$, we obtain a sheaf

$$\mathcal{F}_l = R^1 h_* \mathbb{Q}_l$$

on $X_{\Gamma'}^0$. The $i_*$ map then transports it to a sheaf $i_* \mathcal{F}_l$ on $X_{\Gamma'}$. The action of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ on the $\mathbb{Q}_l$-space

$$W_l = H^1(X_{\Gamma'} \otimes \overline{\mathbb{Q}}, i_* \mathcal{F}_l) \quad (23)$$

defines an $l$-adic representation, denoted by $\rho_l^*$, of the Galois group of $\mathbb{Q}$.

The following diagram depicts the relationship of the curves and surfaces involved. For a scheme $X$ and a nonzero integer $M$, we use $X[1/M]$ to denote $X \times_{\text{Spec} \mathbb{Z}} \text{Spec} \mathbb{Z}[1/M]$. The integer $M$ below is chosen so that both modular curves $X(N)$ and $X_{\Gamma'}$ are smooth and proper over the ring $\mathbb{Z}[1/M]$. The maps $\pi_0$ and $\pi_{\Gamma'}$ are the natural projections.

```
E^{univ} \downarrow \pi_0^{univ} \quad \pi_0 \quad \pi_{\Gamma'} \quad \pi_{\Gamma'}
\downarrow \quad V^0(N)[1/M] \quad i_N \quad V(N)[1/M] \quad X_{\Gamma'}^0[1/M]
X(N)^0 \leftarrow \quad X_{\Gamma'} \times_{X(1)} X(N)^0[1/M] \quad \rightarrow \quad X_{\Gamma'} \times_{X(1)} X(N)[1/M] \quad \rightarrow \quad X_{\Gamma'}[1/M]
```

**Proposition 5.1.** The two representations $\rho_l^*$ and $\rho_1$ are isomorphic up to a twist by a character $\phi_l$ of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ of order at most 2.

**Proof.** Using the projection $\pi_{\Gamma'}$ from the fiber product $V(N)[1/M]$ to the factor $X_{\Gamma'}^0[1/M]$, we pull back the sheaf $i_* \mathcal{F}_l$ on $X_{\Gamma'}[1/M]$ to the sheaf $\pi_{\Gamma'}^* i_* \mathcal{F}_l$ on $V(N)[1/M]$. It follows from an argument similar to Section 1.3 of [Sch88] that the sheaf $(i_N)_* \pi_0^* \mathcal{F}_l^{univ}$ is isomorphic to the sheaf $\pi_{\Gamma'}^* i_* \mathcal{F}_l \otimes \mathcal{L}$, where $\mathcal{L}$ is a rank
one sheaf on \(V(N)[1/M]\) with \(L^{\otimes 2}\) isomorphic to the constant sheaf \(\mathbb{Q}_l\). Consequently, \(H^1(V(N) \otimes \overline{\mathbb{Q}}, \pi^*_{l,i_*F_l})\) and \(H^1(V(N) \otimes \overline{\mathbb{Q}}, (i_N)_*\pi^*_{0,F_l}^{\text{univ}})\) are isomorphic \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})\) modules up to a twist by a character \(\phi_l\) of \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})\) of order at most 2. Since \(G(N)\) acts only on \(X(N)\), the \(G(N)\)-invariant part of \(H^1(V(N) \otimes \overline{\mathbb{Q}}, \pi^*_{l,i_*F_l})\) is isomorphic to \(W_l\). Therefore, the representation \(\rho^*_l\) on \(W_l\) is isomorphic to \(\rho_l\) on \(H^1(V(N) \otimes \overline{\mathbb{Q}}, (i_N)_*\pi^*_{0,F_l}^{\text{univ}})^{G(N)}\) up to a twist by \(\phi_l\).

Apply the above discussion to the case where \(I' = I\). We shall show later in Section 6 that our \(\rho^*_l\) is in fact isomorphic to Scholl’s representation \(\rho_l\). The simpler description of \(\rho^*_l\) allows us to get more information about the representation, and eventually leading to a finer congruence result than the one provided by Theorem 1.2.

By Scholl’s result in [Sch85], \(\dim_{\mathbb{Q}_2} W_l = h^{2,0}(\mathcal{E}) + h^{0,2}(\mathcal{E}) = 4\). As remarked at the end of the previous section, the action of \(A\) on \(X\) and on \(\mathcal{E}\) are both \(\mathbb{Q}\)-rational, thus \(A\) commutes with the action of \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})\) on the space \(W_l\). Moreover, the action of \(A\) on the sheaf \(\mathcal{F}_l\) or the representing space \(W_l\) has order 4. This makes \(W_l\) a 2-dimensional module over the algebra \(K = \mathbb{Q}_2(A)\).

Now fix the prime \(l = 2\). Then \(K = \mathbb{Q}_2(A)\) is a degree two field extension of \(\mathbb{Q}_2\) and \(\rho^*_2\) is a degree two representation of \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})\) over the 2-dimensional vector space \(W_2\) over \(K\). We calculate the characteristic polynomial \(H_p(T)\) (resp. \(H'_p(T)\)) of \(\rho^*_2(\text{Frob}_p)\) over \(\mathbb{Q}_2\) (resp. \(\mathbb{Q}_2(A) = K = \mathbb{Q}(A)_{1+A}\)) for varying primes \(p \neq 2, 3\). By Scholl’s work [Sch85], the characteristic polynomial \(H_p(T)\) of \(\rho^*_2(\text{Frob}_p)\) can be factored as

\[
H_p(T) = (T - \alpha_p)(T - \beta_p)(T - p^2/\alpha_p)(T - p^2/\beta_p)
= T^4 - C_1(p)T^3 + C_2(p)T^2 - p^2C_1(p)T + p^4 \in \mathbb{Z}[T],
\]

where

\[
C_1(p) = \alpha_p + \beta_p + p^2/\alpha_p + p^2/\beta_p = \text{Tr}(\rho^*_2(\text{Frob}_p)),
C_2(p) = \frac{1}{2} (C_1^2 - \text{Tr}(\rho^*_2(\text{Frob}_p))) = \frac{1}{2} \left( (\text{Tr}(\rho^*_2(\text{Frob}_p)))^2 - \text{Tr}(\rho^*_2(\text{Frob}_p)) \right).
\]

Since

\[
H^j(X_{T'} \otimes \overline{F}_p, i_*F_l) = 0 \quad \text{for} \quad j \neq 1,
\]

we can easily get from the Lefschetz formula the following trace formula of \(\text{Tr}(\rho^*_2(\text{Frob}_q))\) for \(q = p^r\):

\[
\text{Tr}(\rho^*_2(\text{Frob}_q)) = - \sum_{x \in X(\overline{F}_q)} \text{Tr}(x),
\]
where

\[ \text{Tr}(x) = \text{Tr} \left( (Frob_q)_x : (i_* F_2)_x \right) \]

is the trace of Frob\(_q\) restricted to the stalk at \(x\) of the sheaf \(i_* F_2\).

We proceed to compute \(\text{Tr}(x)\). Denote by \(E_x\) the fibre of the elliptic surface \(f : E \to X\) at the point \(x \in X(F_q)\); it is a curve of genus 0 or 1 depending on whether the discriminant of \(E_x\) vanishes in the field \(F_q\) or not. Note that there is no need to treat the cases \(j = 0\) and 1728 separately as in [Sch88] since we are computing the traces by using the explicit equation of the elliptic surface \(E \to X\). In our case, for \(x \in X(F_q)\) we always have

\[ \text{Tr}(x) = 1 + q - \#E_x(F_q). \]

A computer program yields the following table for \(\text{Tr}_q\), the traces of Frob\(_q\):

<table>
<thead>
<tr>
<th>(p)</th>
<th>5</th>
<th>7</th>
<th>11</th>
<th>13</th>
<th>17</th>
<th>19</th>
<th>23</th>
<th>29</th>
<th>31</th>
</tr>
</thead>
<tbody>
<tr>
<td>(\text{Tr}_p)</td>
<td>0</td>
<td>10</td>
<td>0</td>
<td>-20</td>
<td>0</td>
<td>-32</td>
<td>0</td>
<td>0</td>
<td>-2</td>
</tr>
<tr>
<td>(\text{Tr}_p^2)</td>
<td>82</td>
<td>-146</td>
<td>34</td>
<td>-476</td>
<td>508</td>
<td>-932</td>
<td>1828</td>
<td>1564</td>
<td>-3842</td>
</tr>
</tbody>
</table>

Therefore we obtain the characteristic polynomials \(H_p(T)\) of \(\rho^*_p(Frob_p)\), with \(p\) the primes between 5 and 31. For such a prime \(p\), the characteristic polynomial \(H'_p(T)\) of \(\rho^*_p(Frob_p)\) over \(K\) has degree 2, and it has the property that \(H'_p(T)H''_p(T) = H_p(T)\), where \(H''_p(T)\) is the conjugate of \(H'_p(T)\) under the automorphism of \(K\) over \(Q_2\) sending \(A\) to \(-A\). Since Scholl [Sch85] proved that all roots of \(H_p\) are algebraic integers, the first step towards determining \(H'_p\) is to figure out how to separate the four roots of \(H_p\) into two conjugate pairs to form the roots of \(H'_p\) and \(H''_p\). It turns out that for the primes from 5 to 31 such separation is unique except for \(p = 13\) and 19. For these primes, while we cannot choose between \(H'_p\) and \(H''_p\) without further work, we do know the determinants of the 2-dimensional representation \(\rho^*_2\) over \(K\) at these primes since the constant term of \(H'_p(T)\) is \(\pm p^2\).

Next we prove that the information so far determines the determinants of \(\rho^*_2\), which in turn will allow us to determine \(H'_{13}\) and \(H'_{19}\).

Lemma 5.2. If two integral 1-dimensional representations \(\sigma_1\) and \(\sigma_2\) of the group \(\text{Gal}(\overline{Q}/Q)\) over the field \(Q_2(A)\), which are unramified away from 2 and 3, agree on the elements Frob\(_p\) for \(p = 5, 7, 11, 17\), then they are equal.

Proof. The images of the 1-dimensional representations are in \((\mathbb{Z}_2[A])^* = (\mathbb{Z}(A)_p)^*\), where \(p = (1 + A)\) is the maximal ideal of the local ring \(\mathbb{Z}_2(A)\). Note that \((\mathbb{Z}(A)_p)^* = \langle A \rangle \times (1 + p^3)\). Let \(\log\) denote the \((1 + A)\)-adic logarithm on \((\mathbb{Z}_2[A])^*\). More precisely,
it has kernel the group \( \langle A \rangle \) of roots of unity in \( \mathbb{Q}_2(A) \), and it maps \( 1 + x \in 1 + p^3 \) to \( \sum_{n=1}^{\infty} \frac{(-1)^{n-1}}{n} x^n \in \mathbb{Z}_2[A] \). As such, \( \log \) gives an isomorphism between the multiplicative group \( 1 + p^3 \) and the additive group \( p^3 \). Consider

\[
\psi = \log \circ \sigma_1 - \log \circ \sigma_2,
\]

which is a homomorphism from \( \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \) to \( p^3 \).

If \( \psi \neq 0 \), then

\[
n_0 = \min \left\{ \text{ord}_p \left( \psi(\tau) \right) : \tau \in \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \right\}
\]

is finite. Then

\[
\tilde{\psi} := \frac{1}{(1 + A)^{n_0}} \psi \pmod{p}
\]

is a continuous surjective homomorphism from \( \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \) to \( \mathbb{F}_2 \) which is trivial at the Frobenius elements at primes \( p = 5, 7, 11, 17 \) by assumption. This representation of \( \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \) factors through a quadratic extension of \( \mathbb{Q} \) unramified outside 2 and 3. Such fields are extensions of \( \mathbb{Q} \) by adjoining square roots of 2, 3, 6, \(-1\), \(-2\), \(-3\), \(-6\), respectively. It is easy to check that the prime \( p = 5, 7, 11, 17 \) is inert in the respective field, and thus \( \hat{\psi} \) at such \( \text{Frob}_p \) would be nontrivial, a contradiction. Therefore \( \psi = 0 \), in other words, the image of the representation \( \sigma := \sigma_1(\sigma_2)^{-1} \) is a subgroup of \( \langle A \rangle \), a cyclic group of order 4. Hence we consider all Galois extensions of \( \mathbb{Q} \) with group equal to a subgroup of a cyclic group of order 4, unramified away from 2 and 3, and in which \( p = 5, 7, 11, 17 \) split completely. If a nontrivial such extension exists, then it contains a quadratic subextension unramified outside 2 and 3, and in which \( p = 5, 7, 11, 17 \) split completely. As shown above, this is impossible. Therefore the image of \( \sigma \) can only be \( \{1\} \), in other words, \( \sigma_1 \) and \( \sigma_2 \) are equal. \( \square \)

Denote by \( \chi_{-3} \) the quadratic character attached to the field \( \mathbb{Q}(\sqrt{-3}) \), that is, \( \chi_{-3}(x) \) is equal to the Legendre symbol \( \left( \frac{-3}{x} \right) \). The 1-dimensional representation \( \sigma \) of \( \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \) over \( K \) given by \( \tilde{\chi}_{-3}(\text{Frob}_p) = \chi_{-3}(p)p^2 \) for primes \( p \neq 2, 3 \) agrees with \( \text{det}(\rho_2^*) \) at \( \text{Frob}_p \) for \( p = 5, 7, 11, 17 \) by checking the constant term of \( H'_p(T) \), and hence we conclude from Lemma 5.2 that the two degree one representations agree.

**Corollary 5.3.** Let \( \tilde{\chi}_{-3} \) be as above. We have \( \text{det}(\rho_2^*) = \tilde{\chi}_{-3} \).

In particular, we know that the constant term of \( H'_{13} \) (resp. \( H'_{19} \)) is \((13)^2\) (resp. \((19)^2\)). This information enables us to determine \( H'_{13}(T) \) and \( H'_{19}(T) \). We record the result so far in the following proposition.
the 4-dimensional 2-adic representation of the Galois group of \( \overline{\mathbb{Q}} \).

Proposition 5.4. The characteristic polynomials \( H_p(T) \) and \( H'_p(T) \) of the element \( \rho_2^*({\text{Frob}}_p) \) over \( \mathbb{Q}_2 \) and \( K = \mathbb{Q}_2(A) = \mathbb{Q}(A)_{1+1} \), respectively, for primes \( 5 \leq p \leq 31 \) are given in Table 1.

Note that in case \( p \equiv 2 \pmod{3} \), the coefficient of \( T \) in \( H'_p(T) \) is determined up to sign.

6. Comparison with the representation \( \tilde{\rho}_2 \) attached to certain cusp forms in \( S_3(\Gamma_1(27)) \)

In the space of weight 3 level 27 cusp forms \( S_3(\Gamma_1(27)) \), we find from Stein’s website [Stein] two Hecke eigenforms \( g_a \) whose \( q \) expansion (\( q = e^{2\pi i z} \)) to order 31 are as follows:

\[
g_a = q + aq^2 - 5q^4 - aq^5 + 5q^7 - aq^8 + 9q^{10} - 5aq^{11} - 10q^{13} + 5aq^{14} - 11q^{16} + 6aq^{17} - 16q^{19} + 5aq^{20} + 45q^{22} - 4aq^{23} + 16q^{25} - 10aq^{26} - 25q^{28} + 10aq^{29} - q^{31} + O(q^{32}),
\]

where \( a \) is a root of \( x^2 + 9 \). The character of this modular form is \( \chi_{-3} \). Denote by \( \tilde{\rho}_2 \) the 4-dimensional 2-adic representation of the Galois group of \( \mathbb{Q} \) attached to \( g_a + g_{-a} \), established by Deligne [Del]. The Atkin–Lehner operator \( H_{27} \) acts on the curve \( X_1(27) \) as an involution, and it is \( \mathbb{Q} \)-rational. Further, it induces an action of order 4 on the representation space of \( \tilde{\rho}_2 \) so that \( \tilde{\rho}_2 \) may be regarded as a 2-dimensional representation over \( \mathbb{Q}_2(i) = \mathbb{Q}(i)_{1+i} \). In view of Corollary 5.3, we have

Corollary 6.1. \( \det(\rho_2^*) = \det(\tilde{\rho}_2) \).

Fix an isomorphism from \( K \) to \( \mathbb{Q}_2(i) = \mathbb{Q}(i)_{1+i} \) such that the characteristic polynomial of \( \rho_2^*({\text{Frob}}_5) \) agrees with that of \( \tilde{\rho}_2({\text{Frob}}_5) \).
Denote by $\rho'_2$ the representation $\rho^\ast_2$ viewed over $\mathbb{Q}(i)_{1+i}$. Our goal in this section is to show that $\rho'_2$ and $\tilde{\rho}_2$ are isomorphic.

To compare two representations from $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ to $\text{GL}_2(\mathbb{Z}[i]_{1+i})$, we will apply Serre’s method [Ser84].

**Theorem 6.2** (Serre). Let $\rho_1$ and $\rho_2$ be representations of the group $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ to $\text{GL}_2(\mathbb{Z}[i]_{1+i})$. Assume they satisfy the following two conditions:

1. $\text{det}(\rho_1) = \text{det}(\rho_2)$;
2. the two homomorphisms from $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ to $\text{GL}_2(\mathbb{F}_2)$, obtained from the reductions of $\rho_1$ and $\rho_2$ modulo $1+i$, are surjective and equal.

If $\rho_1$ and $\rho_2$ are not isomorphic, then there exists a pair $(\tilde{G}, t)$, where $\tilde{G}$ is a quotient of the Galois group $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ isomorphic to either $S_4 \times \{\pm 1\}$, or $S_4$, or $S_3 \times \{\pm 1\}$, and the map $t : \tilde{G} \to \mathbb{F}_2$ has value 0 on the elements of $\tilde{G}$ of order $\leq 3$, and 1 on the other elements.

This result, explained in detail in [Ser84] and in a letter from Serre to Tate, is a specialization of a general idea to determine the “deviation” of two non-isomorphic representations over a local field. When the residue field of the local field is small, it gives a feasible method to determine an $l$-adic representation by checking the traces of Frobenii at a small number of primes. Originated from Faltings work [Fal83], this idea was made effective by Serre [Ser84]. In [Liv87] Livné described how to use Serre’s method to prove two representations with even traces to be isomorphic. We sketch Serre’s proof below.

Suppose two representations $\rho_1$ and $\rho_2$ are not isomorphic. Then the traces $\text{Tr}(\rho_1)$ and $\text{Tr}(\rho_2)$ are not identical on $G := \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$. Write $\pi$ for the uniformizer $1+i$ of the local field $\mathbb{Q}_2(i)$ for brevity. There exists a highest power $\pi^n$ for some integer $n$ ($\geq 1$ by condition (2)) such that

$$\text{Tr}(\rho_1(s)) \equiv \text{Tr}(\rho_2(s)) \mod \pi^n \quad \text{for all } s \in \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}).$$

This yields a nonconstant (and hence surjective) map

$$t : G \to \mathbb{F}_2,$$

$$s \mapsto \left(\text{Tr}(\rho_2(s)) - \text{Tr}(\rho_1(s))\right) / \pi^n \mod \pi,$$

which records the difference between the two representations $\rho_1$ and $\rho_2$. Next one seeks to pass this information to a manageable finite quotient $\tilde{G}$ of the Galois group $G$ so that the pair $(\tilde{G}, t)$ measures the difference of $\rho_1$ and $\rho_2$.

The three cases of “deviation” listed in Theorem 6.2 were derived by explicitly computing possible $\tilde{G}$ as follows. The hypothesis $\rho_1 \equiv \rho_2 \mod \pi^n$ implies that we may write, for $s \in G$,

$$\rho_2(s) = (1 + \pi^n a(s)) \rho_1(s) \quad \text{with } a(s) \in M_2(\mathcal{O}_K).$$
showing \( t(s) = \text{Tr}(a(s)\rho_1(s)) \mod \pi \). Hence it suffices to find a quotient \( \tilde{G} \) capturing \( a(s) \mod \pi \) and \( \rho_1(s) \mod \pi \) for all \( s \) in \( G \). The relation \( \rho_2(s_1s_2) = \rho_2(s_1)\rho_2(s_2) \) yields \( a(s_1s_2) = a(s_1) + \rho_1(s_1)a(s_2)\rho_1(s_1)^{-1} \mod \pi \) for all \( s_1, s_2 \) in \( G \). In other words, the map \( s \mapsto a(s) \mod \pi \) from \( G \) to \( M_2(\mathbb{F}_2) \) is a 1-cocycle under the adjoint action of the Galois group on \( M_2(\mathbb{F}_2) \) through \( \rho_1 \) modulo \( \pi \). The map \( \theta : s \mapsto (a(s) \mod \pi, \rho_1(s) \mod \pi) \) is a homomorphism from \( G \) to the semi-direct product \( M_2(\mathbb{F}_2) \rtimes GL_2(\mathbb{F}_2) \), where the group law is \((a, g) \cdot (b, h) = (a + gb^{-1}, gh)\). The desired group \( \tilde{G} \) is isomorphic to the image of \( \theta \). It remains to figure out the possible group structure of \( \tilde{G} \). The projection of \( \tilde{G} \to GL_2(\mathbb{F}_2) \cong S_3 \) is surjective by condition (2). Condition (1) implies that the trace of \( a(s) \mod \pi \) is zero. Therefore the projection of \( \tilde{G} \) to \( M_2(\mathbb{F}_2) \) is a subgroup of the trace zero elements in \( M_2(\mathbb{F}_2) \), which is generated by \( \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 0 & 0 \end{pmatrix}, \text{ and } \begin{pmatrix} 0 & 0 \\ 1 & 0 \end{pmatrix} \). Put together, one finds three possibilities for \( \tilde{G} \): they are \( \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 0 & 1 \\ 0 & 0 \end{pmatrix}, \text{ and } \begin{pmatrix} 0 & 0 \\ 1 & 0 \end{pmatrix} \) \rtimes GL_2(\mathbb{F}_2) \), \( \begin{pmatrix} 1 & 1 \\ 0 & 1 \end{pmatrix}, \begin{pmatrix} 1 & 0 \\ 1 & 1 \end{pmatrix} \) \rtimes GL_2(\mathbb{F}_2) \), and \( \begin{pmatrix} 1 & 0 \\ 0 & 1 \end{pmatrix} \) \rtimes GL_2(\mathbb{F}_2) \), corresponding to the three cases given in the theorem.

It follows from Serre’s Theorem that two nonisomorphic representations must have different traces at elements in \( \tilde{G} \) of order at least 4. Therefore, to show that two representations \( \rho_2 \) and \( \tilde{\rho}_2 \) are isomorphic (and necessarily ramify at the same places), our strategy is to search for Galois extensions of \( \mathbb{Q} \) with Galois groups isomorphic to those listed in Theorem 6.2 and unramified where the representations are unramified. In each of such extensions, if we can find a Frobenius element of order \( \geq 4 \) in the Galois group at which the two representations have the same trace, then they must be isomorphic.

Corollary 6.1 shows that condition (1) holds. We now proceed to prove that both representations satisfy condition (2) as well. Note that the residue field of \( \mathbb{Q}(i) \) at \( 1+i \) is \( \mathbb{F}_2 \), hence the reductions of \( \rho_2' \) and \( \tilde{\rho}_2 \) mod \( 1+i \) yield two representations of the Galois group of \( \mathbb{Q} \) to \( GL_2(\mathbb{F}_2) \). While we do not know the characteristic polynomial of the Frobenius at almost all primes \( p \geq 5 \) for the representation \( \rho_2' \), we do know them modulo \( 1+i \) for \( 5 \leq p \leq 31 \) from Table 1, and it is easy to check that they agree with those from representation \( \tilde{\rho}_2 \) mod \( 1+i \). Thus condition (2) for our two representations will follow from

**Lemma 6.3.** There is only one representation \( \rho \) from \( \text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \) to \( GL_2(\mathbb{F}_2) \), unramified outside 2 and 3, such that the characteristic polynomial of \( \rho(Frob_p) \) is equal to \( H'_p \mod (1 + A) \) from Table 1 for primes \( p = 5, 7, 13 \). Further \( \rho \) is surjective.

**Proof.** The existence is obvious. We prove the uniqueness. Note that \( GL_2(\mathbb{F}_2) \) is isomorphic to the symmetric group on three letters \( S_3 \), which is generated by an element of order 2 and an element of order 3. Denote by \( \sigma \) the sign homomorphism from \( S_3 \) to \( \{\pm 1\} \), and by \( \varepsilon \) the composition \( \sigma \circ \rho \). Then \( \varepsilon \) is a character of the Galois group of \( \mathbb{Q} \) of order at most two and it is unramified outside 2 and 3.
We see from Table 1 that the characteristic polynomials of \(\rho(\text{Frob}_5)\) and \(\rho(\text{Frob}_7)\) are \(T^2 + T + 1\), hence \(\rho(\text{Frob}_5)\) and \(\rho(\text{Frob}_7)\) both have order 3. The characteristic polynomial of \(\rho(\text{Frob}_{13})\) is \(T^2 + 1\), hence \(\rho(\text{Frob}_{13})\) has order 1 or 2. In particular, \(\psi(\text{Frob}_p) = 1\) for \(p = 5, 7\). If \(\psi\) is nontrivial, then it arises from a quadratic extension of \(\mathbb{Q}\) unramified outside 2 and 3. Such extensions are \(\mathbb{Q}(\sqrt{d})\) with \(d = 2, 3, 6, -1, -2, -3, -6\). Since 5 is inert in \(\mathbb{Q}(\sqrt{d})\) with \(d = 6, -1\), this leaves \(\psi = \left(\frac{\cdot}{6}\right)\) or \(\psi = 1\) as the only possibilities.

Assume \(\psi = 1\). Then \(\rho\) factors through \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q}) \twoheadrightarrow \mathbb{S}_3 \subset \text{GL}_2(\mathbb{F}_2)\). But the unique \(\mathbb{S}_3\) extension of \(\mathbb{Q}\) unramified outside of \(\{2, 3\}\) is \(\mathbb{Q}(\zeta_9 + \zeta_9^{-1})\), in which primes \(p \equiv \pm 1 \mod 9\) split completely. Hence

\[
\text{ord} (\rho(\text{Frob}_p)) = \begin{cases} 1 & \text{if } p \equiv \pm 1 \mod 9, \\ 3 & \text{if } p \not\equiv \pm 1 \mod 9. \end{cases}
\]

This contradicts the fact that \(\rho(\text{Frob}_{13})\) has order at most 2.

Therefore \(\psi = \left(\frac{\cdot}{6}\right)\) and \(\rho\) is surjective.

Finally, we know that \(\ker(\rho) = \text{Gal}(\overline{\mathbb{Q}}/K_6)\) for some \(S_3\)-Galois extension \(K_6\) of \(\mathbb{Q}\) containing the quadratic field \(\mathbb{Q}(\sqrt{-6})\) and with discriminant of type \(\pm 2^{12} 3^6\). Such a field \(K_6\) is unique by applying the class field theory to the extension \(K_6\) over \(\mathbb{Q}(\sqrt{-6})\). And it is given as the splitting field of a cubic polynomial over \(\mathbb{Q}\):

\[K_6 = \text{Split}(x^3 + 3x - 2)\]

from the table by Cohen [Coh]. The uniqueness of \(\rho\) then follows from the uniqueness of \(K_6\) and the uniqueness of degree 2 irreducible representation over \(\mathbb{F}_2\) of \(S_3\). \(\square\)

Now we are ready to apply Serre’s theorem to our representations \(\rho'_2\) and \(\tilde{\rho}_2\). It follows from the above proof that the fixed field of the possible deviation group \(\tilde{G}\) contains the field \(K_6\). We start by finding all quartic fields \(M\) with \(\text{Gal}(M/\mathbb{Q}) = S_4\), which contain \(K_6\) and are unramified outside 2 and 3. There are three such fields. Listed below are their defining equations, discriminants, and certain primes \(p\) such that each \(\text{Frob}_p\) is of order 4 in the group \(\text{Gal}(M/\mathbb{Q})\).

<table>
<thead>
<tr>
<th>Defining equation</th>
<th>Discriminant</th>
<th>(p) with order 4 Frobenius</th>
</tr>
</thead>
<tbody>
<tr>
<td>(x^4 - 4x - 3 = 0)</td>
<td>((-216) \cdot 8^2 = -2^9 \cdot 3^3)</td>
<td>13, 17, 19, 23</td>
</tr>
<tr>
<td>(x^4 - 8x + 6 = 0)</td>
<td>((-216) \cdot 16^2 = -2^{13} \cdot 3^3)</td>
<td>13, 17</td>
</tr>
<tr>
<td>(x^4 - 12x^2 - 16x + 12 = 0)</td>
<td>((-216) \cdot 16^2 = -2^{13} \cdot 3^3)</td>
<td>19, 23</td>
</tr>
</tbody>
</table>

As \(H'_S(T)\) agrees with the characteristic polynomial of \(\tilde{\rho}_2(\text{Frob}_p)\) for primes \(p = 13, 19\), the value of \(t\) at such \(\text{Frob}_p\) is zero. Hence we may rule out two possible deviation groups \(\tilde{G} = S_4 \times \{\pm 1\}\) and \(\tilde{G} = S_4\). For the case \(\tilde{G} = S_3 \times \{\pm 1\}\), elements of interest are those primes \(p \neq 2, 3\) such that
The two representations $\rho_2^*$ and $\tilde{\rho}_2$ are isomorphic.

It is worth pointing out that this theorem uses the information of $H_p$ for primes $5 \leq p \leq 19$ only. Further, the $H_p(T)$’s for $p \equiv 1 \pmod{3}$ in Table 1 are uniquely determined, given by the characteristic polynomial of $\tilde{\rho}_2(Frob_p)$ as in the list below.

**Corollary 6.5.** The characteristic polynomials $H_p(T)$ and $H'_p(T)$ of $\rho_2^*(Frob_p)$ over $\mathbb{Q}_2$ and $K = \mathbb{Q}_2(A) = \mathbb{Q}(A)_{1+\phi}$, respectively, for primes $5 \leq p \leq 31$ are listed in (Table 2).

Finally we prove

**Theorem 6.6.** The representations $\rho_2$, $\rho_2^*$, and $\tilde{\rho}_2$ are isomorphic to each other.

**Proof.** We know from Theorem 6.4 that $\rho_2^*$ and $\tilde{\rho}_2$ are isomorphic. Further, by Proposition 5.1, $\rho_2$ is isomorphic to $\rho_2^* \otimes \phi_2$ for some character $\phi_2$ of order at most 2. Therefore it remains to determine $\phi_2$. Since $\rho_2$ and $\rho_2^*$ are unramified outside 2 and 3, the character $\phi_2$, if nontrivial, is associated to a quadratic field $\mathbb{Q}(\sqrt{d})$ with $d = 2, 3, 6, -1, -2, -3, -6$. Write $\phi(p)$ for $\phi_2(Frob_p)$ for brevity. For odd primes $p \geq 5$,
the characteristic polynomial of $\rho_2(\text{Frob}_p)$ is
\[ \tilde{H}_p(T) := T^4 - C_1(p)\phi(p)T^3 + C_2T^2 - p^2C_1(p)\phi(p)T + p^4, \]
where
\[ H_p(T) = T^4 - C_1(p)T^3 + C_2T^2 - p^2C_1(p)T + p^4 \]
is the characteristic polynomial of $\rho_2^*(\text{Frob}_p)$. By Theorem 1.2, the cusp form
\[ f_+(\tau) = \sum_{n \geq 1} a(n)q^{n/15} = q^{1/15} + iq^{2/15} - \frac{11}{3}q^{4/15} - i\frac{16}{3}q^{5/15} - \frac{4}{9}q^{7/15} \]
\[ + i\frac{71}{9}q^{8/15} + \frac{932}{81}q^{10/15} + \cdots, \]
in $S_3(\Gamma)$ satisfies the congruence relation
\[ \text{ord}_p((a(np^2) - C_1(p)\phi(p)a(np) + C_2(p)a(n) - p^2C_1(p)\phi(p)a(n/p) + p^4a(n/p^2))} \]
\[ \geq 2(\text{ord}_pn + 1) \]
for all $n \geq 1$ and $p \geq 5$. Applying this congruence relation to $n = 1$, $p = 7, 13$ and using the explicit values of $C_1(p), C_2(p)$ from Proposition 5.4 as well as the known Fourier coefficients of $f_+$ from Section 4, we find that $\phi(7) = \phi(13) = 1$. Therefore either $\phi_2$ is trivial or $\phi_2$ is $\chi_{-3}$, the quadratic character attached to the field $\mathbb{Q}(\sqrt{-3})$. On the other hand, the two newforms $g_{3i}$ and $g_{-3i}$ are twist of each other by $\chi_{-3}$, which in turn implies that the representation $\tilde{\rho}_2$ is invariant under twisting by $\chi_{-3}$, and hence so is $\rho_2^*$. Therefore in both cases of $\phi_2$ we have $\rho_2$ isomorphic to $\rho_2^*$. □

7. The Atkin–Swinnerton-Dyer congruence relations

Let $\tilde{\xi}$ be the map on the elliptic surface $E$ (given by (8)) sending the base parameter $t$ to $\omega^2 t$, where $\omega = e^{2\pi i/3}$ is a primitive cubic root of one. It induces an action on the weight 3 cusp forms of $\Gamma$ via $\tilde{\xi}(f_1) = \omega f_1, \tilde{\xi}(f_2) = \omega^2 f_2$.

Fix a prime $p \neq 2, 3$. Following the notation in [Sch85], denote by $F$ the canonical endomorphism of $L_k(X, \mathbb{Z}_p)$ coming from an $F$-crystal with logarithmic singularities. In our case, $k = 1$ so that $k + 2 = 3$ is the weight. Recall that $L_k(X, \mathbb{Z}_p)$ is the direct sum of the module $S_3(X, \mathbb{Z}_p)$ of weight 3 cusp forms with Fourier coefficients in $\mathbb{Z}_p$ with its dual $S_3(X, \mathbb{Z}_p)^\vee$. The basis $\{f_1, f_2\}$ of $S_3(X, \mathbb{Z}_p)$ gives rise to a dual basis $\{f_1^\vee, f_2^\vee\}$ of $S_3(X, \mathbb{Z}_p)^\vee$. Denote the space $L_k(X, \mathbb{Z}_p)$ by $V$ for brevity. The two
operators $A$ and $\zeta$ act on $V$ as follows:

\[
\begin{align*}
A(f_1) &= f_2, \quad A(f_2) = -f_1, \quad A(f_1^\vee) = f_2^\vee, \quad A(f_2^\vee) = -f_1^\vee; \\
\zeta(f_1) &= \omega f_1, \quad \zeta(f_2) = \omega^2 f_2, \quad \zeta(f_1^\vee) = \omega^2 f_1^\vee, \quad \zeta(f_2^\vee) = \omega f_2^\vee.
\end{align*}
\]  

(25)

Here the actions of $A$ and $\zeta$ on cusp forms were given before, and their actions on the dual space are obtained by applying the dual action of an operator $T$:

\[
T(h^\vee)(v) = h^\vee(T^{-1} \cdot v) \quad \text{for } h^\vee \in S_3(X, \mathbb{Z}_p)\,^\vee, \quad \text{and } v \in S_3(X, \mathbb{Z}_p).
\]

Denote by $W$ the representation space of the 4-dimensional 2-adic representation $\rho_2$ of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ described at the beginning of Section 5. The operators $A$, $\zeta$ and $F_p$, the Frobenius at $p$, act on $W$.

To prove the Atkin–Swinnerton-Dyer congruence relations at $p$, we shall compare the 2-adic theory and $p$-adic theory. We begin by summarizing the actions of the operators involved.

**Proposition 7.1.** The operators $F_p$, $F$, $A$, $\zeta$ satisfy the following relations on their respective representation spaces:

(i) $AF = FA, \quad F_p A = AF_p$;

(ii) $A \zeta = \zeta^2 A$;

(iii) If $p \equiv 1 \mod 3$, then $\zeta F = F \zeta, \quad \zeta F_p = F_p \zeta$; If $p \equiv 2 \mod 3$, then $\zeta F = F \zeta^2$, $\zeta F_p = F_p \zeta^2$, hence in this case, $\zeta(AF) = (AF)\zeta$, and $\zeta(AF_p) = (AF_p)\zeta$;

(iv) $A^4 = 1, \quad \zeta^3 = 1$.

**Proof.** Assertion (i) follows from the fact that $f_1$ and $f_2$ have rational Fourier coefficients; assertions (ii) and (iv) follow from (25). We need only prove (iii). The action of $\zeta$ on the base curve is given by

\[
\zeta(t) = \omega^2 \cdot t
\]

and $\zeta$ acts trivially on the elliptic curve over $X$. If $p \equiv 1 \mod 3$, then

\[
F_p \zeta(t) = F_p(\omega^2 \cdot t) = (\omega^2)^p \cdot t^p = \omega^2 \cdot t^p = \zeta F_p(t).
\]  

(26)

Therefore on the 2-adic representation space $W$, the action of $F_p$ commutes with $\zeta$ by functoriality. On the $p$-adic representation space $V$, the action of $F$ is semi-linear, hence

\[
F \zeta(t) = F(\omega^2 \cdot t) = (\omega^2)^p \cdot F(t) = \zeta F(t)
\]
and hence the commutativity of $F$ and $\zeta$. In the case $p \equiv 2 \mod 3$, we check the commutativity in the same way as in Eq. (26):

$$F_p \zeta^2(t) = F_p(\omega \cdot t) = (\omega)^p F_p(t) = (\omega)^2 F_p(t) = \zeta F_p(t),$$

$$F \zeta^2(t) = F(\omega \cdot t) = (\omega)^p F(t) = (\omega)^2 F(t) = \zeta F(t),$$

as described in (iii).

Denote by $\text{char}(U, T)$ the characteristic polynomial of an operator $T$ on the space $U$. The representation spaces $W$ and $V$ decompose as direct sums of eigenspaces of the two operators $A$ and $\zeta$, respectively, with eigenvalues appearing as subindices:

$$W = W_i \oplus W_{-i}, \quad V = V_i \oplus V_{-i}; \quad (27a)$$

$$W = W_\omega \oplus W_{\omega^2}, \quad V = V_\omega \oplus V_{\omega^2}. \quad (27b)$$

We have shown in Section 6 that $\rho_2$ is isomorphic to the 2-adic representation $\tilde{\rho}_2$ acting on the 4-dimensional space $\tilde{W}$ attached to the newforms $g_+$ and $g_-$ of weight 3 and level 27. The operator

$$H_{27} = \begin{pmatrix} 0 & -1 \\ 27 & 0 \end{pmatrix}$$

on $\tilde{W}$ plays the same role as $A$ on $W$. Hence the space $\tilde{W}$ decomposes as the sum of two eigenspaces of $H_{27}$:

$$\tilde{W} = \tilde{W}_i \oplus \tilde{W}_{-i}. \quad (28a)$$

Denote by $\text{Frob}_p$ the action of the Frobenius at $p$ on the space $\tilde{W}$. The isomorphism between $\rho_2$ and $\tilde{\rho}_2$ yields immediately

$$\text{char}(W_{\pm i}, F_p) = \text{char}(\tilde{W}_{\pm i}, \text{Frob}_p). \quad (28b)$$

On the other hand, it was shown in [Sch85] that

$$\text{char}(V, F) = \text{char}(W, F_p).$$

Consequently we have

$$\text{char}(V_i, F)\text{char}(V_{-i}, F) = \text{char}(\tilde{W}_i, \text{Frob}_p)\text{char}(\tilde{W}_{-i}, \text{Frob}_p)$$

$$= \text{char}(W_i, F_p)\text{char}(W_{-i}, F_p).$$

We proceed to prove a refinement of this relation.
Theorem 7.2. There hold

\[ \text{char}(V_i, F) = \text{char}(W_i, F_p) = \text{char}(\tilde{W}_i, \text{Frob}_p), \]  
\[ \text{char}(V_{-i}, F) = \text{char}(W_{-i}, F_p) = \text{char}(\tilde{W}_{-i}, \text{Frob}_p). \]  

Proof. In view of the above analysis, it suffices to prove \( \text{char}(V_i, F) = \text{char}(W_i, F_p) \) under the assumption that not all eigenvalues of \( \text{char}(V, F) \) are equal. The eigenspaces of \( V \) occurred in (27) have the following explicit bases:

\[ V_i = (f_1 - if_2, f_1^\vee - if_2^\vee), \quad V_{-i} = (f_1 + if_2, f_1^\vee + if_2^\vee), \]  
\[ V_\omega = (f_1, f_2^\vee), \quad V_\omega^2 = (f_2, f_1^\vee). \]  

The proof is divided into two cases.

Case 1: \( p \equiv 1 \pmod{3} \). In decomposition (27b), the operator \( A \) commutes with \( F \), it maps \( V_\omega \) isomorphically onto \( V_\omega^2 \), and \( V_\omega \) and \( V_\omega^2 \) are \( F \)-invariant by Proposition 7.1. This implies that

\[ \text{char}(V, F) = \text{char}(V_\omega, F) \cdot \text{char}(V_\omega^2, F) = \text{char}(V_\omega, F)^2. \]

Therefore \( \text{char}(V_\omega, F) \) is uniquely determined by \( \text{char}(V, F) \). Let \( \mu \) and \( \nu \) be the two roots of \( \text{char}(V_\omega, F) \); then \( \mu \neq \nu \) by assumption.

As \( F \) commutes with \( A \), the operator \( F \) also fixes \( V_{\pm i} \). We may assume that \( \mu \) is a root of \( \text{char}(V_i, F) \) and call the other root \( \lambda \). Let \( v_1 = a(f_1 - if_2) + b(f_1^\vee - if_2^\vee) \) be an eigenvector of \( F \) on \( V_i \) with eigenvalue \( \mu \). Since \( \mu \) is also an eigenvalue of \( F \) on \( V_\omega \) and \( V_\omega^2 \), using the explicit bases of these two spaces, we may express \( v_1 \) as a sum of two \( \mu \)-eigenvectors \( \alpha f_1 + \beta f_2^\vee \in V_\omega \) and \( \gamma f_2 + \delta f_1^\vee \in V_\omega^2 \):

\[ a(f_1 - if_2) + b(f_1^\vee - if_2^\vee) = (\alpha f_1 + \beta f_2^\vee) + (\gamma f_2 + \delta f_1^\vee). \]

Without loss of generality, we may assume \( a = 1 \). By comparing the coefficients of both sides, we conclude that \( f_1 - if_2^\vee \) and \( f_2 + ibf_1^\vee \) form a basis of the \( \mu \)-eigenspace of \( F \) and \( v_1 = (f_1 - if_2^\vee) - i(f_2 + ibf_1^\vee) \). Note that \( v_2 = (f_1 - ibf_2^\vee) + i(f_2 + ibf_1^\vee) = (f_1 + if_2) - b(f_1^\vee + if_2^\vee) \) is a \( \mu \)-eigenvector of \( F \) on \( V_{-i} \). This shows that the other eigenvalue \( \lambda \) on \( V_i \) is equal to \( \nu \). Therefore

\[ \text{char}(V_i, F) = \text{char}(V_{-i}, F) = \text{char}(V_\omega, F) = \text{char}(V_\omega^2, F). \]

On the 2-adic side, we see this from the space \( \tilde{W} \). More precisely, the space \( \tilde{W}_{\pm i} \) is the space attached to the newform \( g_\pm \), and \( g_- \) is the twist of \( g_+ \) by the quadratic
character $\chi_{-3}$ as noted in Section 6. For $p \equiv 1 \mod 3$, we have $\chi_{-3}(p) = 1$ and hence $\text{char} (\tilde{W}_i, \text{Frob}_p) = \text{char} (\tilde{W}_{-i}, \text{Frob}_p)$, and consequently

$$\text{char} (\tilde{W}_i, \text{Frob}_p) = \text{char} (W_i, F_p) = \text{char} (V_i, F),$$

as desired.

Case 2: $p \equiv 2 \mod 3$. In this case, by Proposition 7.1, neither $F$ nor $F_p$ commutes with the operator $\zeta$, but $FA = AF$ and $F_p A = AF_p$ do. Therefore we use $FA$ and $F_p A$ instead. The same analysis as in Case 1 yields

$$\text{char} (V_i, FA) = \text{char} (V_{-i}, FA) = \text{char} (V_\omega, FA) = \text{char} (V_{\omega^2}, FA).$$

On the 2-adic side, as explained in Case 1, the action of $\text{Frob}_p$ on $W_i$ is twisted by $\chi_{-3}(p) = -1$ of the action of $\text{Frob}_p$ on $W_i$, and hence

$$\text{char} (\tilde{W}_i, \text{Frob}_p H_{27}) = \text{char} (\tilde{W}_{-i}, \text{Frob}_p H_{27})$$

or equivalently,

$$\text{char} (W_i, F_p A) = \text{char} (W_{-i}, F_p A).$$

Applying (4.4.1) of [Sch85] to the matrix $A$, we obtain

$$\text{char} (V, FA) = \text{char} (W, F_p A),$$

hence

$$\text{char} (V_i, FA)^2 = \text{char} (W_i, F_p A)^2$$

from which it follows that

$$\text{char} (V_i, FA) = \text{char} (W_i, F_p A).$$

Since the action of $A$ on both spaces is multiplication by $i$, this yields the desired equality

$$\text{char} (V_i, F) = \text{char} (W_i, F_p). \quad \Box$$

The above theorem allows us to conclude that $f_+$ and $g_+$ (resp. $f_-$ and $g_-$) satisfy the Atkin–Swinnerton-Dyer congruence relations at $p \neq 2, 3$ by the same argument as
Scholl’s proof of Theorem 1.2 presented in Section 5 of [Sch85]. This completes the proof of Theorem 1.4.

**Remark 7.3.** The same congruence relations at \( p \) can be concluded from Remark 5.8 of [Sch85] provided that the action of \( F \) is ordinary.

### 8. The Atkin–Swinnerton-Dyer congruence relations and elliptic modular K3 surfaces over \( \mathbb{Q} \)

In this section we derive similar results for the Atkin–Swinnerton-Dyer congruence relations arising from K3 surfaces over \( \mathbb{Q} \).

First consider an explicit example. Let \( \Gamma_2 \) denote the group associated to the algebraic equation (12). A similar discussion shows that this is a noncongruence subgroup. The space \( S_3(\Gamma_2) \) of weight 3 cusp forms for \( \Gamma_2 \) is 1-dimensional and it is generated by

\[
h_2 = \sqrt{E_1E_2} = q^{1/10} - \frac{3^2}{2} q^{3/10} + \frac{3^3}{2^3} q^{5/10} + \frac{3 \cdot 7^2}{2^4} q^{7/10} - \frac{3^2 \cdot 7 \cdot 19}{2^7} q^{9/10} + O(q^{10/10}),
\]

where \( E_1 \) and \( E_2 \) are given by (17) and (18), respectively. It is clear that the Fourier expansion of \( h_2 \) at the cusp \( \infty \) has coefficients in the ring \( \mathbb{Z}[1/2] \). Let \( g_2 \) be a level 16 newform with the first few terms of its Fourier expansion (provided by Stein’s data base [Stein]) as

\[
g_2 = q - 6q^5 + 9q^9 + 10q^{13} - 30q^{17} + 11q^{25} + 42q^{29} + O(q^{32}).
\]

Using Serre’s method, we can show that the Atkin–Swinnerton-Dyer congruence relations hold for \( h_2 \) and \( g_2 \). In fact, a more general result regarding this situation can be proved.

Recall that a K3 surface \( S \) is a simply connected compact complex surface with trivial canonical bundle. Its Hodge diamond, as the one defined in Section 2, is

\[
\begin{array}{ccc}
1 & 0 & 0 \\
0 & 20 & 1 \\
0 & 0 & 1 \\
1
\end{array}
\]

Hence as a compact complex surface, its Picard number \( \rho(S) \leq h^{1,1} = 20 \). The cohomology group \( A = H^2(S, \mathbb{Z}) \) is a rank 22 free \( \mathbb{Z} \)-module, called a K3 lattice. As a lattice, \( A \) is unimodular due to the Poincaré duality, even by Wu’s formula, with
signature (3, 19) by the Hodge index theorem. Its Néron–Severi group $NS(S)$, defined as in Section 2, is a sublattice of $A$ of signature $(1, \rho(S) - 1)$ again by the Hodge index theorem. An elliptic surface $\pi : S \to C$ is a 2-dimensional complex variety over the base curve $C$ such that every fiber $\pi^{-1}(t)$ is a smooth genus one curve except for finitely many points $t$ in $C$. A compact elliptic surface is called an elliptic modular surface if its monodromy group $\Gamma_S$ is a finite index subgroup of $SL_2(\mathbb{Z})$ and $-I \notin \Gamma_S$ [Shi72].

**Theorem 8.1.** Let $S$ be an elliptic modular K3 surface defined over $\mathbb{Q}$ with $\Gamma_S$ being the associated modular group. Let $f_{\Gamma_S}$ be a nonzero $M$-integral form in the 1-dimensional space $S_3(\Gamma_S)$ for some integer $M$. Then there is a weight 3 cusp form $g_{\Gamma_S}$ with integral Fourier coefficients for some congruence subgroup such that $f_{\Gamma_S}$ and $g_{\Gamma_S}$ satisfy the Atkin–Swinnerton-Dyer congruence relations.

**Proof.** Since $S$ is an elliptic modular K3 surface, by Shioda’s result [Shi72], its Picard number $\rho(S) = 20$. Further, the work of Shioda and Inose [SI77] on K3 surfaces with Picard number 20 shows that the Hasse–Weil $L$-function attached to $S$ contains a factor $L(s, \chi^2)$, where $\chi$ is a Grossencharacter of some imaginary quadratic extension of $\mathbb{Q}$ associated to an elliptic curve over $\mathbb{Q}$ with complex multiplications, arising from the K3 lattice of $S$. Combining the analytic behavior of $L(s, \chi^2)$ proved by Hecke and the converse theorem for $GL_2$ proved by Weil, we know that $L(s, \chi^2)$ is also the $L$-function attached to a weight 3 cuspidal newform $h_{\Gamma_S}$ with integral coefficients for a congruence subgroup.

Let $\rho_l(S)$ be the $l$-adic representation of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ associated to $h_{\Gamma_S}$, which exists for almost all primes $l$. In particular, this representation is isomorphic to $\rho_l^*$ defined in Section 5 induced by the action of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ on $W_l$. This isomorphism of representations can be seen explicitly via the characteristic polynomials of Frobenius elements. Due to a trace formula of Monsky [Mon71], for almost all prime $p$, the Euler $p$-factor $P_{21,p}(p^{-s})$ of the $L$-function of $h_{\Gamma_S}$ appears as one part of the characteristic polynomial of the Frobenius endomorphism $F_p$ acting on the crystalline cohomology $H^2_{cris}(S_p/\mathbb{Z}_p)$, where $S_p$ is the Néron minimal model of $S$ over $\mathbb{F}_p$. In particular, $P_{21,p}(p^{-s})$ corresponds to the characteristic polynomial of $\overline{\mathbb{F}}_p$ acting on the orthogonal complement of the Néron–Severi group $NS(S_p \otimes \overline{\mathbb{F}}_p) \otimes \mathbb{Q}$ in $H^2_{cris}(S_p/\mathbb{Z}_p) \otimes \mathbb{Q}$, where $\overline{\mathbb{F}}_p$ is the algebraic closure of $\mathbb{F}_p$ (for details see [SB85, Section 12]). By the Lefschetz fixed point theorem, the coefficients of $P_{21,p}(p^{-s})$ can be calculated by the same trace formula [Sch88, Section 3] for the Frobenius endomorphism $\text{Frob}_p$ on the $l$-adic cohomology $W_l$.

Denote by $\rho_l(\Gamma_S)$ the $l$-adic representation of the Galois group of $\mathbb{Q}$ associated to the space $S_3(\Gamma_S)$ constructed by Scholl [Sch85]. By Proposition 5.1, $\rho_l(\Gamma_S)$ and $\rho_l(S)$ are isomorphic up to a quadratic twist. Let $\phi$ be a character of $\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})$ of order at most two such that $\rho_l(\Gamma_S)$ is isomorphic to $\rho_l(S) \otimes \phi$. By class field theory, we may regard $\phi$ as a Dirichlet character of order at most two. Let $g_{\Gamma_S}$ be the newform having the same eigenvalues as the twist of $h_{\Gamma_S}$ by $\phi$. Then $g_{\Gamma_S}$ also has integral Fourier coefficients and $\rho_l(\Gamma_S)$ is isomorphic to the $l$-adic representation attached to
Applying Theorem 1.2 to \( \rho_l(\Gamma_S) \) and \( f\Gamma_S \), we conclude that \( f\Gamma_S \) and \( g\Gamma_S \) satisfy the Atkin–Swinnerton-Dyer relations. □

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**References**


