# Infinite transitivity on the Calogero-Moser space $\mathcal{C}_{2}{ }^{*}$ 

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Abstract. We prove a particular case of the conjecture of Berest-Eshmatov-Eshmatov by showing that the group of unimodular automorphisms of $\mathbb{C}[x, y]$ acts in an infinitely-transitive way on the Calogero-Moser space $\mathcal{C}_{2}$.

## 1. Introduction

Let $M_{n}$ be the $\mathbb{C}$-algebra of $n \times n$ matrices over $\mathbb{C}$. The group $\mathrm{GL}_{n}(\mathbb{C})$ acts on the direct product $M_{n} \times M_{n}$ in the natural way:

$$
\begin{equation*}
g \cdot(X, Y)=\left(g X g^{-1}, g Y g^{-1}\right), \quad g \in \mathrm{GL}_{n}(\mathbb{C}) \tag{1}
\end{equation*}
$$

For an integer $n \geqslant 0$, let $\hat{\mathcal{C}}_{n}$ be the subset of $M_{n} \times M_{n}$ defined as

$$
\left\{(X, Y) \in M_{n} \times M_{n}: \operatorname{rank}\left([X, Y]+I_{n}\right)=1\right\}
$$

where $I_{n}$ is the $n \times n$ identity matrix. The action of (1) on $M_{n} \times M_{n}$ restricts to an action on $\hat{\mathcal{C}}_{n}$, and we can then define the $n$-th CalogeroMoser space $\mathcal{C}_{n}$ to be the quotient $\hat{\mathcal{C}}_{n} / / \mathrm{GL}_{n}$. These spaces were studied in detail by Wilson [4], where it was shown, among other things, that $\mathcal{C}_{n}$ is a smooth, affine, irreducible, complex, symplectic variety of dimension $2 n$.

[^0]The group of unimodular automorphisms of $\mathbb{C}[x, y]$ acts on $\mathcal{C}_{n}$, and it is proved in [1] that this action is doubly transitive. Additionally, a conjecture that this action is infinitely transitive is stated. Recently, in [3], this conjecture was proved.

The goal of this paper is to give another proof of infinite transitivity for the case $n=2$. The proofs here are more constructive and shed more light on the action on $\mathcal{C}_{2}$. We do this inductively by first choosing distinct points $x_{1}, \ldots, x_{n}, x_{n+1} \in \mathcal{C}_{2}$. Then, for any tuple of distinct elements $\left(y_{1}, \ldots, y_{n}, y_{n+1}\right)$, we use the inductive hypothesis to move $\left(y_{1}, \ldots, y_{n}, y_{n+1}\right)$ to $\left(x_{1}, \ldots, x_{n}, \tilde{y}_{n+1}\right)$. If we can then find elements of $G$ that stabilize $x_{1}, \ldots, x_{n}$ while acting transitively on the rest of the elements of $\mathcal{C}_{2}$, we can then move the element $\tilde{y}_{n+1}$ to the predetermined element $x_{n+1}$, while keeping $x_{1}, \ldots, x_{n}$ fixed. This will show that any tuple $\left(y_{1}, \ldots, y_{n+1}\right)$ is in the same orbit as $\left(x_{1}, \ldots, x_{n+1}\right)$, thus establishing $(n+1)$-transitivity. For this approach to work, we see that we will require information about the stabilizers of specific elements in $\mathcal{C}_{2}$, which we collect in future sections.

In general, an explicit representation for the coordinate ring, $\mathbb{C}\left[\mathcal{C}_{n}\right]$, of a Calogero-Moser space is not known. However, for $n=2$, it is not difficult to find. Let $A=X-\frac{1}{2} \operatorname{Tr}(X) I_{2}$ and $B=Y-\frac{1}{2} \operatorname{Tr}(Y) I_{2}$ be traceless matrices associated to $X$ and $Y$, respectively. In this case, using the generators $\left\{\operatorname{Tr}(X), \operatorname{Tr}(Y), \operatorname{Tr}\left(X^{2}\right), \operatorname{Tr}(X Y), \operatorname{Tr}\left(Y^{2}\right)\right\}$ of $\mathbb{C}\left[\left(M_{2} \times M_{2}\right) / / \mathrm{GL}_{2}\right]$ found in [2], we define the following generators of $\mathbb{C}\left[\mathcal{C}_{2}\right]$ :

$$
a_{1}=\operatorname{Tr}(X), a_{2}=\operatorname{Tr}(Y), a_{3}=\operatorname{Tr}\left(A^{2}\right), a_{4}=\operatorname{Tr}(A B), a_{5}=\operatorname{Tr}\left(B^{2}\right)
$$

Using the fact that a non-zero $2 \times 2$ matrix is of rank one if and only if its determinant is zero, we find that

$$
\mathbb{C}\left[\mathcal{C}_{2}\right]=\mathbb{C}\left[a_{1}, a_{2}, a_{3}, a_{4}, a_{5}\right] /\left(a_{4}^{2}-a_{3} a_{5}-1\right)
$$

Note that there is a one to one correspondence between a point $(X, Y) \in \mathcal{C}_{2}$ and a point $\left(a_{1}, \ldots, a_{5}\right) \in \mathbb{C}^{5}$ such that $a_{4}^{2}-a_{3} a_{5}=1$, given by

$$
\begin{equation*}
(X, Y) \mapsto\left(\operatorname{Tr}(X), \operatorname{Tr}(Y), \operatorname{Tr}\left(A^{2}\right), \operatorname{Tr}(A B), \operatorname{Tr}\left(B^{2}\right)\right) \tag{2}
\end{equation*}
$$

## 2. Preliminaries

Denote by $G$ the group generated by the following two kinds of automorphisms of $M_{n} \times M_{n}$ :
(i) $\Phi_{p}:(X, Y) \mapsto(X, Y+p(X))$, where $p \in \mathbb{C}[t]$,
(ii) $\Psi_{q}:(X, Y) \mapsto(X+q(Y), Y)$, where $q \in \mathbb{C}[t]$.

It is known [5], G is isomorphic to

$$
\operatorname{SAut}(C[x, y])=\left\{f=\left(f_{1}, f_{2}\right) \in \operatorname{Aut}(C[x, y]) \mid \operatorname{Jac}\left(f_{1}, f_{2}\right)=1\right\}
$$

where $\operatorname{Jac}\left(f_{1}, f_{2}\right)$ is the determinant of the Jacobian matrix of the map $\left(f_{1}, f_{2}\right)$. Note that $\operatorname{Aut}(\mathbb{C}[x, y])$ is isomorphic to a semidirect product $\operatorname{SAut}(\mathbb{C}[x, y]) \rtimes G_{m}$, where $G_{m}$ is a multiplicative group of the field $\mathbb{C}$ which acts on $\mathbb{C}[x, y]$ by scalar multiplication on variables $x$ and $y$. From the correspondence given in (2), we obtain an easy way of computing the action of the above group, $G$, using the following component-wise rules:

$$
\begin{align*}
& \Phi_{p}\left(a_{1}\right):=a_{1} \\
& \Phi_{p}\left(a_{2}\right):=a_{2}+\operatorname{Tr}(p(X)) \\
& \Phi_{p}\left(a_{3}\right):=a_{3}  \tag{3}\\
& \Phi_{p}\left(a_{4}\right):=a_{4}+\operatorname{Tr}(A p(X)) \\
& \Phi_{p}\left(a_{5}\right):=a_{5}+\operatorname{Tr}\left(p^{2}(X)\right)+2 \operatorname{Tr}(B \cdot p(X))-\frac{1}{3} \operatorname{Tr}^{2}(p(X)) .
\end{align*}
$$

The action of $\Psi_{q}$ on $\mathcal{C}_{2}$ is similar, and is symmetric to (3).
For a matrix $M=\left(\begin{array}{cc}\alpha & \beta \\ \lambda & \mu\end{array}\right) \in \mathrm{SL}_{2}$ (so that $\alpha \mu-\beta \lambda=1$ ) consider $\Theta_{M}: \mathcal{C}_{2} \rightarrow \mathcal{C}_{2}$ defined by

$$
(X, Y) \mapsto(\alpha X+\beta Y, \lambda X+\mu Y)
$$

One can easily find that the action $\Theta_{M}$ is a composition of the automorphisms of type (i) and (ii) using some linear polynomials $p$ and $q$. Under this action, a point will change as follows:

$$
\begin{align*}
& \Theta_{M}\left(a_{1}\right)=\alpha a_{1}+\beta a_{2} \\
& \Theta_{M}\left(a_{2}\right)=\lambda a_{1}+\mu a_{2} \\
& \Theta_{M}\left(a_{3}\right)=\alpha^{2} a_{3}+2 \alpha \beta a_{4}+\beta^{2} a_{5}  \tag{4}\\
& \Theta_{M}\left(a_{4}\right)=\alpha \lambda a_{3}+(\alpha \mu+\beta \lambda) a_{4}+\beta \mu a_{5} \\
& \Theta_{M}\left(a_{5}\right)=\lambda^{2} a_{3}+2 \lambda \mu a_{4}+\mu^{2} a_{5}
\end{align*}
$$

We now remind the following definitions concerning group actions on sets. To do this, let $\mathcal{G}$ be a group acting on a set $S$.

Definition 1. We say the group $\mathcal{G}$ acts transitively, or that the action is transitive, if for every pair of elements $s, r \in S$, there is a $g \in \mathcal{G}$ such that $g \cdot s=r$.

Definition 2. The group $\mathcal{G}$ acts $n$-transitively, or the action is $n$-transitive, if it can map any $n$-tuple of distinct points of the set to any other $n$-tuple of distinct points. In other words, if $\left(s_{1}, \ldots, s_{n}\right)$ and $\left(r_{1}, \ldots, r_{n}\right)$ are $n$ tuples of distinct elements in $S$, then there is some $g \in \mathcal{G}$ such that $g \cdot\left(s_{1}, \ldots, s_{n}\right)=\left(g \cdot s_{1}, \ldots, g \cdot s_{n}\right)=\left(r_{1}, \ldots, r_{n}\right)$.

Definition 3. Lastly, we say that the action of $\mathcal{G}$ on $S$ is infinitely transitive if it is $n$-transitive for every positive integer $n$.

We now claim that the action of $G$ on $\mathcal{C}_{2}$ defined above is infinitelytransitive.

## 3. Base cases: $n=1,2,3$

As stated previously, we plan to prove this main result by induction, and so we begin by proving the base cases for $n=1,2,3$. We start with $n=1$ :

Proposition 1. The action of $G$ on $\mathcal{C}_{2}$ is a transitive group action.
Proof. Let $A=\left(a_{1}, a_{2}, a_{3}, a_{4}, a_{5}\right) \in \mathcal{C}_{2}$ be an arbitrary point. Note that this proof does not require us to stabilize any elements, and so we may use $p(t)=-\frac{a_{2}}{2}$ and $q(t)=-\frac{a_{1}}{2}$, to get that

$$
\left(\Psi_{q} \circ \Phi_{p}\right)(A)=\left(0,0, a_{3}, a_{4}, a_{5}\right)
$$

From here we use the action of $\Theta_{M}$, defined in (4), with either the matrix $M_{+}:=\left[\begin{array}{cc}-\frac{a_{5}}{2\left(a_{4}+1\right)} & \frac{1}{2} \\ a_{4}+1 & a_{3}\end{array}\right]$ or the matrix $M_{-}:=\left[\begin{array}{cc}-\frac{a_{5}}{2\left(a_{4}-1\right)} & \frac{1}{2} \\ a_{4}-1 & a_{3}\end{array}\right]$ to reach the point

$$
\Theta_{M}\left(0,0, a_{3}, a_{4}, a_{5}\right)=(0,0,0,1,0)
$$

More specifically, if $a_{3}=0$ or $a_{5}=0$, then, since $a_{4}^{2}-a_{3} a_{5}=1$, we must have that $a_{4}= \pm 1$. If $a_{4}=1$, then we use the matrix $M_{+}$. If $a_{4}=-1$, we use the matrix $M_{-}$. If $a_{3} a_{5} \neq 0$, then either matrix $M_{+}$or $M_{-}$will suffice. Thus we have that all elements $A \in \mathcal{C}_{2}$ are in the orbit of the point $(0,0,0,1,0) \in \mathcal{C}_{2}$.

Next, we prove 2- and 3-transitivity, since they differ from the general $n$ case by requiring us only to focus on stabilizing nilpotent points. We will need the following two lemmas:

Lemma 1. Let $A \in \mathcal{C}_{2} \backslash\{(0,0,0, \pm 1,0)\}$. Then there is a $g \in$ $\operatorname{Stab}\{(0,0,0, \pm 1,0)\}$ such that $A^{\prime}=g A$ satisfies $a_{1}^{\prime} a_{3}^{\prime} \neq 0$, where $A^{\prime}=$ $\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}, a_{4}^{\prime}, a_{5}^{\prime}\right)$.

Proof. We may assume at least one of $a_{1}, a_{2}, a_{3}, a_{5}$ is nonzero. We will proceed by case work.

Case 1: $a_{1} \neq 0$. If $a_{3} \neq 0$ we are already done, so suppose $a_{3}=0$. Without loss of generality, we may assume that $a_{2} \neq 0$, since if $a_{2}=0$, we may apply $\Phi_{t^{2}}$ to arrive at the point $\left(b_{1}, b_{2}, b_{3}, b_{4}, b_{5}\right)$ with $b_{2}=\frac{a_{1}^{2}}{2} \neq 0$. Additionally, since $a_{3}=0$, we must have $a_{4}$ is non-zero, so that there exists an $\alpha \in \mathbb{C}$ such that $\alpha a_{2}\left(\alpha a_{2} a_{5}+2 a_{4}\right) \neq 0$ and such that $\alpha\left(a_{5}+\frac{a_{2}^{2}}{2}\right)+a_{1} \neq 0$, since this is a non-zero polynomial in $\alpha$. We now apply $\Psi_{\alpha t^{2}}$ to arrive at $A^{\prime}$ where $a_{1}^{\prime} a_{3}^{\prime} \neq 0$.

Case 2: $a_{1}=0$. From here, we will show that we can move into case one.

Case 2.1: $a_{5}+\frac{a_{2}^{2}}{2} \neq 0$. We can calculate explicitly that applying $\Psi_{t^{2}}$ gives $a_{1}^{\prime}=a_{5}+\frac{a_{2}^{2}}{2} \neq 0$, so that we are back in Case 1 .

Case 2.2: $a_{5}+\frac{a_{2}^{2}}{2}=0$.
Case 2.2.1: $a_{3} \neq 0$. We can map $a_{2}$ and $a_{5}$ to $a_{2}^{\prime}$ and $a_{5}^{\prime}$ such that $a_{5}^{\prime}+\frac{\left(a_{2}^{\prime}\right)^{2}}{2} \neq 0$ by $\Phi_{\beta t^{2}}$, since $a_{2}^{\prime}=a_{2}+\beta a_{3}$ and $a_{5}^{\prime}=a_{5}$. This moves us back to Case 2.1.

Case 2.2.2: $a_{3}=0$. Since $a_{1}=0$ and $a_{5}+\frac{a_{2}^{2}}{2}=0$ with $a_{2} \neq 0$ after we are at the point $\left(0, a_{2}, 0, \pm 1, \frac{-a_{2}^{2}}{2}\right)$. By applying $\Psi_{t^{3}}$, we can send $a_{1}$ to $\frac{-a_{2}^{3}}{2}$, thus showing that we can send $a_{1}$ to a nonzero value, returning us to Case 1 and completing the proof. It is easy to check that all of the elements of $G$ used above are indeed in $\operatorname{Stab}\{(0,0,0, \pm 1,0)\}$.

Lemma 2. Let $A \in \mathcal{C}_{2}$ with $a_{1} a_{3} \neq 0$. Then there is a $g \in$ $\operatorname{Stab}\{(0,0,0, \pm 1,0)\}$ such that $A^{\prime}=g A$ satisfies $a_{1}^{\prime} a_{3}^{\prime} \neq 0$ and $a_{3}^{\prime} \neq \frac{a_{1}^{2}}{2}$.

Proof. Let $A=\left(a_{1}, a_{2}, a_{3}, a_{4}, a_{5}\right)$ with $a_{1} a_{3} \neq 0$ be given. We also assume $a_{3}=\frac{a_{1}^{2}}{2}$, since otherwise we are done.

Case 1: $a_{2} a_{5} \neq 0$. Then we can apply $\Psi_{\alpha t^{2}}$ to get

$$
\begin{aligned}
\Psi_{\alpha t^{2}}\left(a_{1}, a_{2}, \frac{a_{1}^{2}}{2}, a_{4}, a_{5}\right)= & \left(a_{1}+\left(\frac{a_{2}}{2}+a_{5}\right) \alpha, a_{2}, \frac{a_{1}^{2}}{2}+2 a_{2} a_{4} \alpha+a_{2}^{2} a_{5} \alpha^{2},\right. \\
& \left.a_{4}+a_{2} a_{5} \alpha, a_{5}\right)
\end{aligned}
$$

From this we can see that $a_{1}^{\prime}, a_{3}^{\prime}$ are non-zero polynomials of $\alpha$, so that there are at most finitely many values of $\alpha$ such that $a_{1}^{\prime} a_{3}^{\prime}=0$. Additionally, plugging into $a_{1}^{\prime 2}-2 a_{3}^{\prime}$, we obtain the polynomial

$$
\left(-4 a_{2} a_{4}+2 a_{1}\left(a_{2}^{2} / 2+a_{5}\right)\right) \alpha+\left(-2 a_{2}^{2} a_{5}+\left(a_{2}^{2} / 2+a_{5}\right)^{2}\right) \alpha^{2}
$$

We claim this is a non-zero polynomial in $\alpha$. To see this, assume that it is the zero polynomial, so that $-4 a_{2} a_{4}+2 a_{1}\left(a_{2}^{2} / 2+a_{5}\right)=0$ and $-2 a_{2}^{2} a_{5}+\left(a_{2}^{2} / 2+a_{5}\right)^{2}=0$. This implies that $a_{5}=\frac{-a_{1} a_{2}^{2}+4 a_{2} a_{4}}{2 a_{1}}$ and $a_{5}=\frac{a_{2}^{2}}{2}$. Setting these equal means that we must have $a_{4}=\frac{a_{1} a_{2}}{2}$. (Note that this is where we have used the assumption that $a_{2} a_{5} \neq 0$, so that we can actually solve for $a_{4}$ in this way.) It leads to a contradiction. Thus we have that $a_{1}^{\prime 2}-2 a_{3}^{\prime}$ is a non-zero polynomial in $\alpha$, and hence, since there are uncountably many $\alpha \in \mathbb{C}$, we can choose an $\alpha$ such that $a_{1}^{\prime} a_{3}^{\prime} \neq 0$ and $a_{1}^{\prime 2}-2 a_{3}^{\prime} \neq 0$, as desired.

Case 2: $a_{2} a_{5}=0$. We again assume that $a_{3}=\frac{a_{1}^{2}}{2}$. Applying $\Phi_{\alpha t^{2}}$, we get
$\Phi_{\alpha t^{2}}\left(a_{1}, a_{2}, \frac{a_{1}^{2}}{2}, a_{4}, a_{5}\right)=\left(a_{1}, a_{2}+a_{1}^{2} \alpha, \frac{a_{1}^{2}}{2}, a_{4}+\frac{a_{1}^{3} \alpha}{2}, a_{5}+2 a_{1} a_{4} \alpha+\frac{a_{1}^{4} \alpha^{2}}{2}\right)$
From this, since $a_{1} \neq 0$, we can conclude that $a_{2}^{\prime}$ and $a_{5}^{\prime}$ are non-zero polynomials in $\alpha$, and hence we can choose some $\alpha \in \mathbb{C}$ such that $a_{2}^{\prime} a_{5}^{\prime} \neq 0$, landing us back in the Case 1. (Again, it is easy to check that all elements of $G$ used in the above proof are elements of $\operatorname{Stab}\{(0,0,0, \pm 1,0)\}$.)

We are now able to prove 2 - and 3 - transitivity:
Proposition 2. The action of $G$ on $\mathcal{C}_{2}$ is a 2-transitive group action.
Proof. Let $(A, B) \in \mathcal{C}_{2} \times \mathcal{C}_{2}$ be a pair of distinct points. Then, since $G$ acts on $\mathcal{C}_{2}$ transitively by Proposition 1 , there is a $g \in G$ such that $g(A, B)=\left((0,0,0,1,0), B^{\prime}\right)$ for some $B^{\prime} \in \mathcal{C}_{2}$. Thus, if there is an $h \in$ $\operatorname{Stab}\{(0,0,0,1,0)\}$ such that $h B^{\prime}=(0,0,0,-1,0)$, we are done. In particular, this reduces the problem to showing that for any $A \in \mathcal{C}_{2} \backslash\{(0,0,0,1,0)\}$, there is a $g \in \operatorname{Stab}\{(0,0,0,1,0)\}$ such that $g A=(0,0,0,-1,0)$.

Thus, let $A \in \mathcal{C}_{2} \backslash\{(0,0,0,1,0)\}$ be an arbitrary point. If $A=(0,0,0,-1,0)$, then we are already at the point we desire. Otherwise, we have that $A \in \mathcal{C}_{2} \backslash\{(0,0,0, \pm 1,0)\}$. Using Lemmas 1 and 2, we may also assume that $a_{1} a_{3} \neq 0$ and $a_{1}^{2}-2 a_{3} \neq 0$. Applying $\Phi_{\alpha t^{3}} \in$ $\operatorname{Stab}\{(0,0,0,1,0)\}$ to $A$, we reach the point

$$
\begin{aligned}
& A^{\prime}=\left(a_{1}, a_{2}+\left(a_{1} a_{3}+\frac{1}{2} a_{1}\left(\frac{a_{1}^{2}}{2}+a_{3}\right)\right) \alpha, a_{3}\right. \\
& a_{4}+\left(\frac{a_{1}^{2} a_{3}}{2}+\frac{1}{2} a_{3}\left(\frac{a_{1}^{2}}{2}+a_{3}\right)\right) \alpha \\
&\left.a_{5}+\left(\frac{3 a_{1}^{2}}{2}+a_{3}\right) a_{4} \alpha+\frac{1}{4} a_{3}\left(\frac{3 a_{1}^{2}}{2}+a_{3}\right)^{2} \alpha^{2}\right)
\end{aligned}
$$

Then we can calculate that

$$
\begin{aligned}
& -4 a_{1}^{\prime 3} a_{2}^{\prime} a_{4}^{\prime}-8 a_{1}^{\prime} a_{2}^{\prime} a_{3}^{\prime} a_{4}^{\prime}+a_{1}^{\prime 4} a_{5}^{\prime}+4 a_{3}^{\prime 2} a_{5}^{\prime}+4 a_{1}^{\prime 2} a_{3}^{\prime}\left(a_{2}^{\prime 2}+a_{5}^{\prime}\right) \\
& =4 a_{1}^{2} a_{2}^{2} a_{3}-4 a_{1}^{3} a_{2} a_{4}-8 a_{1} a_{2} a_{3} a_{4}+a_{1}^{4} a_{5}+4 a_{1}^{2} a_{3} a_{5}+4 a_{3}^{2} a_{5} \\
& \\
& \quad+\left(-a_{1}^{5} a_{2} a_{3}+4 a_{1}^{3} a_{2} a_{3}^{2}-4 a_{1} a_{2} a_{3}^{3}+\frac{a_{1}^{6} a_{4}}{2}-a_{1}^{4} a_{3} a_{4}\right. \\
& \\
& \left.\quad-2 a_{1}^{2} a_{3}^{2} a_{4}+4 a_{3}^{3} a_{4}\right) \alpha+\left(\frac{a_{1}^{8} a_{3}}{16}-\frac{a_{1}^{6} a_{3}^{2}}{2}+\frac{3 a_{1}^{4} a_{3}^{3}}{2}-2 a_{1}^{2} a_{3}^{4}+a_{3}^{5}\right) \alpha^{2} .
\end{aligned}
$$

One can check that if the coefficient of the $\alpha^{2}$ term is zero, then either $a_{3}=0$ or $a_{3}=a_{1}^{2} / 2$. Since we know neither of these is true, it follows that the last term is non-zero, and hence this is a non-zero polynomial in $\alpha$. Thus we can choose a $\alpha \in \mathbb{C}$ such that this polynomial does not vanish. Since this is the case, we can consider the polynomials

$$
p(t)=\frac{-a_{4}^{\prime}+1}{a_{1}^{\prime} a_{3}^{\prime}} t^{2}=\frac{-a_{4}^{\prime}+1}{a_{1} a_{3}} t^{2}
$$

and

$$
q(t)=\frac{a_{1}^{\prime} a_{3}^{\prime}\left(2 a_{1}^{\prime} a_{2}^{\prime} a_{3}^{\prime}-a_{1}^{\prime 2}\left(a_{4}^{\prime}-1\right)-2 a_{3}^{\prime}\left(a_{4}^{\prime}-1\right)\right)}{-4 a_{1}^{\prime 3} a_{2}^{\prime} a_{4}^{\prime}-8 a_{1}^{\prime} a_{2}^{\prime} a_{3}^{\prime} a_{4}^{\prime}+a_{1}^{\prime 4} a_{5}^{\prime}+4 a_{3}^{\prime 2} a_{5}^{\prime}+4 a_{1}^{\prime 2} a_{3}^{\prime}\left(a_{2}^{\prime 2}+a_{5}^{\prime}\right)} t^{2}
$$

These satisfy $\Phi_{p(t)}, \Psi_{q(t)} \in \operatorname{Stab}\{(0,0,0,1,0)\}$, and we can calculate that

$$
\left(\Psi_{q} \circ \Phi_{p}\right)\left(a_{1}^{\prime}, a_{2}^{\prime}, a_{3}^{\prime}, a_{4}^{\prime}, a_{5}^{\prime}\right)=\left(b_{1}, b_{2}, 0,-1,0\right)
$$

for some $b_{1}, b_{2} \in \mathbb{C}$.
Case 1: $b_{1} b_{2} \neq 0$. Using $\lambda=\frac{b_{1}}{b_{2}}$ and $\mu=\frac{b_{2}^{2}}{4 b_{1}}$, we can apply the following composition to get

$$
\left(\Psi_{\lambda t} \circ \Phi_{\mu t^{2}+\lambda \mu} \circ \Psi_{-\lambda t}\right)\left(b_{1}, b_{2}, 0,-1,0\right)=(0,0,0,-1,0)
$$

as desired. One can easily check that this composition is in $\operatorname{Stab}\{(0,0,0,1,0)\}$.

Case 2: $b_{1} b_{2}=0$. If $b_{1}=b_{2}=0$, then we already have that $\left(b_{1}, b_{2}, 0,-1,0\right)=(0,0,0,-1,0)$. If not, we have that either $b_{1} \neq 0$ or $b_{2} \neq 0$. In these cases we use the element $\Phi_{t^{2}-\frac{4}{3 a_{1}^{\prime}} t^{3}}$ or $\Psi_{t^{2}-\frac{4}{3 a_{2}^{\prime}}}$, respectively, in order to map

$$
\left(b_{1}, 0,0,-1,0\right) \mapsto\left(b_{1}, b_{1}^{2} / 6,0,-1,0\right)
$$

or

$$
\left(0, b_{2}, 0,-1,0\right) \mapsto\left(b_{2}^{2} / 6, b_{2}, 0,-1,0\right)
$$

landing us back in Case 1.
Thus we have shown that all points in $\mathcal{C}_{2} \backslash(0,0,0,1,0)$ are in the same orbit as $(0,0,0,-1,0)$ under the action of $\operatorname{Stab}\{(0,0,0,1,0)\}$, so that $G$ acts 2 -transitively on $\mathcal{C}_{2}$.

Proposition 3. The action of $G$ on $\mathcal{C}_{2}$ is a 3-transitive group action.
Proof. Since the group $G$ acts 2 -transitively by Proposition 2, we can reduce the problem to showing that for any $A \in C_{2} \backslash\{(0,0,0, \pm 1,0)\}$, there is some $g \in S=\operatorname{Stab}\{(0,0,0, \pm 1,0)\}$ such that $g(A)=(0,0,0,-1,2)$. A straightforward computation shows that all of the elements of $\operatorname{Stab}\{(0,0,0,1,0)\}$ used in the proof of Proposition (2) also stabilize the point $(0,0,0,-1,0)$, so that, using Lemmas 1 and 2 , and then proceeding analogously to the proof of Proposition 2, we are able to find some $g \in S$ such that $g(A)=\left(b_{1}, b_{2}, 0,-1,0\right)$.

We may assume that $b_{1}$ and $b_{2}$ are not simultaneously zero.
Case 1: $b_{1} \neq 0$. For the polynomials $p_{1}(t)=\frac{4\left(b_{1}-4 b_{2}\right)}{b_{1}^{3}} t^{3}-\frac{8\left(b_{1}-3 b_{2}\right)}{b_{1}^{4}} t^{4}$ and $q_{1}(t)=\frac{-b_{1}}{2} t^{2}$, we get that $\left(\Psi_{q_{1}} \circ \Phi_{p_{1}}\right) \in \operatorname{Stab}\{(0,0,0, \pm 1,0)\}$ and

$$
\left(\Psi_{q_{1}} \circ \Phi_{p_{1}}\right)\left(b_{1}, b_{2}, 0,-1,0\right)=(0,0,0,-1,2)
$$

as desired.
Case 2: $b_{2} \neq 0$. Similarly, using the polynomials $p_{2}(t)=\frac{-b_{2}}{2} t^{2}$ and $q_{2}(t)=\frac{4\left(b_{2}-4 b_{1}\right)}{b_{2}^{3}} t^{3}-\frac{8\left(b_{2}-3 b_{1}\right)}{b_{2}^{4}} t^{4}$, we obtain that $\left(\Phi_{p_{2}} \circ \Psi_{q_{2}}\right) \in$ $\operatorname{Stab}\{(0,0,0, \pm 1,0)\}$ and

$$
\left(\Phi_{p_{2}} \circ \Psi_{q_{2}}\right)\left(b_{1}, b_{2}, 0,-1,0\right)=(0,0,2,-1,0)
$$

Then, we can apply $\left(\Psi_{t^{3}} \circ \Phi_{\frac{1}{3} t^{2}}\right) \in \operatorname{Stab}(0,0,0, \pm 1,0)$ to get to the point

$$
\left(\Psi_{t^{3}} \circ \Phi_{\frac{1}{3} t^{2}}\right)(0,0,2,-1,0)=(2 / 3,2,0,-1,0)
$$

landing us back in the case where $b_{1} \neq 0$.

## 4. Stabilizer elements

While proving the base cases, we could easily check that the elements of $G$ being used were in the desired stabilizers; unfortunately it is not as easy to do this as the sets of points we wish to stabilize get larger. This section is concerned with determining which elements of $G$ are in the stabilizers of larger subsets of $\mathcal{C}_{2}$.

Proposition 4. Let $A=\left(0,0,0, \pm 1, a_{5}\right)$ be a point in $\mathcal{C}_{2}$. Then $\Phi_{p}$ stabilizes $A$ if and only if $t^{2} \mid p(t)$.

Proof. Let $(X, Y)$ be the pair of matrices in $\mathcal{C}_{2}$ that corresponds to the point $A$, and recall that the action of $\Phi_{p}$ on $A$ corresponds to the action on $(X, Y)$ defined by $\Phi_{p}(X, Y)=(X, Y+p(X))$ in $\mathcal{C}_{2}$.

Assume $t^{2} \mid p(t)$ and let us show that $\Phi_{p}(A)=A$. Since $\Phi_{p} \circ \Phi_{q}=\Phi_{p+q}$, it is enough to show that for any monomial $\alpha t^{n}$ with $n \geqslant 2$ and $\alpha \in \mathbb{C}$, we have that $\Phi_{\alpha t^{n}}$ stabilizes $A$. Since $a_{1}=a_{3}=0$, we can conclude that $X$ is nilpotent, meaning that $X^{n}=0$ for any $n \geqslant 2$. It follows that $p(X)=\alpha X^{n} \equiv 0$, so that $\Phi_{p}(X, Y)=(X, Y+p(X))=(X, Y)$, and hence $\Phi_{p}(A)=A$.

Conversely, suppose that $t^{2} \nmid p(t)$. This implies that $p$ has non-zero linear or constant terms. Also, since $a_{1}=a_{3}=0, X$ is still nilpotent, so that we can assume $p(t)=\alpha t+\beta$ where one of the parameters $\alpha, \beta$ is nonzero. Then, using (3), we have that

$$
\Phi_{p}(A)=\left(0,2 \beta, 0, \pm 1, a_{5} \pm 2 \alpha\right)
$$

so that $\Phi_{p}$ does not fix $A$, thus proving the contrapositive.
Now, consider the point $A=\left(0,0,0,-1,2^{k}\right)$ for some $k \in \mathbb{Z}_{+}$. We want to determine which $q \in \mathbb{C}[t]$ will satisfy $\Psi_{q}(A)=A$. Again, since $a_{1}=a_{3}=0, X$ is nilpotent, and hence we may assume that $X=\left(\begin{array}{ll}0 & 1 \\ 0 & 0\end{array}\right)$. Also, $\operatorname{Tr}(Y)=0$, so that we can write $Y=\left(\begin{array}{cc}a & b \\ c & -a\end{array}\right)$.

Next, we note that $\operatorname{Tr}(X Y)=-1$, which implies that $c=-1$. Now we note that under the group action of $\mathrm{GL}_{2}$ by the matrix $M=\left(\begin{array}{cc}1 & a \\ 0 & 1\end{array}\right)$, we may assume that $a=0$, since $M X M^{-1}=X$ and $M Y M^{-1}=\left(\begin{array}{cc}0 & b^{\prime} \\ -1 & 0\end{array}\right)$. Thus, we have that $Y=\left(\begin{array}{cc}0 & b \\ -1 & 0\end{array}\right)$ in $\mathcal{C}_{2}$.

Lastly, consider the fact that $\operatorname{Tr}\left(Y^{2}\right)=2^{k}$, we have $b=-2^{k-1}$, giving us $Y=\left(\begin{array}{cc}0 & -2^{k-1} \\ -1 & 0\end{array}\right)$.

Now that we have nice formulas for $X$ and $Y$, we can explicitly see how our group action, defined by $(X, Y) \mapsto(X+q(Y), Y)$, acts on this specific pair of matrices. We wish to determine which $q \in \mathbb{C}[t]$ will satisfy
$X+q(Y)=X$, thus stabilizing the point $\left(0,0,0,-1,2^{k}\right)$. We start with $q(t)=\sum_{i=1}^{n} \alpha_{2 i} t^{2 i}$. Then we have that

$$
X \mapsto X+\sum_{i=1}^{n} \alpha_{2 i}\left(2^{k-1}\right)^{i} I_{2} .
$$

Thus, to stabilize the point, we must have that $\sum_{i=1}^{n}\left(2^{k-1}\right)^{i} \alpha_{2 i}=0$.
A similar argument shows that if $q(t)=\sum_{i=1}^{n} \alpha_{2 i+1} t^{2 i+1}$, then we must have that $\sum_{i=1}^{n}\left(2^{k-1}\right)^{i} \alpha_{2 i+1}=0$.

Therefore, concerning the set of points

$$
\left\{(0,0,0, \pm 1,0),(0,0,0,-1,2),(0,0,0,-1,4), \cdots,\left(0,0,0,-1,2^{k}\right)\right\}
$$

for some $k \in \mathbb{Z}_{+}$, if we have that $\sum_{i=1}^{n}\left(2^{j-1}\right)^{i} \alpha_{2 i}=0$ and $\sum_{i=1}^{n}\left(2^{j-1}\right)^{i} \alpha_{2 i+1}=0$ for all $1 \leqslant j \leqslant k$, then all the above points will be stabilized under the action of $\Psi_{q}$. This brings us to the following Lemma, which deals with finding solutions to this necessary system of equations obtained from the previous discussion:

Lemma 3. The solution set to the system of equations given by

$$
\left\{\begin{array}{l}
a_{1}+a_{2}+\ldots+a_{n}=0 \\
2 a_{1}+2^{2} a_{2}+\ldots+2^{n} a_{n}=0 \\
4 a_{1}+4^{2} a_{2}+\ldots+4^{n} a_{n}=0 \\
\vdots \\
2^{n-2} a_{1}+\left(2^{n-2}\right)^{2} a_{2}+\ldots+\left(2^{n-2}\right)^{n} a_{n}=0
\end{array}\right.
$$

can be expressed in terms of $a_{n}$ as

$$
\left\{\begin{array}{l}
a_{1}=S(n-1, n-2) a_{n} \\
a_{2}=S(n-2, n-2) a_{n} \\
a_{3}=S(n-3, n-2) a_{n} \\
\vdots \\
a_{n-1}=S(1, n-2) a_{n} \\
a_{n}=a_{n}
\end{array}\right.
$$

where $S(i, j)$ is the $i^{\text {th }}$ symmetric sum on the set $\left\{-1,-2,-4, \ldots,-2^{j}\right\}$.
Proof. The base of induction for $n=2$ is straightforward.

For our inductive hypothesis, assume the claim is true for $n=k$, and let us consider the case $n=k+1$. We begin by noting that for $1 \leqslant i \leqslant k$, $S(i, k-1)=S(i-1, k-2)-2^{k-1} S(i, k-2)$.

If $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ is the solution set of the $(k-1)$-dimensional system, then by the inductive hypothesis we have that

$$
\begin{equation*}
a_{i}=S(k-i, k-2) a_{k} \quad 1 \leqslant i \leqslant k-1 \tag{5}
\end{equation*}
$$

Define $b_{i}$ for $1 \leqslant i \leqslant k+1$ such that $b_{i}=a_{i-1}-2^{k-1} a_{i}$, where $a_{0}$ and $a_{k+1}$ are defined to be 0 . Substituting in (5), we get that

$$
b_{i}=S(k-i+1, k-2) a_{k}-2^{k-1} S(k-i, k-2) a_{k}=S(i, k-1) a_{k}
$$

for $1 \leqslant i \leqslant k$. Furthermore, we get that $b_{k+1}=a_{k}$. Thus, to prove the claim, it suffices to show that

$$
\left(b_{1}, b_{2}, \ldots, b_{k+1}\right)=\left(0-2^{k-1} a_{1}, a_{1}-2^{k-1} a_{2}, \ldots, a_{i}-2^{k-1} a_{i+1}, \ldots, a_{k}\right)
$$

is a solution to the system of $k$ equations,

$$
\left\{\begin{array}{l}
x_{1}+x_{2}+\ldots+x_{k+1}=0 \\
2 x_{1}+2^{2} x_{2}+\ldots+2^{n} x_{k+1}=0 \\
4 x_{1}+4^{2} x_{2}+\ldots+4^{n} x_{k+1}=0 \\
\vdots \\
2^{k-1} x_{1}+\left(2^{k-1}\right)^{2} x_{2}+\ldots+\left(2^{k-1}\right)^{n} x_{k+1}=0
\end{array}\right.
$$

For the equation

$$
2^{i} x_{1}+\left(2^{i}\right)^{2} x_{2}+\ldots+\left(2^{i}\right)^{n} x_{k+1}=0
$$

we plug in $b_{j}$ for $x_{j}$ to get

$$
\begin{aligned}
& 2^{i} b_{1}+\left(2^{i}\right)^{2} b_{2}+\ldots+\left(2^{i}\right)^{n} b_{k+1} \\
& \quad=2^{k-1}\left(-2^{i} a_{1}-2^{2 i} a_{2}-\ldots-2^{k i} a_{k}\right)+2^{i}\left(2^{i} a_{1}+2^{2 i} a_{2}+\ldots+2^{k i} a_{k}\right)
\end{aligned}
$$

Now, for $i<k-1$, by the inductive hypothesis,

$$
2^{i} a_{1}+2^{2 i} a_{2}+\ldots+2^{j i} a_{j}+\ldots+2^{k i} a_{k}=0
$$

so both terms become 0 . For $i=k-1$, we have by direct substitution that

$$
\begin{aligned}
& 2^{k-1}\left(-2^{k-1} a_{1}-2^{2(k-1)} a_{2}-\ldots-2^{j(k-1)} a_{j}-\ldots-2^{k(k-1)} a_{k}\right) \\
& \quad+2^{k-1}\left(2^{k-1} a_{1}+2^{2(k-1)} a_{2}+\ldots+2^{j(k-1)} a_{j}+\ldots+2^{k(k-1)} a_{k}\right)=0
\end{aligned}
$$

Thus, we can conclude that $\left(b_{1}, b_{2}, \ldots, b_{k+1}\right)$ satisfy the system of equations for $n=k+1$, and since $b_{i}=S(i, k-1) a_{k}=S(i, k-1) b_{k+1}$, we are done.

From this and the preceding discussion, we immediately conclude that if we wish to stabilize all the points

$$
\left\{(0,0,0, \pm 1,0),(0,0,0,-1,2),(0,0,0,-1,4), \cdots,\left(0,0,0,-1,2^{k}\right)\right\}
$$

using a polynomial with only even powers, we can use
$q_{k}^{E}(t):=\alpha t^{2(k+1)}+\alpha S(1, k-1) t^{2 k}+\cdots \alpha S(k-1, k-1) t^{4}+\alpha S(k, k-1) t^{2}$.
Similarly, if we wish to use a polynomial with only odd powers, we can use
$q_{k}^{O}(t):=\alpha t^{2 k+3}+\alpha S(1, k-1) t^{2 k+1}+\cdots \alpha S(k-1, k-1) t^{5}+\alpha S(k, k-1) t^{3}$.
Notation. Let $q_{k}^{E}(t)$ and $q_{k}^{O}(t)$ be defined as above. We then define the following notation: $\Psi_{k}^{E}:=\Psi_{q_{k}^{E}}$ and $\Psi_{k}^{O}:=\Psi_{q_{k}}$.

From the above arguments, we can conclude that $\Psi_{k}^{E}$ and $\Psi_{k}^{O}$ will stabilize the set

$$
\left\{(0,0,0, \pm 1,0),(0,0,0,-1,2),(0,0,0,-1,4), \cdots,\left(0,0,0,-1,2^{k}\right)\right\}
$$

Now, consider elements of the form $\Psi_{q} \circ \Phi_{p} \circ \Psi_{-q}$ and $\Phi_{p} \circ \Psi_{q} \circ \Phi_{-p}$, which we will call conjugation by $\Psi_{q}$ and $\Phi_{p}$, respectively. We will use the following four lemmas without proof, since the proofs are not difficult.

Lemma 4. The action $\Psi_{t} \circ \Phi_{p(t)} \circ \Psi_{-t}$ stabilizes the point $\left(0,0,0,1, a_{5}\right)$ if and only if $t^{2}-\frac{a_{5}-2}{2}$ divides $p(t)$ and it stabilizes $\left(0,0,0,-1, a_{5}\right)$ if and only if $t^{2}-\frac{a_{5}+2}{2}$ divides $p(t)$.

Lemma 5. The action $\Phi_{t} \circ \Psi_{q(t)} \circ \Phi_{-t}$ stabilizes the point $\left(0,0,0,1, a_{5}\right)$ if and only if $t^{2}-\frac{a_{5}-2}{2}$ divides $q(t)$ and it stabilizes $\left(0,0,0,-1, a_{5}\right)$ if and only if $t^{2}-\frac{a_{5}+2}{2}$ divides $q(t)$.

Lemma 6. The action $\Psi_{-t} \circ \Phi_{p(t)} \circ \Psi_{t}$ stabilizes the point $(0,0,0,1, a)$ if and only if $t^{2}-\frac{a+2}{2}$ divides $p(t)$ and it stabilizes $(0,0,0,-1, b)$ if and only if $t^{2}-\frac{b-2}{2}$ divides $p(t)$.

Lemma 7. The action $\Phi_{-t} \circ \Psi_{q(t)} \circ \Phi_{t}$ stabilizes the point $(0,0,0,1, a)$ if and only if $t^{2}-\frac{a+2}{2}$ divides $q(t)$ and it stabilizes $(0,0,0,-1, b)$ if and only if $t^{2}-\frac{b-2}{2}$ divides $q(t)$.

The following conjugations

$$
\begin{array}{ll}
\Phi_{t} \circ \Psi_{q_{k}^{C}} \circ \Phi_{-t} ; & \Psi_{t} \circ \Phi_{p_{k}^{C}} \circ \Psi_{-t} ; \\
\Phi_{-t} \circ \Psi_{\tilde{q}_{k}^{C}} \circ \Phi_{t} ; & \Psi_{-t} \circ \Phi_{\tilde{p}_{k}^{C}} \circ \Psi_{t}
\end{array}
$$

will all stabilize the set

$$
\begin{equation*}
\left\{(0,0,0, \pm 1,0),(0,0,0,-1,2), \cdots,\left(0,0,0,-1,2^{k}\right)\right\} \backslash\{(0,0,0,-1,8)\} \tag{6}
\end{equation*}
$$

where

$$
\begin{aligned}
& q_{k}^{C}(t)=\alpha\left(t^{2}+1\right)\left(t^{2}-1\right)\left(t^{2}-2\right)\left(t^{2}-3\right)\left(t^{2}-9\right) \cdots\left(t^{2}-\left(2^{k-1}+1\right)\right) \\
& p_{k}^{C}(t)=\alpha\left(t^{2}+1\right)\left(t^{2}-1\right)\left(t^{2}-2\right)\left(t^{2}-3\right)\left(t^{2}-9\right) \cdots\left(t^{2}-\left(2^{k-1}+1\right)\right), \\
& \tilde{p}_{k}^{C}(t)=\alpha\left(t^{2}-1\right)\left(t^{2}+1\right) t^{2}\left(t^{2}-7\right)\left(t^{2}-15\right) \cdots\left(t^{2}-\left(2^{k-1}-1\right)\right) \\
& \tilde{q}_{k}^{C}(t)=\alpha\left(t^{2}-1\right)\left(t^{2}+1\right) t^{2}\left(t^{2}-7\right)\left(t^{2}-15\right) \cdots\left(t^{2}-\left(2^{k-1}-1\right)\right) .
\end{aligned}
$$

## 5. More details concerning $\Psi_{k}^{E}$

This section is still concerned with determining the structure of the stabilizers. Specifically, we consolidate information about the element $\Psi_{k}^{E}$ that will be useful for the proof of $n$-transitivity presented in the next section. We are especially interested in how this element acts on large subsets of $\mathcal{C}_{2}$, which will appear later. The first result we need is the following formula for $2 \times 2$ matrices being raised to integer powers. It is an easy calculation by induction, and so we omit the proof.

Proposition 5. Let $M \in M_{2}(\mathbb{C})$ be an arbitrary matrix with two distinct eigenvalues. Let $\mu=\operatorname{Tr}(M)$ and $\nu=-\operatorname{det}(M)$. Then, for any $k \in \mathbb{Z}_{+}$, $M^{k}=\mu_{k} M+\nu \mu_{k-1} I$, where

$$
\mu_{k}=\frac{1}{\sqrt{\mu^{2}+4 \nu}}\left(\left(\frac{\mu+\sqrt{\mu^{2}+4 \nu}}{2}\right)^{k}-\left(\frac{\mu-\sqrt{\mu^{2}+4 \nu}}{2}\right)^{k}\right)
$$

for any $k \in \mathbb{Z}_{+}$.
We can now use this result to determine how $\Psi_{\gamma t^{2 n}}$ will act on points with $a_{2}=0$ :

Lemma 8. Applying $\Psi_{\gamma t^{2 n}}$ for any $n \in \mathbb{Z}_{+}$to the point $\left(a_{1}, 0, a_{3}, a_{4}, a_{5}\right)$, we arrive at the point $\left(a_{1}+\gamma a_{5}\left(\frac{a_{5}}{2}\right)^{n-1}, 0, a_{3}, a_{4}, a_{5}\right)$.

Proof. Let $A=\left(a_{1}, 0, a_{3}, a_{4}, a_{5}\right)$. We note that by definition, $a_{2}, a_{5}$ are fixed by the action of $\Psi$. Now consider $a_{1}$. Again by definition of $\Psi$, we know that

$$
\begin{equation*}
a_{1}^{\prime}=\Psi_{\gamma t^{2 n}}\left(a_{1}\right)=a_{1}+\operatorname{Tr}\left(\gamma Y^{2 n}\right)=a_{1}+\gamma \operatorname{Tr}\left(Y^{2 n}\right) \tag{7}
\end{equation*}
$$

If $a_{5} \neq 0$, we can now use the formula in Proposition 5 to get that $Y^{2 n}=\frac{a_{5}^{n}}{2^{n}} I$. Plugging into (7), we get

$$
a_{1}^{\prime}=a_{1}+\gamma \operatorname{Tr}\left(\frac{a_{5}^{n}}{2^{n}} I_{2}\right)=a_{1}+\frac{\gamma a_{5}^{n}}{2^{n-1}} .
$$

Next we consider the action of $\Psi_{\gamma t^{2 n}}$ on $a_{3}$ :

$$
\begin{aligned}
\Psi_{\gamma t^{2 n}}\left(a_{3}\right) & =a_{3}+2 \operatorname{Tr}(A q(Y))+\operatorname{Tr}\left(q^{2}(Y)\right)-\frac{1}{2} \operatorname{Tr}^{2}(q(Y)) \\
& =a_{3}+2 \operatorname{Tr}\left(A \cdot \frac{\gamma a_{5}^{n}}{2^{n}} I_{2}\right)+\operatorname{Tr}\left(\gamma^{2} \frac{a_{5}^{2 n}}{2^{2 n}} I_{2}\right)-\frac{1}{2} \operatorname{Tr}^{2}\left(\gamma \frac{a_{5}^{n}}{2^{n}} I_{2}\right) \\
& =a_{3}+\gamma^{2} \frac{a_{5}^{2 n}}{2^{2 n-1}}-\frac{1}{2}\left(\gamma \frac{a_{5}^{n}}{2^{n-1}}\right)^{2} \\
& =a_{3}+\gamma^{2} \frac{a_{5}^{2 n}}{2^{2 n-1}}-\frac{1}{2} \gamma^{2} \frac{a_{5}^{2 n}}{2^{2 n-2}} \\
& =a_{3}
\end{aligned}
$$

as claimed. We lastly consider the action on $a_{4}$ :

$$
\Psi_{\gamma t^{2 n}}\left(a_{4}\right)=a_{4}+\operatorname{Tr}\left(Y \cdot \gamma \frac{a_{5}^{n}}{2^{n}} I_{2}\right)=a_{4}+\gamma \frac{a_{5}^{n}}{2^{n}} \operatorname{Tr}(Y)=a_{4}
$$

Thus, we see that $\Psi_{\gamma t^{2 n}}(A)=\left(a_{1}+\gamma \frac{a_{5}^{n}}{2^{n-1}}, 0, a_{3}, a_{4}, a_{5}\right)$, as desired.
If $a_{5}=0$, then we know that $Y$ is nilpotent, and hence $Y^{2 n} \equiv 0$ for all $n \in \mathbb{Z}_{+}$. It follows that $\Psi_{\gamma t^{2 n}}(A)=A$, so that the formula holds in this case as well.

Lemma 9. Applying $\Psi_{\gamma t^{2 n+1}}$ for any $n \in \mathbb{Z}_{+}$to the point $\left(a_{1}, 0, a_{3}, a_{4}, a_{5}\right)$, we arrive at the point

$$
\left(a_{1}, 0, a_{3}+\gamma^{2} \frac{a_{5}^{2 n+3}}{2^{2(n+1)}}, a_{4}+\gamma \frac{a_{5}^{n+2}}{2^{n+1}}, a_{5}\right)
$$

Proof. The proof is analogous to that of Lemma 8.
Corollary 1. Let $A=\left(a_{1}, 0, a_{3}, a_{4}, a_{5}\right)$. Then $A^{\prime}=\Psi_{k}^{E}(A)$ (respectively, $\left.A^{\prime}=\Psi_{k}^{O}(A)\right)$ satisfies $a_{1}^{\prime} \equiv a_{1}$ (respectively, $a_{3}^{\prime} \equiv a_{3}$ ) if and only if $a_{5} \in R_{2}^{k}:=\left\{0,2,4, \ldots, 2^{k}\right\}$.

Proof. The proof is a straightforward application of Vieta's formula, using Lemma 8 (respectively, Lemma 9) and the fact that $\Psi_{f} \circ \Psi_{g}=\Psi_{f+g}$.

## 6. $n$-Transitivity

We now have everything we need in order to prove our main result: infinite-transitivity. We start with two supporting lemmas, analogous to Lemmas 1 and 2, and then proceed to the final theorem.
Notation. For any $k \in \mathbb{Z}_{+}$, we define the following notation:

$$
\begin{gathered}
C_{k}:=\left\{(0,0,0, \pm 1,0), \ldots,\left(0,0,0,-1,2^{k}\right)\right\} \backslash\{(0,0,0,-1,8)\} \subseteq \mathcal{C}_{2} \\
S_{k}:=\operatorname{Stab}\left[C_{k}\right] \subseteq G
\end{gathered}
$$

Lemma 10. Let $A=\left(a_{1}, a_{2}, a_{3}, a_{4}, a_{5}\right) \in \mathcal{C}_{2} \backslash C_{k}$ for some $k \in \mathbb{Z}_{+}$. Then there is an element $g \in S_{k}$ such that $A^{\prime}=g(A)$ satisfies $a_{1}^{\prime} a_{3}^{\prime} \neq 0$.

Proof. Case 1: $a_{1} \neq 0$. If $a_{3} \neq 0$ as well, then we are done, so we assume $a_{3}=0$. Hence, the point is $A=\left(a_{1}, a_{2}, 0, \pm 1, a_{5}\right)$. Applying $\Phi_{-\frac{2 a_{2}}{a_{1}^{2}}-\frac{a_{1}^{2}}{4} \beta}$ to $A$, we arrive at the point $\left(a_{1}, 0,0, \pm 1, \mp \frac{4 a_{2}}{a_{1}}+a_{5} \pm \frac{1}{2} a_{1}^{3} \beta\right)$. Since $a_{1} \neq 0$ we can choose the parameter $\beta$ such that $a_{5} \notin R_{2}^{k}$. Then using Lemma 9 and Corollary 1, we know that after applying $\Psi_{k}^{O}$, we will have that $a_{1}^{\prime}=a_{1} \neq 0$ and $a_{3}^{\prime}$ is non-constant polynomial in $\gamma$. Thus we can choose some $\gamma$ such that $a_{3}^{\prime} \neq 0$, as desired.
Case 2: $a_{3} \neq 0$. Similarly to Case 1 , we assume that $a_{1}=0$, since otherwise we are done. Then, we apply the element $\Phi_{\alpha t^{2}}$ to arrive at the point

$$
A^{\prime}=\left(0, a_{2}+\alpha a_{3}, a_{3}, a_{4}, a_{5}\right)
$$

We then see that

$$
\begin{aligned}
\left(a_{5}^{\prime}-\frac{1}{2}\right. & \left.\left(2-3 a_{2}^{\prime 2}-2 \sqrt{1-2 a_{2}^{\prime 2}+2 a_{2}^{\prime 4}}\right)\right) \\
& \cdot\left(a_{5}^{\prime}-\frac{1}{2}\left(2-3 a_{2}^{\prime 2}+2 \sqrt{1-2 a_{2}^{\prime 2}+2 a_{2}^{\prime 4}}\right)\right) \\
= & -a_{2}^{2}+\frac{a_{2}^{4}}{4}-2 a_{5}+3 a_{2}^{2} a_{5}+a_{5}^{2}+\left(-a_{2} a_{3}+a_{2}^{3} a_{3}+6 a_{2} a_{3} a_{5}\right) \alpha \\
& +\left(-a_{3}^{2}+\frac{3 a_{2}^{2} a_{3}^{2}}{2}+3 a_{3}^{2} a_{5}\right) \alpha^{2}+a_{2} a_{3}^{3} \alpha^{3}+\frac{a_{3}^{4} \alpha^{4}}{4}
\end{aligned}
$$

Since $a_{3} \neq 0$, we conclude that the coefficient of $\alpha^{4}$ is non-zero, and hence this is a non-constant polynomial of $\alpha$. We then choose $\alpha \in \mathbb{C}$ such that

$$
\begin{equation*}
a_{5}^{\prime}-\frac{1}{2}\left(2-3 a_{2}^{\prime 2} \pm 2 \sqrt{1-2 a_{2}^{\prime 2}+2 a_{2}^{\prime 4}}\right) \neq 0 \tag{8}
\end{equation*}
$$

Then, we can apply the element $\Phi_{t^{4}-t^{2}}$ to $A^{\prime}$ to reach a point with

$$
\tilde{a}_{1}=\frac{1}{8}\left(a_{2}^{\prime 4}+4\left(-2+a_{5}^{\prime}\right) a_{5}^{\prime}+4 a_{2}^{\prime 2}\left(-1+3 a_{5}^{\prime}\right)\right)
$$

This is zero if and only if $a_{5}^{\prime}=\frac{1}{2}\left(2-3 a_{2}^{\prime 2} \pm 2 \sqrt{1-2 a_{2}^{\prime 2}+2 a_{2}^{\prime 4}}\right)$, but by equation (8), we know this is not the case. Hence we are at a point with $\tilde{a}_{1} \neq 0$, so that we are back in Case 1.
Case 3: $a_{1}=a_{3}=0$. Since $a_{3}=0$, we know that $a_{4}= \pm 1$ and at least one of $a_{2}, a_{5} \neq 0$, as otherwise the point would be nilpotent, and hence in $C_{k}$.

Case 3.1: $a_{5} \notin R_{2}^{k}$. Recall from Corollary 1 that $R_{2}^{k}=\left\{0,2,4, \ldots, 2^{k}\right\}$. If $a_{5} \notin R_{2}^{k}$, then this same corollary tells us that we can apply $\Psi_{k}^{E}$ for some value of $\alpha$ to get that $a_{1}^{\prime} \neq 0$. Then we are back in Case 1 .

Case 3.2: $a_{5} \in R_{2}^{k}$. This case requires the use of conjugation. For ease of notation, let $T=\left\{\left(0, a_{2}, 0, \pm 1, a_{5}\right)\right\}$ be the set of points that satisfy $a_{1}=a_{3}=0$.

Case 3.2.1: $a_{2} \neq 0$. Let $T_{0}=\left\{\left(0, a_{2}, 0, \pm 1, a_{5}\right)\right\}$ be the set of points we are considering here, so that $a_{2} \neq 0$ in $T_{0}$. We want to show that, using conjugation, we can move any element of $T_{0}$ to a point outside of $T$ (i.e. a point $A^{\prime}$ with $a_{1}^{\prime}$ or $a_{3}^{\prime} \neq 0$ ).

To do this, let $A=\left(0, a_{2}, 0,-1, a_{5}\right) \in T$ be an arbitrary point, and assume, for the sake of a contradiction, that none of the conjugations defined in section 4 move $A$ out of $T$. Then we have that for any choice of $\alpha$ in the polynomials $\tilde{q}_{k}^{C}$ and $\tilde{p}_{k}^{C}$, we must have that

$$
\left(\Phi_{t} \circ \Psi_{\tilde{q}_{k}^{C}} \circ \Phi_{t}\right)(A)=\left(0, a_{2}^{\prime}, 0, \pm 1, a_{5}^{\prime}\right) \in T
$$

and

$$
\left(\Psi_{-t} \circ \Phi_{\tilde{p}_{k}^{C}} \circ \Psi_{t}\right)(A)=\left(0, \hat{a}_{2}, 0, \pm 1, \hat{a}_{5}\right) \in T .
$$

Moving $\Phi_{t}$ and $\Psi_{t}$ to the other side of the equations, we see that these reduce to

$$
\begin{gather*}
\Psi_{\tilde{q}_{k}^{C}}\left(0, a_{2}, 0,-1, a_{5}-2\right)=\left(0, a_{2}^{\prime}, 0, \pm 1, a_{5}^{\prime} \pm 2\right)  \tag{9}\\
\Phi_{\tilde{p}_{k}^{C}}\left(a_{2}, a_{2}, a_{5}-2,-1+a_{5}, a_{5}\right)=\left(\hat{a}_{2}, \hat{a}_{2}, \hat{a}_{5} \pm 2, \pm 1+\hat{a}_{5}, \hat{a}_{5}\right) . \tag{10}
\end{gather*}
$$

We will use these equations to prove the following claim:
Claim 1. Let $C_{Y}(t)$ denote the characteristic polynomial of $Y$. Then we have that $C_{Y}(t) \mid \tilde{q}_{k}^{C}(t)$.

For now, we will assume this claim is true. This implies, since $\tilde{p}_{k}^{C}(t)=$ $\tilde{q}_{k}^{C}(t)$, that the eigenvalues of $X$ and $Y$ are roots of $\tilde{p}_{k}^{C}(t)$.

We now consider the element $\Psi_{k}^{E}$. We know from Lemma 8 and Corollary 1 that $\Psi_{k}^{E}$ fixes the point $A$. Thus, using the definition of the action by $\Psi_{q}$, we see that $\operatorname{Tr}\left(q_{k}^{E}(Y)\right)=\Psi_{k}^{E}\left(a_{1}\right)=0$ and $2 \operatorname{Tr}(A$. $\left.q_{k}^{E}(Y)\right)+\operatorname{Tr}\left(q_{k}^{E}(Y)^{2}\right)=\Psi_{k}^{E}\left(a_{3}\right)=0$. This is true for all $\alpha$, and hence, since $2 \operatorname{Tr}\left(A \cdot q_{k}^{E}(Y)\right)$ is linear in $\alpha$ while $\operatorname{Tr}\left(q_{k}^{E}(Y)^{2}\right)$ is quadratic in $\alpha$, we must have that their coefficients are 0 separately. In particular, we get that $\operatorname{Tr}\left(q_{k}^{E}(Y)\right)=0$ and $\operatorname{Tr}\left(q_{k}^{E}(Y)^{2}\right)=0$, so that $q_{k}^{E}(Y)$ is nilpotent, and hence the eigenvalues of $q_{k}^{E}(Y)$ are both 0 . If we denote the eigenvalues of $Y$ by $\lambda_{1}, \lambda_{2}$, then we obtain the fact that $q_{k}^{E}\left(\lambda_{1}\right)=0=q_{k}^{E}\left(\lambda_{2}\right)$ for all $\alpha \in \mathbb{C}$, so that $\lambda_{1}$ and $\lambda_{2}$ are roots of $q_{k}^{E} / \alpha$.

Then, combining this with the preceding statement, we see that the eigenvalues $\lambda_{1}$ and $\lambda_{2}$ are roots of both

$$
\begin{gathered}
q_{k}^{E}(t) / \alpha=t^{2(k+1)}+S(1, k-1) t^{2 k}+\cdots+S(k-1, k-1) t^{4}+S(k, k-1) t^{2} \\
\tilde{p}_{k}^{C}(t) / \alpha=\left(t^{2}-1\right)\left(t^{2}+1\right) t^{2}\left(t^{2}-7\right)\left(t^{2}-15\right) \cdots\left(t^{2}-\left(2^{k-1}-1\right)\right)
\end{gathered}
$$

The only roots these equations share are given by 0 and $\pm 1$, and since $\tilde{p}_{k}^{C}$ has no double roots, we know that $\lambda_{1} \neq \lambda_{2}$. Additionally, since $\lambda_{1}+\lambda_{2}=$ $\operatorname{Tr}(Y)=a_{2} \neq 0$ by the case assumption, we also know that $\lambda_{1} \neq-\lambda_{2}$. Thus we can conclude that if both $\left(\Phi_{-t} \circ \Psi_{\tilde{q}_{k}^{C}} \circ \Phi_{t}\right)(A) \in T$ and $\left(\Psi_{-t} \circ\right.$ $\left.\Phi_{\tilde{p}_{k}^{C}} \circ \Psi_{t}\right)(A) \in T$, then $A$ satisfies $a_{2}= \pm 1$.

Running the same argument with $\Phi_{t} \circ \Psi_{q_{k}^{C}} \circ \Phi_{-t}$ and $\Psi_{t} \circ \Phi_{p_{k}^{C}} \circ \Psi_{-t}$, we get that if both $\left(\Phi_{t} \circ \Psi_{q_{k}^{C}} \circ \Phi_{-t}\right)(A) \in T$ and $\left(\Psi_{t} \circ \Phi_{p_{k}^{C}} \circ \Psi_{-t}\right)(A) \in T$, then $A$ must satisfy $a_{2} \in\{1 \pm \sqrt{2},-1 \pm \sqrt{2}\}$.

Since, by assumption we have that all of these conjugations land $A$ back in $T$, we thus conclude that $a_{2} \in\{ \pm 1\}$ and $a_{2} \in\{1 \pm \sqrt{2},-1 \pm \sqrt{2}\}$, so that $a_{2} \in\{ \pm 1\} \cap\{1 \pm \sqrt{2},-1 \pm \sqrt{2}\}$ which is a contradiction. Thus, we must have that at least one of the conjugations moves the point $A$ out of the set $T$, thus landing us back in a previous case where $a_{1} \neq 0$ or $a_{3} \neq 0$. An analogous argument gives the same result if we start with a point $A$ where $a_{4}=1$.

Case 3.2.2: $a_{2}=0$. Since we also have that $a_{1}=a_{3}=0$, in order for our point $A$ to not be $C_{k}$, we see that we are reduced to considering the set of points

$$
\begin{gather*}
\{(0,0,0,1,2),(0,0,0,1,4),(0,0,0, \pm 1,8) \\
\left.(0,0,0,1,16), \ldots,\left(0,0,0,1,2^{k}\right)\right\} \tag{11}
\end{gather*}
$$

We will now show that using the element $\Psi_{-t} \circ \Phi_{\tilde{p}_{k}^{C}} \circ \Psi_{t}$, we can move all of these points out of this set. To do this, we first note that applying any number of times repeatedly gives a conjugation with $\Phi_{n \tilde{p}_{k}^{C}}$ in the middle. Now let $A$ be a point in (11), and assume that $\left(\Psi_{-t} \circ \Phi_{n \tilde{p}_{k}^{C}} \circ \Psi_{t}\right)(A)$ is in (11) for all $n \in \mathbb{Z}_{+}$. Then, since there are only finitely many points in (11), this means that there must be some $n, m \in \mathbb{Z}_{+}$with $m>n$ such that $\left(\Psi_{-t} \circ \Phi_{n \tilde{p}_{k}^{C}} \circ \Psi_{t}\right)(A)=\left(\Psi_{-t} \circ \Phi_{m \tilde{p}_{k}^{C}} \circ \Psi_{t}\right)(A)=\left(\left(\Psi_{-t} \circ \Phi_{(m-n) \tilde{p}_{k}^{C}} \circ \Psi_{t}\right) \circ\right.$ $\left(\Psi_{-t} \circ \Phi_{n \tilde{p}}^{k}\right.$ © $\left.\left.\circ \Psi_{t}\right)\right)(A)$. In particular, we see that $\Psi_{-t} \circ \Phi_{(m-n) \tilde{p}_{k}^{C}} \circ \Psi_{t}$ fixes the point $\left(\Psi_{-t} \circ \Phi_{n \tilde{p}}^{k}\right.$ $\left.\circ \Psi_{t}\right)(A)$ in (11). However, this is a contradiction to Lemma 6, which guarantees that no points in (11) are fixed by conjugation of the form $\Psi_{-t} \circ \Phi_{n \tilde{p}_{k}^{C}} \circ \Psi_{t}$ for $n \in \mathbb{Z}_{+}$. Thus we must have that successive usage of the element $\Psi_{-t} \circ \Phi_{\tilde{p}_{k}^{C}} \circ \Psi_{t}$ will move any point in (11) out of this set and into a previous case.

To finish the proof, we need to prove Claim 1.
Proof of Claim 1 We start by showing that $a_{2}=a_{2}^{\prime}=\hat{a}_{2}$ and $a_{5}=a_{5}^{\prime}=$ $\hat{a}_{5}$. To do this, we first consider equation (9). Since $\Psi$ fixes $a_{2}$ and $a_{5}$, this equation tells us immediately that $a_{2}=a_{2}^{\prime}$. Thus we only need to consider $a_{5}$. There are two cases:

$$
\begin{gathered}
\Psi_{\tilde{q}_{k}^{C}}\left(0, a_{2}, 0,-1, a_{5}-2\right)=\left(0, a_{2}^{\prime}, 0,1, a_{5}^{\prime}+2\right) \\
\Psi_{\tilde{q}_{k}^{C}}\left(0, a_{2}, 0,-1, a_{5}-2\right)=\left(0, a_{2}^{\prime}, 0,-1, a_{5}^{\prime}-2\right)
\end{gathered}
$$

If the first of these is true, then we must have, by examining the 4th term, that $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right)=2$ for that specific choice of $\alpha$, and if the second is true we must have that, for those specific $\alpha \in \mathbb{C}$, that $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right)=0$. We easily see from the definition of $\tilde{q}_{k}^{C}$ that if $\alpha=0$, then $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right)=0$. This implies either that, as a function of $\alpha$, we have that $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right) \equiv 0$ or $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right)$ is non-constant. If the second of these options is true, then we can choose some $\alpha \in \mathbb{C}$ such that $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right) \neq 0,2$. However, for such an $\alpha$, we then have that $a_{4}^{\prime} \neq \pm 1$, which contradicts the fact that we must have $a_{3}^{\prime} a_{5}^{\prime}+a_{4}^{\prime 2}=1$. Thus we must have that $\operatorname{Tr}\left(B \cdot \tilde{q}_{k}^{C}(Y)\right) \equiv 0$, so that $a_{5}-2=a_{5}^{\prime}-2$ for all $\alpha \in \mathbb{C}$. Thus we conclude that $a_{5}=a_{5}^{\prime}$, as claimed. An analogous argument using equation (10) gives that $a_{2}=\hat{a}_{2}$ and $a_{5}=\hat{a}_{5}$.

Thus, using the second equation above that we have found to be the case, we must have that

$$
\Psi_{\tilde{q}_{k}^{C}}\left(0, a_{2}, 0,-1, a_{5}-2\right)=\left(0, a_{2}, 0,-1, a_{5}-2\right)
$$

so that $\Psi_{\tilde{q}_{k}^{C}}$ fixes the point $\left(0, a_{2}, 0,-1, a_{5}-2\right)$. We can then write that $\tilde{q}_{k}^{C}(t)=C_{Y}(t) f(t)+r(t)$, where $r(t)$ has degree $\leqslant 1$, since $C_{Y}(t)$ has degree 2. However, since Cayley-Hamilton guarantess that $C_{Y}(Y)=0$, we conclude that $\tilde{q}_{k}^{C}(Y)=r(Y)$, so that $\Psi_{\tilde{q}_{k}^{C}}=\Psi_{r}$, and hence $\Psi_{r}$ fixes the point $\left(0, a_{2}, 0,-1, a_{5}-2\right)$. However, it is easy to see that the only polynomial of degree less than or equal to 1 that fixes this point is the zero polynomial, and hence $r(t) \equiv 0$. This shows that $C_{Y}(t) \mid \tilde{q}_{k}^{C}(t)$, as claimed.

We now move onto the second lemma that will be necessary in proving infinite-transitivity.

Lemma 11. Let $A \in \mathcal{C}_{2} \backslash C_{k}$ for some $k \in \mathbb{Z}_{+}$satisfy $a_{1} a_{3} \neq 0$. Then there is a $g \in S_{k} \subseteq G$ such that $A^{\prime}=g(A)$ satisfies both $a_{1}^{\prime} a_{3}^{\prime} \neq 0$ and $a_{1}^{4}-4 a_{3}^{\prime 2} \neq 0$.

Proof. We may assume that $a_{1}^{2}-4 a_{3}^{2}=0$, since otherwise we are done.
Case 1: $a_{1}^{2}=-2 a_{3}$. Applying the element $\Phi_{\gamma t^{2}+\beta t^{4}}$ with $\beta=-\frac{2 a_{2}}{a_{1}^{4}}$ we have that $a_{2}^{\prime}=0$ and

$$
a_{5}^{\prime}=a_{5}+a_{1} \gamma\left(2 a_{4}-\frac{a_{1}^{3} \gamma}{2}\right)=a_{5}+2 a_{1} a_{4} \gamma-\frac{a_{1}^{4} \gamma^{2}}{2}
$$

while $a_{1}, a_{3}$ stay fixed. We note that since $a_{1} \neq 0$, the coefficient of $\gamma^{2}$ is non-zero, and hence $a_{5}^{\prime}$ is a non-constant polynomial in $\gamma$, so that we can choose some $\gamma \in \mathbb{C}$ such that $a_{5}^{\prime} \notin\left\{0,2,4, \ldots, 2^{k}\right\}$. Let $\tilde{A}=\Psi_{k}^{E}\left(A^{\prime}\right)$. Using Lemma 8 and Corollary 1, we know that $\tilde{a}_{1}$ is a non-constant polynomial in $\alpha$ while $\tilde{a}_{3}=a_{3}=\frac{a_{1}^{2}}{2}$ is non-zero and constant. Thus $\tilde{a}_{1}^{2} \pm 2 \tilde{a}_{3}$ will be a non-constant polynomials in $\alpha$, since a constant function cannot cancel higher order non-constant terms. Thus we can choose some $\alpha \in \mathbb{C}$ such that $\tilde{a}_{1} \tilde{a}_{3} \neq 0$ and $\tilde{a}_{1}^{2} \pm 2 \tilde{a}_{3} \neq 0$, as desired.

Case 2: $a_{1}^{2}=2 a_{3}$. We start at the point $A=\left(a_{1}, a_{2}, \frac{a_{1}^{2}}{2}, a_{4}, a_{5}\right)$. Consider the polynomial

$$
\begin{aligned}
q(t) & =t^{2}\left(t^{2}-1\right)\left(t^{2}-2\right)\left(t^{2}-4\right) \ldots\left(t^{2}-2^{k-1}\right) \\
& =t^{2(k+1)}+S(1, k-1) t^{2 k}+\cdots S(k-1, k-1) t^{4}+S(k, k-1) t^{2}
\end{aligned}
$$

and recall that, by construction, $\Psi_{\alpha q(t)}$ will stabilize $C_{k}$ for any $\alpha \in \mathbb{C}$.
First, using $\Phi_{\frac{\varepsilon-a_{2}}{a_{1}^{2}} t^{2}}$ we also stabilize the points in $C_{k}$ and send $A$ to $C=$ $\left(a_{1}, \varepsilon, \frac{a_{1}^{2}}{2}, a_{4}^{\prime}, a_{5}^{\prime}\right)$, where $a_{4}^{\prime}=a_{4}+\frac{\varepsilon-a_{2}}{2} a_{1}$ and $a_{5}^{\prime}=a_{5}+\frac{\left(\varepsilon-a_{2}\right)^{2}}{2}+\frac{2 a_{4}\left(\varepsilon-a_{2}\right)}{a_{1}}$.

Now, for $\Psi_{\alpha q(t)}(C)=\left(a_{1, \alpha}, a_{2, \alpha}, a_{3, \alpha}, a_{4, \alpha}, a_{5, \alpha}\right)$, we calculate the following (keep in mind that this family of automorphisms fixes points in $C_{k}$ ):

$$
\begin{aligned}
a_{1, \alpha} & =\alpha \cdot \operatorname{Tr}(q(Y))+a_{1} \\
a_{3, \alpha}=\alpha^{2} \cdot\left(\operatorname{Tr}\left(q^{2}(Y)\right)\right. & \left.-\frac{1}{2} \operatorname{Tr}^{2}(q(Y))\right)+2 \alpha \cdot \operatorname{Tr}(A q(Y))+\frac{a_{1}^{2}}{2}
\end{aligned}
$$

Our goal is to show that there exists a nonzero $\alpha$ such that $a_{3, \alpha}=-\frac{1}{2} a_{1, \alpha}^{2}$ with $a_{1, \alpha} \neq 0$, since we will then be back in Case 1. Let $f(\alpha)=a_{3, \alpha}+\frac{1}{2} a_{1, \alpha}^{2}$. We now need to find the roots of $f(\alpha)$.

$$
f(\alpha)=\alpha^{2} \cdot \operatorname{Tr}\left(q^{2}(Y)\right)+\alpha\left(2 \operatorname{Tr}(A q(Y))+a_{1} \operatorname{Tr}(q(Y))\right)+a_{1}^{2}
$$

Let $\alpha_{1}$ and $\alpha_{2}$ be roots of $f(\alpha)$. Then both of them are not zero since $a_{1} \neq 0$. Since $\operatorname{Tr}(Y)=\varepsilon$ and $\operatorname{Tr}\left(Y^{2}\right)=\operatorname{Tr}\left(B^{2}\right)+\frac{1}{2} \operatorname{Tr}^{2}(Y)=a_{5}^{\prime}+\frac{1}{2} \varepsilon^{2}$ we find the eigenvalues of $Y$ to be given by $\mu_{1}=\frac{\varepsilon}{2}+\frac{\sqrt{2 a_{5}^{\prime}}}{2}$ and $\mu_{2}=\frac{\varepsilon}{2}-\frac{\sqrt{2 a_{5}^{\prime}}}{2}$.

We also have that

$$
\operatorname{Tr}(q(Y))=q\left(\mu_{1}\right)+q\left(\mu_{2}\right) \quad \text { and } \quad \operatorname{Tr}\left(q^{2}(Y)\right)=q^{2}\left(\mu_{1}\right)+q^{2}\left(\mu_{2}\right)
$$

so we can calculate that

$$
\begin{aligned}
& \operatorname{Tr}(q(Y))=q\left(\mu_{1}\right)+q\left(\mu_{2}\right) \\
& \begin{aligned}
&= \mu_{1}{ }^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i} \mu_{1}{ }^{2 i}+\mu_{2}{ }^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i} \mu_{2}{ }^{2 i} \\
&= \mu_{1}{ }^{2(k+1)}+\mu_{2}{ }^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i}\left(\mu_{1}{ }^{2 i}+\mu_{2}{ }^{2 i}\right) \\
&=\left(\frac{\varepsilon}{2}+\frac{\sqrt{2 a_{5}^{\prime}}}{2}\right)^{2(k+1)}+\left(\frac{\varepsilon}{2}-\frac{\sqrt{2 a_{5}^{\prime}}}{2}\right)^{2(k+1)} \\
&+\sum_{i=1}^{k} \alpha_{i}\left(\left(\frac{\varepsilon}{2}+\frac{\sqrt{2 a_{5}^{\prime}}}{2}\right)^{2 i}+\left(\frac{\varepsilon}{2}-\frac{\sqrt{2 a_{5}^{\prime}}}{2}\right)^{2 i}\right) \\
&=\left.\frac{1}{2^{k+1}}\left(\varepsilon^{2}+2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{(k+1)}+\left(\varepsilon^{2}-2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{(k+1)}\right) \\
& \quad+\sum_{i=1}^{k} \frac{\alpha_{i}}{2^{2 i}}\left(\left(\varepsilon^{2}+2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{i}+\left(\varepsilon^{2}-2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{i}\right) .
\end{aligned}
\end{aligned}
$$

From this we know $\operatorname{Tr}(q(Y))$ is a polynomial in $\varepsilon$ of degree $2(k+1)$; call it $g(\varepsilon)$.

Similarly, we get that

$$
\begin{aligned}
\operatorname{Tr}\left(q^{2}(Y)\right)= & \frac{1}{4^{k+1}}\left(\left(\varepsilon^{2}+2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{2(k+1)}+\left(\varepsilon^{2}-2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{2(k+1)}\right) \\
& +\sum_{i=1}^{2 k+1} \frac{\beta_{i}}{4^{2 i}}\left(\left(\varepsilon^{2}+2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{i}+\left(\varepsilon^{2}-2 \varepsilon \sqrt{2 a_{5}^{\prime}}+2 a_{5}^{\prime}\right)^{i}\right)
\end{aligned}
$$

This implies that $\operatorname{Tr}\left(q^{2}(Y)\right)$ is a polynomial in $\varepsilon$ of degree $4(k+1)$, say $h(\varepsilon)$.

Let us show that for either $\alpha=\alpha_{1}$ or $\alpha=\alpha_{2}$, we must have that $a_{1, \alpha} \neq 0$. If $a_{1, \alpha}$ is zero in both cases, it implies that

$$
0=-\left(2 \operatorname{Tr}(A q(Y))+a_{1} \operatorname{Tr}(q(Y))\right) \operatorname{Tr}(q(Y))+2 a_{1} \operatorname{Tr}\left(q^{2}(Y)\right)
$$

This shows that the discriminant of $f(\alpha)$ (say $D$ ) is zero, i.e. $D=0$. Since we can choose $\varepsilon$ such that $g(\varepsilon) h(\varepsilon) \neq 0$, we then get that

$$
\begin{gathered}
\left(2 \operatorname{Tr}(A q(Y))+a_{1} \operatorname{Tr}(q(Y))\right)=\frac{2 a_{1} \operatorname{Tr}\left(q^{2}(Y)\right)}{\operatorname{Tr}(q(Y))}, \\
0=D=\frac{4 a_{1}^{2} \operatorname{Tr}\left(q^{2}(Y)\right)\left(\operatorname{Tr}\left(q^{2}(Y)\right)-\operatorname{Tr}^{2}(q(Y))\right)}{\operatorname{Tr}^{2}(q(Y))}
\end{gathered}
$$

We have

$$
\begin{aligned}
& \left.\operatorname{Tr}\left(q^{2}(Y)\right)-\operatorname{Tr}^{2}(q(Y))\right)=-2 q\left(\mu_{1}\right) q\left(\mu_{2}\right) \\
& \quad=\left(\mu_{1}^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i} \mu_{1}^{2 i}\right)\left(\mu_{2}^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i} \mu_{2}^{2 i}\right) \\
& =\left(\mu_{1} \mu_{2}\right)^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i} \mu_{1}^{2 i} \mu_{2}^{2(k+1)}+\sum_{i=1}^{k} \alpha_{i} \mu_{2}^{2 i} \mu_{1}^{2(k+1)} \\
& \quad+\sum_{i=1}^{k} \alpha_{i} \mu_{1}^{2 i} \sum_{i=1}^{k} \alpha_{i} \mu_{2}^{2 i}
\end{aligned}
$$

From this it is not difficult to see that this is a polynomial in $\varepsilon$, which we will denote $x(\varepsilon)$. Taking $\varepsilon$ such that $g(\varepsilon) h(\varepsilon) x(\varepsilon) \neq 0$, we reach a contradiction. Thus, using the above actions, we can arrive at the point $\left(b_{1}, \varepsilon,-\frac{b_{1}^{2}}{2}, b_{4}, b_{5}\right)$ with $b_{1} \neq 0$; we then use the proof of Case 1 to achieve the desired result.

Finally, we can prove infinite transitivity:
Theorem 1. The action of $G$ on $\mathcal{C}_{2}$ is a n-transitive group action for all $n \in \mathbb{Z}_{+}$, and hence infinitely transitive.

Proof. Let $A \in \mathcal{C}_{2} \backslash C_{k}$ be an arbitrary point. To prove the theorem, it is sufficient to show that there is a $g \in S_{k}$ such that $g(A)=\left(0,0,0,-1,2^{k+1}\right)$. Using Lemmas 10 and 11, we may assume that $a_{1} a_{3} \neq 0$ and $a_{1}^{2} \pm 2 a_{3} \neq 0$. Then, since $\Phi_{p} \in S_{k}$ when $p(t)=\alpha t^{2}+\beta t^{3}$, we can let

$$
\begin{aligned}
\alpha= & \frac{2}{a_{3}\left(a_{1}^{2}+2 a_{3}\right)\left(a_{1}^{2}-2 a_{3}\right)^{4}}\left(3 a_{1}^{8} a_{2} a_{3}-4 a_{1}^{6} a_{2} a_{3}^{2}-16 a_{1}^{4} a_{2} a_{3}^{3}+16 a_{1}^{2} a_{2} a_{3}^{4}\right. \\
& +16 a_{2} a_{3}^{5}-a_{1}^{9} a_{4}-4 a_{1}^{7} a_{3} a_{4}+16 a_{1}^{5} a_{3}^{2} a_{4}+16 a_{1}^{3} a_{3}^{3} a_{4}-48 a_{1} a_{3}^{4} a_{4} \\
& \left.+\left(a_{1}^{3}+6 a_{1} a_{3}\right) \sqrt{\left(a_{1}^{2}-2 a_{3}\right)^{4}\left(a_{1}^{2}+2 a_{3}\right)^{2}\left(1+2^{k+1} a_{3}\right)}\right)
\end{aligned}
$$

and

$$
\begin{array}{r}
\beta=\frac{4}{a_{3}\left(a_{1}^{2}-2 a_{3}\right)^{4}}\left(2 a_{1}^{5} a_{2} a_{3}-8 a_{1}^{3} a_{2} a_{3}^{2}+8 a_{1} a_{2} a_{3}^{3}-a_{1}^{6} a_{4}+2 a_{1}^{4} a_{3} a_{4}\right. \\
\left.\quad+4 a_{1}^{2} a_{3}^{2} a_{4}-8 a_{3}^{3} a_{4}+\sqrt{\left(a_{1}^{2}-2 a_{3}\right)^{4}\left(a_{1}^{2}+2 a_{3}\right)^{2}\left(1+2^{k+1} a_{3}\right)}\right)
\end{array}
$$

to move $A$ to the point

$$
A^{\prime}=\Phi_{p}(A)=\left(a_{1}, 0, a_{3}, a_{4}^{\prime}, 2^{k+1}\right)=\left(a_{1}, 0, \frac{a_{4}^{\prime 2}-1}{2^{k+1}}, a_{4}^{\prime}, 2^{k+1}\right)
$$

Then, letting $q(t)=\sum_{i=2}^{m} \beta_{i} t^{i}$, we will show that $\Psi_{q} \in S_{k}$.
We first write $q(t)=\sum_{i=2}^{\left[\frac{m}{2}\right]} \beta_{2 i} t^{2 i}+\sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1} t^{2 i+1}$. Then by Lemma 8 and Lemma 9 we get that

$$
\left\{\begin{array}{l}
\Psi_{q}(0,0,0,1,0)=(0,0,0,1,0), \\
\Psi_{q}(0,0,0,-1,0)=(0,0,0,1,0), \\
\Psi_{q}\left(0,0,0,-1,2^{j}\right)=\left(\sum_{i=2}^{\left[\frac{m}{2}\right]} 2^{i(j-1)+1} \beta_{2 i}, 0, \sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1}^{2} 2^{2 i(j-1)+3 j-2},\right. \\
\left.\sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1} 2^{i(j-1)+2 j-1}, 2^{j}\right), \quad j=\overline{1, k}, j \neq 3, \\
\Psi_{q}\left(a_{1}, 0, \frac{a_{4}^{\prime 2}-1}{2^{k+1}}, a_{4}, 2^{k+1}\right)=\left(a_{1}+\sum_{i=2}^{\left[\frac{m}{2}\right]} 2^{k i+1} \beta_{2 i}, 0,\right. \\
\left.\frac{a_{4}^{\prime 2}-1}{2^{k+1}}+\sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1}^{2} 2^{2 k(i+1)+k+1}, a_{4}^{\prime}+\sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1} 2^{k(i+2)}, 2^{k+1}\right)
\end{array}\right.
$$

From these we obtain the following systems:

$$
\begin{gather*}
\left\{\begin{array}{l}
\sum_{i=2}^{\left[\frac{m}{2}\right]} 2^{i(j-1)} \beta_{2 i}=0 \\
\sum_{i=2}^{\left[\frac{m}{2}\right]} 2^{k i+1} \beta_{2 i}=-a_{1}
\end{array} \quad j=\overline{1, k}, j \neq 3\right.  \tag{12}\\
\left\{\begin{array}{l}
\sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1} 2^{i(j-1)+2 j-1}=-1 \\
\sum_{i=1}^{\left[\frac{m}{2}\right]-1} \beta_{2 i+1} 2^{k(i+2)}=-1-a_{4}^{\prime}
\end{array}\right. \tag{13}
\end{gather*}
$$

We need to show the systems (12) and (13) have solutions. For the systems (12) and (13) we have the following matrix:

$$
\left(\begin{array}{cccccc}
1 & 1 & 1 & 1 & \ldots & 1 \\
1 & 2 & 2^{2} & 2^{3} & \ldots & 2^{k-1} \\
1 & 2^{3} & 2^{6} & 2^{9} & \ldots & 2^{3(k-1)} \\
1 & 2^{4} & 2^{8} & 2^{12} & \ldots & 2^{4(k-1)} \\
\vdots & \vdots & \vdots & \vdots & \ddots & \vdots \\
1 & 2^{k} & 2^{2 k} & 2^{3 k} & \ldots & 2^{k(k-1)}
\end{array}\right)
$$

The determinant of this matrix is Vandermonde's determinant and is not zero. Thus we can choose coefficients so that $\Psi_{q} \in S_{k}$, and applying this element to $A^{\prime}$, we get that $\Psi_{q}\left(A^{\prime}\right)=\left(0,0,0,-1,2^{k+1}\right)$. Therefore, all elements in $\mathcal{C}_{2} \backslash C_{k}$ are in the same orbit as $\left(0,0,0,-1,2^{k+1}\right)$ under the action of $S_{k}$, and hence $G$ acts $n$-transitively on $\mathcal{C}_{2}$, as desired.

## Conclusion

While Berest-Eshmatov-Eshmatov's conjecture has been recently proved, our approach for the case $n=2$ brings more clarity of the action on $\mathcal{C}_{2}$, which could be useful for future studies.

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